# Massachusetts Institute of Technology 

## Notes for $\mathbf{1 8 . 7 2 1}$

# INTRODUCTION TO ALGEBRAIC GEOMETRY 

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## TABLE OF CONTENTS

## Chapter 1: PLANE CURVES

### 1.1 The Affine Plane

1.2 The Projective Plane
1.3 Plane Projective Curves
1.4 Tangent Lines
1.5 Transcendence Degree
1.6 The Dual Curve
1.7 Resultants
1.8 Nodes and Cusps
1.9 Hensel's Lemma
1.10 Bézout's Theorem
1.11 The Plücker Formulas

## Chapter 2: AFFINE ALGEBRAIC GEOMETRY

2.1 Rings and Modules
2.2 The Zariski Topology
2.3 Some Affine Varieties
2.4 The Nullstellensatz
2.5 The Spectrum
2.6 Localization
2.7 Morphisms of Affine Varieties
2.8 Finite Group Actions

## Chapter 3: PROJECTIVE ALGEBRAIC GEOMETRY

3.1 Projective Varieties
3.2 Homogeneous Ideals
3.3 Product Varieties
3.4 Morphisms and Isomorphisms
3.5 Affine Varieties
3.6 Lines in Projective Three-Space

Chapter 4: INTEGRAL MORPHISMS OF AFFINE VARIETIES
4.1 The Nakayama Lemma
4.2 Integral Extensions
4.3 Normalization
4.4 Geometry of Integral Morphisms
4.5 Dimension
4.6 Krull's Theorem
4.7 Chevalley's Finiteness Theorem
4.8 Double Planes
5.1 Modules, review
5.2 Valuations
5.3 Smooth Curves
5.4 Constructible sets
5.5 Closed Sets
5.6 Fibred Products
5.7 Projective Varieties are Proper
5.8 Fibre Dimension

Chapter 6: MODULES
6.1 The Structure Sheaf
$6.2 \mathcal{O}$-Modules
6.3 The Sheaf Property
6.4 Some $\mathcal{O}$-Modules
6.5 Direct Image
6.6 Support
6.7 Twisting
6.8 Proof of Theorem 6.3.2

Chapter 7: COHOMOLOGY
7.1 Cohomology of $\mathcal{O}$-Modules
7.2 Complexes
7.3 Characteristic Properties
7.4 Construction of Cohomology
7.5 Cohomology of the Twisting Modules
7.6 Cohomology of Hypersurfaces
7.7 Three Theorems about Cohomology
7.8 Bézout's Theorem

Chapter 8: THE RIEMANN-ROCH THEOREM FOR CURVES
8.1 Branched Coverings
8.2 Divisors
8.3 The Riemann-Roch Theorem I
8.4 The Birkhoff-Grothendieck Theorem
8.5 Differentials
8.6 Trace
8.7 The Riemann-Roch Theorem II
8.8 Using Riemann-Roch

## Chapter 1 PLANE CURVES

1.1 The Affine Plane

1.2 The Projective Plane
1.3 Plane Projective Curves
1.4 Tangent Lines
1.5 Transcendence Degree
1.6 The Dual Curve
1.7 Resultants and Discriminants
1.8 Nodes and Cusps
1.9 Hensel's Lemma
1.10 Bézout's Theorem
1.11 The Plücker Formulas

Plane curves were the first algebraic varieties to be studied, so we begin with them. They provide helpful examples, and we will see in Chapter ?? how they control higher dimensional varieties. Chapters 2 - 7 are about varieties of arbitrary dimension. We will come back to curves in Chapter 8 .

### 1.1 The Affine Plane

affineplane
affcurve
goober8

The $n$-dimensional affine space $\mathbb{A}^{n}$ is the space of $n$-tuples of complex numbers. The affine plane $\mathbb{A}^{2}$ is the two-dimensional affine space.

Let $f\left(x_{1}, x_{2}\right)$ be an irreducible polynomial in two variables with complex coefficients. The set of points of the affine plane at which $f$ vanishes, the locus of zeros of $f$, is called a plane affine curve. Let's denote this locus by $X$. Using vector notation $x=\left(x_{1}, x_{2}\right)$,

$$
\begin{equation*}
X=\{x \mid f(x)=0\} \tag{1.1.1}
\end{equation*}
$$

The degree of the curve $X$ is the degree of its irreducible defining polynomial $f$.
1.1.2.


The Cubic Curve $y^{2}=x^{3}+x$ (real locus)
1.1.3. Note. In contrast with comples polynomials in one variable, most polynomials in two or more variables are irreducible - they cannot be factored. This can be shown by a method called "counting constants". For instance, quadratic polynomials in $x_{1}, x_{2}$ depend on the six coefficients of the monomials of degree at most two. Linear polynomials $a x_{1}+b x_{2}+c$ depend on three coefficients, but the product of two linear polynomials depends on only five parameters, because a scalar factor can be moved from one of the linear polynomials to the other. So the quadratic polynomials cannot all be written as products of linear polynomials. This reasoning is fairly convincing. It can be justified formally in terms of dimension, which will be discussed in Chapter ??.

We will get an understanding of the geometry of a plane curve as we go along, and we mention just one point here. A plane curve is called a curve because it is defined by one equation in two variables. Its algebraic dimension is one. But because our scalars are complex numbers, it will be a surface, geometrically. This is analogous to the fact that the affine line $\mathbb{A}^{1}$ is the plane of complex numbers.

One can see that a plane curve $X$ is a surface by inspecting its projection to the affine $x_{1}$-line. One writes the defining polynomial as a polynomial in $x_{2}$, whose coefficients $c_{i}=c_{i}\left(x_{1}\right)$ are polynomials in $x_{1}$ :

$$
f\left(x_{1}, x_{2}\right)=c_{0} x_{2}^{d}+c_{1} x_{2}^{d-1}+\cdots+c_{d}
$$

Let's suppose that $d$ is positive, i.e., that $f$ isn't a polynomial in $x_{1}$ alone (in which case, since it is irreducible, it would be linear).

The fibre of a map $V \rightarrow U$ over a point $p$ of $U$ is the inverse image of $p$, the set of points of $V$ that map to $p$. The fibre of the projection $X \rightarrow \mathbb{A}^{1}$ over the point $x_{1}=a$ is the set of points $(a, b)$ for which $b$ is a root of the one-variable polynomial

$$
f\left(a, x_{2}\right)=\bar{c}_{0} x_{2}^{d}+\bar{c}_{1} x_{2}^{d-1}+\cdots+\bar{c}_{d}
$$

with $\bar{c}_{i}=c_{i}(a)$. There will be finitely many points in this fibre, and the fibre won't be empty unless $f\left(a, x_{2}\right)$ is a constant. So the curve $X$ covers most of the $x_{1}$-line, a complex plane, finitely often.

## (1.1.4) changing coordinates

We allow linear changes of variable and translations in the affine plane $\mathbb{A}^{2}$. When a point $x$ is written as the column vector $x=\left(x_{1}, x_{2}\right)^{t}$, the coordinates $x^{\prime}=\left(x_{1}^{\prime}, x_{2}^{\prime}\right)^{t}$ after such a change of variable will be related to $x$ by the formula

$$
\begin{equation*}
x=Q x^{\prime}+a \tag{1.1.5}
\end{equation*}
$$

where $Q$ is an invertible $2 \times 2$ matrix with complex coefficients and $a=\left(a_{1}, a_{2}\right)^{t}$ is a complex translation vector. This changes a polynomial equation $f(x)=0$, to $f\left(Q x^{\prime}+a\right)=0$. One may also multiply a polynomial $f$ by a nonzero complex scalar without changing the locus $\{f=0\}$. Using these operations, all lines, plane curves of degree 1 , become equivalent.

An affine conic is a plane affine curve of degree two. Every affine conic is equivalent to one of the loci

$$
\begin{equation*}
x_{1}^{2}-x_{2}^{2}=1 \quad \text { or } \quad x_{2}=x_{1}^{2} \tag{1.1.6}
\end{equation*}
$$

The proof of this is similar to the one used to classify real conics. The two loci might be called a complex 'hyperbola' and 'parabola', respectively. The complex 'ellipse' $x_{1}^{2}+x_{2}^{2}=1$ becomes the 'hyperbola' when one multiplies $x_{2}$ by $i$.

On the other hand, there are infinitely many inequivalent cubic curves. Cubic polynomials in two variables depend on the coefficients of the ten monomials $1, x_{1}, x_{2}, x_{1}^{2}, x_{1} x_{2}, x_{2}^{2}, x_{1}^{3}, x_{1}^{2} x_{2}, x_{1} x_{2}^{2}, x_{2}^{3}$ of degree at most 3 in $x$. Linear changes of variable, translations, and scalar multiplication give us only seven scalars to work with, leaving three essential parameters.

### 1.2 The Projective Plane

projplane equivrel projline projpl pline
eqline
linesmeet
standcov

The $n$-dimensional projective space $\mathbb{P}^{n}$ is the set of equivalence classes of nonzero vectors $x=\left(x_{0}, x_{1}, \ldots, x_{n}\right)$, the equivalence relation being

$$
\begin{equation*}
\left(x_{0}^{\prime}, \ldots, x_{n}^{\prime}\right) \sim\left(x_{0}, \ldots, x_{n}\right) \quad \text { if } \quad\left(x_{0}^{\prime}, \ldots, x_{n}^{\prime}\right)=\left(\lambda x_{0}, \ldots, \lambda x_{n}\right) \tag{1.2.1}
\end{equation*}
$$

or $x^{\prime}=\lambda x$, for some nonzero complex number $\lambda$. The equivalence classes are the points of $\mathbb{P}^{n}$, and one often refers to a point by a particular vector in its class.

Points of $\mathbb{P}^{n}$ correspond bijectively to one-dimensional subspaces of $\mathbb{C}^{n+1}$. When $x$ is a nonzero vector, the vectors $\lambda x$, together with the zero vector, form the one-dimensional subspace of the complex vector space $\mathbb{C}^{n+1}$ spanned by $x$.

The projective plane $\mathbb{P}^{2}$ is the two-dimensional projective space. Its points are equivalence classes of nonzero vectors $\left(x_{0}, x_{1}, x_{2}\right)$.

## (1.2.2) the projective line

Points of the projective line $\mathbb{P}^{1}$ are equivalence classes of nonzero vectors $\left(x_{0}, x_{1}\right)$. If $x_{0}$ isn't zero, we may multiply by $\lambda=x_{0}^{-1}$ to normalize the first entry of $\left(x_{0}, x_{1}\right)$ to 1 , and write the point it represents in a unique way as $(1, u)$, with $u=x_{1} / x_{0}$. There is one remaining point, the point represented by the vector $(0,1)$. The projective line $\mathbb{P}^{1}$ can be obtained by adding this point, called the point at infinity, to the affine $u$-line, which is a complex plane. Topologically, $\mathbb{P}^{1}$ is a two-dimensional sphere.

## (1.2.3) lines in projective space

A line in projective space $\mathbb{P}^{n}$ is determined by a pair of distinct points $p$ and $q$. When $p$ and $q$ are represented by specific vectors, the set of points $\{r p+s q\}$, with $r, s$ in $\mathbb{C}$ not both zero is a line. Its points correspond bijectively to points of the projective line $\mathbb{P}^{1}$, by

$$
\begin{equation*}
r p+s q \quad \longleftrightarrow \quad(r, s) \tag{1.2.4}
\end{equation*}
$$

A line in the projective plane $\mathbb{P}^{2}$ can also be described as the locus of solutions of a homogeneous linear equation

$$
\begin{equation*}
s_{0} x_{0}+s_{1} x_{1}+s_{2} x_{2}=0 \tag{1.2.5}
\end{equation*}
$$

1.2.6. Lemma. In the projective plane, two distinct lines have exactly one point in common and, in a projective space of any dimension, a pair of distinct points is contained in exactly one line.

## (1.2.7) $\quad$ the standard covering of $\mathbb{P}^{2}$

If the first entry $x_{0}$ of a point $p=\left(x_{0}, x_{1}, x_{2}\right)$ of the projective plane $\mathbb{P}^{2}$ isn't zero, we may normalize it to 1 without changing the point: $\left(x_{0}, x_{1}, x_{2}\right) \sim\left(1, u_{1}, u_{2}\right)$, where $u_{i}=x_{i} / x_{0}$. We did the analogous thing for $\mathbb{P}^{1}$ above. The representative vector $\left(1, u_{1}, u_{2}\right)$ is uniquely determined by $p$, so points with $x_{0} \neq 0$ correspond bijectively to points of an affine plane $\mathbb{A}^{2}$ with coordinates $\left(u_{1}, u_{2}\right)$ :

$$
\left(x_{0}, x_{1}, x_{2}\right) \sim\left(1, u_{1}, u_{2}\right) \quad \longleftrightarrow \quad\left(u_{1}, u_{2}\right)
$$

We regard the affine plane as a subset of $\mathbb{P}^{2}$ by this correspondence, and we denote that subset by $\mathbb{U}^{0}$. The points of $\mathbb{U}^{0}$, those with $x_{0} \neq 0$, are the points at finite distance. The points at infinity of $\mathbb{P}^{2}$, those of the form $\left(0, x_{1}, x_{2}\right)$, are on the line at infinity $L^{0}$, the locus $\left\{x_{0}=0\right\}$. The projective plane is the union of the two sets $\mathbb{U}^{0}$ and $L^{0}$. When a point is given by a coordinate vector, we can assume that the first coordinate is either 1 or 0.

There is an analogous correspondence between points $\left(x_{0}, 1, x_{2}\right)$ and points of an affine plane $\mathbb{A}^{2}$, and between points $\left(x_{0}, x_{1}, 1\right)$ and points of $\mathbb{A}^{2}$. We denote the subsets $\left\{x_{1} \neq 0\right\}$ and $\left\{x_{2} \neq 0\right\}$ by $\mathbb{U}^{1}$ and $\mathbb{U}^{2}$, respectively. The three sets $\mathbb{U}^{0}, \mathbb{U}^{1}, \mathbb{U}^{2}$ form the standard covering of $\mathbb{P}^{2}$ by three standard affine open sets. Since the vector $(0,0,0)$ has been ruled out, every point of $\mathbb{P}^{2}$ lies in at least one of the standard affine open sets. Points whose three coordinates are nonzero lie in all of them.
1.2.8. Note. Which points of $\mathbb{P}^{2}$ are at infinity depends on which of the standard affine open sets is taken to be the one at finite distance. When the coordinates are ( $x_{0}, x_{1}, x_{2}$ ), I like to normalize $x_{0}$ to 1 , as above. Then the points at infinity are those of the form $\left(0, x_{1}, x_{2}\right)$. But when coordinates are $(x, y, z)$, I may normalize $z$ to 1 . Then the points at infinity are the points $(x, y, 0)$. I hope this won't cause too much confusion.

## (1.2.9) digression: the real projective plane

The points of the real projective plane $\mathbb{R} \mathbb{P}^{2}$ are equivalence classes of nonzero real vectors $x=\left(x_{0}, x_{1}, x_{2}\right)$, the equivalence relation being $x^{\prime} \sim x$ if $x^{\prime}=\lambda x$ for some nonzero real number $\lambda$. The real projective plane can also be thought of as the set of one-dimensional subspaces of the real vector space $V=\mathbb{R}^{3}$.

The plane $U:\left\{x_{0}=1\right\}$ in $V=\mathbb{R}^{3}$ is analogous to the standard affine open subset $\mathbb{U}^{0}$ in the complex projective plane $\mathbb{P}^{2}$. We can project $V$ from the origin $p_{0}=(0,0,0)$ to $U$, sending a point $x=\left(x_{0}, x_{1}, x_{2}\right)$ of $V$ distinct from $p_{0}$ to the point $\left(1, u_{1}, u_{2}\right)$, with $u_{i}=x_{i} / x_{0}$. The fibres of this projection are the lines through $p_{0}$ and $x$, with $p_{0}$ omitted. The projection to $U$ is undefined at the points $\left(0, x_{1}, x_{2}\right)$, which are orthogonal to the $x_{0}$-axis. The line connecting such a point to $p_{0}$ doesn't meet $U$. The points $\left(0, x_{1}, x_{2}\right)$ correspond to the points at infinity of $\mathbb{R P}^{2}$.

Looking from the origin, $U$ becomes a "picture plane".

### 1.2.10.

pointatinfinity
realprojplane


This illustration is from Dürer's book on perspective

The projection from 3 -space to a picture plane goes back to the the 16th century, the time of Desargues and Dürer. Projective coordinates were introduced by Möbius, but not until 200 years later.
1.2.11.


## A Schematic Representation of the real Projective Plane, with a Conic

This figure shows the plane $W: x+y+z=1$ in the real vector space $\mathbb{R}^{3}$. If $p=(x, y, z)$ is a nonzero vector, the one-dimensional subspace spanned by $p$ will meet $W$ in a single point, unless $p$ is on the line $L: x+y+z=0$. The plane $W$ is a faithful representation of most of $\mathbb{R} \mathbb{P}^{2}$. It contains all points except those on the line $L$.

## (1.2.12) changing coordinates in the projective plane

An invertible $3 \times 3$ matrix $P$ determines a linear change of coordinates in $\mathbb{P}^{2}$. With $x=\left(x_{0}, x_{1}, x_{2}\right)^{t}$ and $x^{\prime}=\left(x_{0}^{\prime}, x_{1}^{\prime}, x_{2}^{\prime}\right)^{t}$ represented as column vectors, the coordinate change is given by

$$
\begin{equation*}
P x^{\prime}=x \tag{1.2.13}
\end{equation*}
$$

As the next proposition shows, four special points, the three points $e_{0}=(1,0,0)^{t}, e_{1}=(0,1,0)^{t}, e_{2}=$ $(0,0,1)^{t}$, together with the point $\epsilon=(1,1,1)^{t}$, determine the coordinates.
1.2.14. Proposition. Let $p_{0}, p_{1}, p_{2}, q$ be four points of $\mathbb{P}^{2}$, no three of which lie on a line. There is, up to scalar factor, a unique linear coordinate change $P x^{\prime}=x$ such that $P p_{i}=e_{i}$ and $P q=\epsilon$.
proof. The hypothesis that the points $p_{0}, p_{1}, p_{2}$ don't lie on a line means that the vectors that represent those points are independent. They span $\mathbb{C}^{3}$. So $q$ will be a combination $q=c_{0} p_{0}+c_{1} p_{1}+c_{2} p_{2}$, and because no three of the points lie on a line, the coefficients $c_{i}$ will be nonzero. We can scale the vectors $p_{i}$ (multiply them by nonzero scalars) to make $q=p_{0}+p_{1}+p_{2}$ without changing the points. Next, the columns of $P$ can be an arbitrary set of independent vectors. We let them be $p_{0}, p_{1}, p_{2}$. Then $P e_{i}=p_{i}$, and $P \epsilon=q$. The matrix $P$ is unique up to scalar factor, as can be verified by looking the reasoning over.

A polynomial $f\left(x_{0}, x_{1}, x_{2}\right)$ is homogeneous, and of degree $d$, if all monomials that appear with nonzero coefficient have (total) degree $d$. For example, $x_{0}^{3}+x_{1}^{3}-x_{0} x_{1} x_{2}$ is a homogeneous cubic polynomial. A homogeneous quadratic polynomal is a combination of the six monomials

$$
x_{0}^{2}, x_{1}^{2}, x_{2}^{2}, x_{0} x_{1}, x_{1} x_{2}, x_{0} x_{2}
$$

A conic is the locus of zeros of an irreducible homogeneous quadratic polynomial.
1.2.16. Proposition. For any conic $C$, there is a choice of coordinates so that $C$ becomes the locus

$$
x_{0} x_{1}+x_{0} x_{2}+x_{1} x_{2}=0
$$

proof. Since the conic $C$ isn't a line, it will contain three points that aren't colinear. Let's leave the verification of this fact as an exercise. We choose three non-colinear points on $C$, and adjust coordinates so that they become the points $e_{0}, e_{1}, e_{2}$. Let $f$ be the quadratic polynomial in those coordinates whose zero locus is $C$. Because $e_{0}$ is a point of $C, f(1,0,0)=0$, and therefore the coefficient of $x_{0}^{2}$ in $f$ is zero. Similarly, the coefficients of $x_{1}^{2}$ and $x_{2}^{2}$ are zero. So $f$ has the form

$$
f=a x_{0} x_{1}+b x_{0} x_{2}+c x_{1} x_{2}
$$

Since $f$ is irreducible, $a, b, c$ aren't zero. By scaling appropriately, we can make $a=b=c=1$. We will be left with the polynomial $x_{0} x_{1}+x_{0} x_{2}+x_{1} x_{2}$.

### 1.3 Plane Projective Curves

The loci in projective space that are studied in algebraic geometry are those that can be defined by systems of homogeneous polynomial equations. The reason for homogeneity is that the vectors $\left(a_{0}, \ldots, a_{n}\right)$ and $\left(\lambda a_{0}, \ldots, \lambda a_{n}\right)$ represent the same point of $\mathbb{P}^{n}$.

To explain this, we write a polynomial $f\left(x_{0}, \ldots, x_{n}\right)$ as a sum of its homogeneous parts:

$$
\begin{equation*}
f=f_{0}+f_{1}+\cdots+f_{d} \tag{1.3.1}
\end{equation*}
$$

homparts
where $f_{0}$ is the constant term, $f_{1}$ is the linear part, etc., and $d$ is the degree of $f$.
1.3.2. Lemma. Let $f$ be a polynomial of degree $d$, and let $a=\left(a_{0}, \ldots, a_{n}\right)$ be a nonzero vector. Then $f(\lambda a)=0$ for every nonzero complex number $\lambda$ if and only if $f_{i}(a)$ is zero for every $i=0, \ldots, d$.
proof. We have $f\left(\lambda x_{0}, \ldots, \lambda x_{n}\right)=f_{0}+\lambda f_{1}(x)+\lambda^{2} f_{2}(x)+\cdot+\lambda^{d} f_{d}(x)$. When we evaluate at some given $x$, the right side of this equation becomes a polynomial of degree at most $d$ in $\lambda$. Since a nonzero polynomial of degree at most $d$ has at most $d$ roots, $f(\lambda x)$ won't be zero for every $\lambda$ unless that polynomial is zero.

Thus we may as well work with homogeneous equations.
1.3.3. Lemma. If a homogeneous polynomial $f$ is a product $g h$ of polynomials, then $g$ and $h$ are homogeneous, and the zero locus $\{f=0\}$ in projective space is the union of the two loci $\{g=0\}$ and $\{h=0\}$.

It is also true that relatively prime homogeneous polynomials $f$ and $g$ have only finitely many common zeros. But this isn't obvious. It will be proved below, in Proposition 1.3.11.

## (1.3.4) loci in the projective line

Before going to plane curves, we describe the zero locus in the projective line $\mathbb{P}^{1}$ of a homogeneous polynomial in two variables.
1.3.5. Lemma. Every nonzero homogeneous polynomial $f(x, y)=a_{0} x^{d}+a_{1} x^{d-1} y+\cdots+a_{d} y^{d}$ with complex coefficients is a product of homogeneous linear polynomials that are unique up to scalar factor.

To prove this, one uses the fact that the field of complex numbers is algebraically closed. A one-variable complex polynomial factors into linear factors in the polynomial ring $\mathbb{C}[y]$. To factor $f(x, y)$, one may factor the one-variable polynomial $f(1, y)$ into linear factors, substitute $y / x$ for $y$, and multiply the result by $x^{d}$. When one adjusts scalar factors, one will obtain the expected factorization of $f(x, y)$. For instance, to factor $f(x, y)=x^{2}-3 x y+2 y^{2}$, substitute $x=1: 2 y^{2}-3 y+1=2(y-1)\left(y-\frac{1}{2}\right)$. Substituting $y=y / x$ and multiplying by $x^{2}, f(x, y)=2(y-x)\left(y-\frac{1}{2} x\right)$. The scalar 2 can be distributed arbitrarily among the linear factors.

Adjusting scalar factors, we may write a homogeneous polynomial as a product of the form

$$
\begin{equation*}
f(x, y)=\left(v_{1} x-u_{1} y\right)^{r_{1}} \cdots\left(v_{k} x-u_{k} y\right)^{r_{k}} \tag{1.3.6}
\end{equation*}
$$

where no factor $v_{i} x-u_{i} y$ is a constant multiple of another, and where $r_{1}+\cdots+r_{k}$ is the degree of $f$. The exponent $r_{i}$ is the multiplicity of the linear factor $v_{i} x-u_{i} y$.

A linear polynomial $v x-u y$ determines a point $(u, v)$ in the projective line $\mathbb{P}^{1}$, the unique zero of that polynomial, and changing the polynomial by a scalar factor doesn't change its zero. Thus the linear factors of the homogeneous polynomial 1.3 .6 determine points of $\mathbb{P}^{1}$, the zeros of $f$. The points $\left(u_{i}, v_{i}\right)$ are zeros of multiplicity $r_{i}$. The total number of those points, counted with multiplicity, will be the degree of $f$.

The zero ( $u_{i}, v_{i}$ ) of $f$ corresponds to a root $x=u_{i} / v_{i}$ of multiplicity $r_{i}$ of the one-variable polynomial $f(x, 1)$, except when the zero is the point $(1,0)$. This happens when the coefficient $a_{0}$ of $f$ is zero, and $y$ is a factor of $f$. One could say that $f(x, y)$ has a zero at infinity in that case.

This sums up the information contained in an algebraic locus in the projective line. It will be a finite set of points with multiplicities.

## (1.3.7) intersections with a line

Let $Z$ be the zero locus of a homogeneous polynomial $f\left(x_{0}, \ldots, x_{n}\right)$ of degree $d$ in projective space $\mathbb{P}^{n}$, and let $L$ be a line in $\mathbb{P}^{n} 1.2 .4$. Say that $L$ is the set of points $r p+s q$, where $p=\left(a_{0}, \ldots, a_{n}\right)$ and $q=\left(b_{0}, \ldots, b_{n}\right)$ are represented by specific vectors, so that $L$ corresponds to the projective line $\mathbb{P}^{1}$ by $r p+s q \leftrightarrow(r, s)$. Let's also assume that $L$ isn't entirely contained in the zero locus $Z$. The intersection $Z \cap L$ corresponds in $\mathbb{P}^{1}$ to the zero locus of the polynomial in $r, s$ that is obtained by substituting $r p+s q$ into $f$. This substitution yields a homogeneous polynomial $\bar{f}(r, s)$ in $r, s$, of degree $d$. For example, if $f=x_{0} x_{1}+x_{0} x_{2}+x_{1} x_{2}$, then with $p=\left(a_{o}, a_{1}, a_{2}\right)$ and $q=\left(b_{0}, b_{1}, b_{2}\right), \bar{f}$ is the following quadratic polynomial in $r, s$ :

$$
\begin{aligned}
\bar{f}(r, s)=f(r p+s q) & =\left(r a_{0}+s b_{0}\right)\left(r a_{1}+s b_{1}\right)+\left(r a_{0}+s b_{0}\right)\left(r a_{2}+s b_{2}\right)+\left(r a_{1}+s b_{1}\right)\left(r a_{2}+s b_{2}\right) \\
& =\left(a_{0} a_{1}+a_{0} a_{2}+a_{1} a_{2}\right) r^{2}+\left(\sum_{i \neq j} a_{i} b_{j}\right) r s+\left(b_{0} b_{1}+b_{0} b_{2}+b_{1} b_{2}\right) s^{2}
\end{aligned}
$$

The zeros of $\bar{f}$ in $\mathbb{P}^{1}$ correspond to the points of $Z \cap L$. There will be $d$ zeros, when counted with multiplicity.
1.3.8. Definition. With notation as above, the intersection multiplicity of $Z$ and $L$ at a point $p$ is the multiplicity of zero of the polynomial $\bar{f}$.
1.3.9. Corollary. Let $Z$ be the zero locus in $\mathbb{P}^{n}$ of a homogeneous polynomial $f$, and let $L$ be a line in $\mathbb{P}^{n}$ that isn't contained in $Z$. The number of intersections of $Z$ and $L$, counted with multiplicity, is equal to the degree of $f$.

## (1.3.10) loci in the projective plane

1.3.11. Proposition. Homogeneous polynomials $f_{1}, \ldots, f_{r}$ in $x, y, z$ with no common factor have finitely many common zeros in $\mathbb{P}^{2}$.

The proof of this proposition is below. It shows that the most interesting type of locus in the projective plane is the zero set of a single equation.

The locus of zeros of an irreducible homogeneous polynomial $f$ is called a plane projective curve. The degree of a plane projective curve is the degree of its irreducible defining polynomial.
1.3.12. Note. Suppose that a homogeneous polynomial is reducible, say $f=g_{1} \cdots g_{k}$, where $g_{i}$ are irreducible, and such that $g_{i}$ and $g_{j}$ don't differ by a scalar factor when $i \neq j$. Then the zero locus $C$ of $f$ is the union of the zero loci $V_{i}$ of the factors $g_{i}$. In this case, $C$ may be called a reducible curve.

When there are multiple factors, say $f=g_{1}^{e_{1}} \cdots g_{k}^{e_{k}}$ and some $e_{i}$ are greater than 1 , it is still true that the locus $C:\{f=0\}$ will be the union of the loci $V_{i}:\left\{g_{i}=0\right\}$, but the connection between the geometry of $C$ and the algebra is weakened. In this situation, the structure of a scheme becomes useful. We won't discuss schemes. The only situation in which we will need to keep track of multiple factors is when counting intersections with another curve $D$. For this purpose, one can define the divisor of $f$ to be the integer combination $e_{1} V_{1}+\cdots+e_{k} V_{k}$.

We need a lemma for the proof of Proposition 1.3.11. The ring $\mathbb{C}[x, y]$ embeds into its field of fractions $F$, which is the field of rational functions $\mathbb{C}(x, y)$ in $x, y$, and the polynomial ring $\mathbb{C}[x, y, z]$ is a subring of the one-variable polynomial ring $F[z]$. It can be useful to study a problem in the principal ideal domain $F[z]$ first because its algebra is simpler.

Recall that the unit ideal of a ring $R$ is the ring $R$ itself.
1.3.13. Lemma. Let $F=\mathbb{C}(x, y)$ be the field of rational functions in $x, y$.
(i) Let $f_{1}, \ldots, f_{k}$ be homogeneous polynomials in $x, y, z$ with no common factor. Their greatest common divisor in $F[z]$ is 1 , and therefore $f_{1}, \ldots, f_{k}$ generate the unit ideal of $F[z]$. There is an equation of the form $\sum g_{i}^{\prime} f_{i}=$ 1 with $g_{i}^{\prime}$ in $F[z]$.
(ii) Let $f$ be an irreducible polynomial in $\mathbb{C}[x, y, z]$ of positive degree in $z$, but not divisible by $z$. Then $f$ is also an irreducible element of $F[z]$.
proof. (i) Let $h^{\prime}$ be an element of $F[z]$ that isn't a unit of $F[z]$, i.e., that isn't an element of $F$. Suppose that, for every $i, h^{\prime}$ divides $f_{i}$ in $F[z]$, say $f_{i}=u_{i}^{\prime} h^{\prime}$. The coefficients of $h^{\prime}$ and $u_{i}^{\prime}$ have denominators that are polynomials in $x, y$. We clear denominators from the coefficients, to obtain elements of $\mathbb{C}[x, y, z]$. This will give us equations of the form $d_{i} f_{i}=u_{i} h$, where $d_{i}$ are polynomials in $x, y$ and $u_{i}, h$ are polynomials in $x, y, z$.

Since $h$ isn't in $F$, it will have positive degree in $z$. Let $g$ be an irreducible factor of $h$ of positive degree in $z$. Then $g$ divides $d_{i} f_{i}$ but doesn't divide $d_{i}$ which has degree zero in $z$. So $g$ divides $f_{i}$, and this is true for every $i$. This contradicts the hypothesis that $f_{1}, \ldots, f_{k}$ have no common factor.
(ii) Say that $f(x, y, z)$ factors in $F[z], f=g^{\prime} h^{\prime}$, where $g^{\prime}$ and $h^{\prime}$ are polynomials of positive degree in $z$ with coefficients in $F$. When we clear denominators from $g^{\prime}$ and $h^{\prime}$, we obtain an equation of the form $d f=g h$, where $g$ and $h$ are polynomials in $x, y, z$ of positive degree in $z$ and $d$ is a polynomial in $x, y$. Neither $g$ nor $h$ divides $d$, so $f$ must be reducible.
proof of Proposition 1.3.11 We are to show that homogeneous polynomials $f_{1}, \ldots, f_{r}$ in $x, y, z$ with no common factor have finitely many common zeros. Lemma 1.3 .13 tells us that we may write $\sum g_{i}^{\prime} f_{i}=1$, with $g_{i}^{\prime}$ in $F[z]$. Clearing denominators from $g_{i}^{\prime}$ gives us an equation of the form

$$
\sum g_{i} f_{i}=d
$$

where $g_{i}$ are polynomials in $x, y, z$ and $d$ is a polynomial in $x, y$. Taking suitable homogeneous parts of $g_{i}$ and $d$ produces an equation $\sum g_{i} f_{i}=d$ in which all terms are homogeneous.

Lemma 1.3 .5 asserts that $d$ is a product of linear polynomials, say $d=\ell_{1} \cdots \ell_{r}$. A common zero of $f_{1}, \ldots, f_{k}$ is also a zero of $d$, and therefore it is a zero of $\ell_{j}$ for some $j$. It suffices to show that, for every $j$, $f_{1}, \ldots, f_{r}$ and $\ell_{j}$ have finitely many common zeros.

Since $f_{1}, \ldots, f_{k}$ have no common factor, there is at least one $f_{i}$ that isn't divisible by $\ell_{j}$. Corollary 1.3 .9 shows that $f_{i}$ and $\ell_{j}$ have finitely many common zeros. Therefore $f_{1}, \ldots, f_{k}$ and $\ell_{j}$ have finitely many common zeros for every $j$.
1.3.14. Corollary. Every locus in the projective plane $\mathbb{P}^{2}$ that can be defined by a system of homogeneous polynomial equations is a finite union of points and curves.

The next corollary is a special case of the Strong Nullstellensatz, which will be proved in the next chapter.
1.3.15. Corollary. Let $f$ be an irreducible homogeneous polynomial in three variables, that vanishes on an infinite set $S$ of points of $\mathbb{P}^{2}$. If another homogeneous polynomial $g$ vanishes on $S$, then $f$ divides $g$. Therefore, if an irreducible polynomial vanishes on an infinite set $S$, that polynomial is unique up to scalar factor.
proof. If the irreducible polynomial $f$ doesn't divide $g$, then $f$ and $g$ have no common factor, and therefore they have finitely many common zeros.

## (1.3.16) the classical topology

The usual topology on the affine space $\mathbb{A}^{n}$ will be called the classical topology. A subset $U$ of $\mathbb{A}^{n}$ is open in the classical topology if, whenever $U$ contains a point $p$, it contains all points sufficiently near to $p$. We call this
relprime
the classical topology to distinguish it from another topology, the Zariski topology, which will be discussed in the next chapter.

The projective space $\mathbb{P}^{n}$ also has a classical topology. A subset $U$ of $\mathbb{P}^{n}$ is open if, whenever a point $p$ of $U$ is represented by a vector $\left(x_{0}, \ldots, x_{n}\right)$, all vectors $x^{\prime}=\left(x_{0}^{\prime}, \ldots, x_{n}^{\prime}\right)$ sufficiently near to $x$ represent points of $U$.
isopts
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## (1.3.17) isolated points

A point $p$ of a topological space $X$ is isolated if both $\{p\}$ and its complement $X-\{p\}$ are closed sets, or if $\{p\}$ is both open and closed. If $X$ is a subset of $\mathbb{A}^{n}$ or $\mathbb{P}^{n}$, a point $p$ of $X$ is isolated in the classical topology if $X$ doesn't contain points $p^{\prime}$ distinct from $p$, but arbitrarily close to $p$.
1.3.18. Proposition Let $n$ be an integer greater than one. The zero locus of a polynomial in $\mathbb{A}^{n}$ or in $\mathbb{P}^{n}$ contains no points that are isolated in the classical topology.
1.3.19. Lemma. Let $f$ be a polynomial of degree $d$ in $n$ variables. When coordinates $x_{1}, \ldots, x_{n}$ are chosen suitably, $f(x)$ will ge a monic polynomial of degree $d$ in the variable $x_{n}$.
proof. We write $f=f_{0}+f_{1}+\cdots+f_{d}$, where $f_{i}$ is the homogeneous part of $f$ of degree $i$, and we choose a point $p$ of $\mathbb{A}^{n}$ at which $f_{d}$ isn't zero. We change variables so that $p$ becomes the point $(0, \ldots, 0,1)$. We call the new variables $x_{1}, \ldots, . x_{n}$ and the new polynomial $f$. Then $f_{d}\left(0, \ldots, 0, x_{n}\right)$ will be equal to $c x_{n}^{d}$ for some nonzero constant $c$. When we adjust $x_{n}$ by a scalar factor to make $c=1$, $f$ will be monic.
proof of Proposition 1.3.18. The proposition is true for loci in affine space and also for loci in projective space. We look at the affine case. Let $f\left(x_{1}, \ldots, x_{n}\right)$ be a polynomial with zero locus $Z$, and let $p$ be a point of $Z$. We adjust coordinates so that $p$ is the origin $(0, \ldots, 0)$ and $f$ is monic in $x_{n}$. We relabel $x_{n}$ as $y$, and write $f$ as a polynomial in $y$. Let's write $f(x, y)=\widetilde{f}(y)$ :

$$
\tilde{f}(y)=f(x, y)=y^{d}+c_{d-1}(x) y^{d-1}+\cdots+c_{0}(x)
$$

where $c_{i}$ is a polynomial in $x_{1}, \ldots, x_{n-1}$. For fixed $x, c_{0}(x)$ is the product of the roots of $\tilde{f}(y)$. Since $p$ is the origin and $f(p)=0, \quad c_{0}(0)=0$. So $c_{0}(x)$ will tend to zero with $x$. Then at least one root $y$ of $\widetilde{f}(y)$ will tend to zero. This gives us points $(x, y)$ of $Z$ that are arbitrarily close to $p$.
1.3.20. Corollary. Let $C^{\prime}$ be the complement of a finite set of points in a plane curve C. In the classical topology, a continuous function $g$ on $C$ that is zero at every point of $C^{\prime}$ is identically zero.

### 1.4 Tangent Lines

## (1.4.1) notation for working locally

We will often want to inspect a plane curve $C:\left\{f\left(x_{0}, x_{1}, x_{2}\right)=0\right\}$ in a neighborhood of a particular point $p$. To do this we may adjust coordinates so that $p$ becomes the point $(1,0,0)$, and look in the standard affine open set $\mathbb{U}^{0}:\left\{x_{0} \neq 0\right\}$. There, $p$ becomes the origin in the affine $x_{1}, x_{2}$-plane, and $C$ becomes the zero locus of the non-homogeneous polynomial $f\left(1, x_{1}, x_{2}\right)$. This will be a standard notation for working locally.

Of course, it doesn't matter which variable we set to 1 . If the variables are $x, y, z$, we may prefer to take for $p$ the point $(0,0,1)$ and work with the polynomial $f(x, y, 1)$.
1.4.2. Lemma. A homogeneous polynomial $f\left(x_{0}, x_{1}, x_{2}\right)$ not divisible by $x_{0}$ is irreducible if and only if its dehomogenization $f\left(1, x_{1}, x_{2}\right)$ is irreducible.

## (1.4.3) homogenizing and dehomogenizing

If $f\left(x_{0}, x_{1}, x_{2}\right)$ is a polynomial, $f\left(1, x_{1}, x_{2}\right)$ is called the dehomogenization of $f$ with respect to the variable $x_{0}$.

A simple procedure, homogenization, inverts dehomogenization. Suppose given a non-homogeneous polynomial $F\left(x_{1}, x_{2}\right)$ of degree $d$. To homogenize $F$, we replace the variables $x_{i}, \quad i=1,2$, by $u_{i}=x_{i} / x_{0}$. Then since $u_{i}$ have degree zero in $x$, so does $F\left(u_{1}, u_{2}\right)$. When we multiply by $x_{0}^{d}$, the result will be a homogeneous polynomial of degree $d$ in $x_{0}, x_{1}, x_{2}$ that isn't divisible by $x_{0}$,

We will come back to homogenization in Chapter 3

## (1.4.4) smooth points and singular points

Let $C$ be the plane curve defined by an irreducible homogeneous polynomial $f\left(x_{0}, x_{1}, x_{2}\right)$, and let $f_{i}$ denote the partial derivative $\frac{\partial f}{\partial x_{i}}$, which can be computed by the usual calculus formula. A point of $C$ at which the partial derivatives $f_{i}$ aren't all zero is called a smooth point of $C$. A point at which all partial derivatives are zero is a singular point. A curve is smooth, or nonsingular, if it contains no singular point; otherwise it is a singular curve.

The Fermat curve

$$
\begin{equation*}
x_{0}^{d}+x_{1}^{d}+x_{2}^{d}=0 \tag{1.4.5}
\end{equation*}
$$

is smooth because the only common zero of the partial derivatives $d x_{0}^{d-1}, d x_{1}^{d-1}, d x_{2}^{d-1}$, which is $(0,0,0)$, doesn't represent a point of $\mathbb{P}^{2}$. The cubic curve $x_{0}^{3}+x_{1}^{3}-x_{0} x_{1} x_{2}=0$ is singular at the point $(0,0,1)$.

The Implicit Function Theorem explains the meaning of smoothness. Suppose that $p=(1,0,0)$ is a point of $C$. We set $x_{0}=1$ and inspect the locus $f\left(1, x_{1}, x_{2}\right)=0$ in the standard affine open set $\mathbb{U}^{0}$. If $f_{2}(p)$ isn't zero, the Implicit Function Theorem tells us that we can solve the equation $f\left(1, x_{1}, x_{2}\right)=0$ for $x_{2}$ locally (i.e., for small $x_{1}$ ) as an analytic function $\varphi$ of $x_{1}$, with $\varphi(0)=0$. Sending $x_{1}$ to $\left(1, x_{1}, \varphi\left(x_{1}\right)\right)$ inverts the projection from $C$ to the affine $x_{1}$-line, locally. So at a smooth point, $C$ is locally homeomorphic to the affine line.
1.4.6. Euler's Formula. Let $f$ be a homogeneous polynomial of degree $d$ in the variables $x_{0}, \ldots, x_{n}$. Then

$$
\sum_{i} x_{i} \frac{\partial f}{\partial x_{i}}=d f
$$

proof. It is enough to check this formula when $f$ is a monomial. As an example, let $f$ be the monomial $x^{2} y^{3} z$, then

$$
x f_{x}+y f_{y}+z f_{z}=x\left(2 x y^{3} z\right)+y\left(3 x^{2} y^{2} z\right)+z\left(x^{2} y^{3}\right)=6 x^{2} y^{3} z=6 f
$$

1.4.7. Corollary. (i) If all partial derivatives of an irreducible homogeneous polynomial $f$ are zero at a point $p$ of $\mathbb{P}^{2}$, then $f$ is zero at $p$, and therefore $p$ is a singular point of the curve $\{f=0\}$.
(ii) At a smooth point of a plane curve, at least two partial derivatives will be nonzero.
(iii) The partial derivatives of an irreducible polynomial have no common (nonconstant) factor.
(iv) A plane curve has finitely many singular points.

## (1.4.8) tangent lines and flex points

Let $C$ be the plane projective curve defined by an irreducible homogeneous polynomial $f$. A line $L$ is tangent to $C$ at a smooth point $p$ if the intersection multiplicity of $C$ and $L$ at $p$ is at least 2 . (See 1.3.8).) There is a unique tangent line at a smooth point.

A smooth point $p$ of $C$ is a flex point if the intersection multiplicity of $C$ and its tangent line at $p$ is at least 3 , and $p$ is an ordinary flex point if the intersection multiplicity is equal to 3 .

Let $L$ be a line through a point $p$ and let $q$ be a point of $L$ distinct from $p$. We represent $p$ and $q$ by specific vectors $\left(p_{0}, p_{1}, p_{2}\right)$ and $\left(q_{0}, q_{1}, q_{2}\right)$, to write a variable point of $L$ as $p+t q$, and we expand the restriction of $f$
to $L$ in a Taylor's Series. The Taylor expansion carries over to complex polynomials because it is an identity. Let $f_{i}=\frac{\partial f}{\partial x_{i}}$ and $f_{i j}=\frac{\partial^{2} f}{\partial x_{i} \partial x_{j}}$. Then
taylor hessianmatrix texp

$$
\begin{equation*}
f(p+t q)=f(p)+\left(\sum_{i} f_{i}(p) q_{i}\right) t+\frac{1}{2}\left(\sum_{i, j} q_{i} f_{i j}(p) q_{j}\right) t^{2}+O(3) \tag{1.4.9}
\end{equation*}
$$

where the symbol $O(3)$ stands for a polynomial in which all terms have degree at least 3 in $t$. (The point $q$ is missing from this parametrization, but this won't be important.)

We rewrite this equation: Let $\nabla$ be the gradient vector $\left(f_{0}, f_{1}, f_{2}\right)$, let $H$ be the Hessian matrix of $f$, the matrix of second partial derivatives:

$$
H=\left(\begin{array}{lll}
f_{00} & f_{01} & f_{02}  \tag{1.4.10}\\
f_{10} & f_{11} & f_{12} \\
f_{20} & f_{21} & f_{22}
\end{array}\right)
$$

and let $\nabla_{p}$ and $H_{p}$ be the evaluations of $\nabla$ and $H$, respectively, at $p$. So $p$ is a smooth point of $C$ if $f(p)=0$ and $\nabla_{p} \neq 0$. Regarding $p$ and $q$ as column vectors, Equation 1.4.9 can be written as

$$
\begin{equation*}
f(p+t q)=f(p)+\left(\nabla_{p} q\right) t+\frac{1}{2}\left(q^{t} H_{p} q\right) t^{2}+O(3) \tag{1.4.11}
\end{equation*}
$$

in which $\nabla_{p} q$ and $q^{t} H_{p} q$ are to be computed as matrix products.
The intersection multiplicity of $C$ and $L$ at $p$ 1.3.8 is the lowest power of $t$ that has nonzero coefficient in $f(p+t q)$. The intersection multiplicity is at least 1 if $p$ lies on $C$, i.e., if $f(p)=0$. If $p$ is a smooth point of $C$, then $L$ is tangent to $C$ at $p$ if the coefficient $\left(\nabla_{p} q\right)$ of $t$ is zero, and $p$ is a flex point if $\left(\nabla_{p} q\right)$ and $\left(q^{t} H_{p} q\right)$ are both zero.

The equation of the tangent line $L$ at a smooth point $p$ is $\nabla_{p} x=0$, or

$$
\begin{equation*}
f_{0}(p) x_{0}+f_{1}(p) x_{1}+f_{2}(p) x_{2}=0 \tag{1.4.12}
\end{equation*}
$$

which tells us that a point $q$ lies on $L$ if the linear term in $t$ of 1.4.11) is zero.
Taylor's formula shows that the restriction of $f$ to every line through a singular point has a multiple zero. However, we will speak of tangent lines only at smooth points of the curve.

The next lemma is obtained by applying Euler's Formula to the entries of $H_{p}$ and $\nabla_{p}$.
1.4.13. Lemma. $\quad p^{t} H_{p}=(d-1) \nabla_{p} \quad$ and $\quad \nabla_{p} p=d f(p)$.

We rewrite Equation 1.4 .9 one more time, using the notation $\langle u, v\rangle$ to represent the symmetric bilinear form $u^{t} H_{p} v$ on $V=\mathbb{C}^{3}$. It makes sense to say that this form vanishes on a pair of points of $\mathbb{P}^{2}$, because the condition $\langle u, v\rangle=0$ doesn't depend on the vectors that represent those points.
1.4.14. Proposition. With notation as above,
(i) Equation (1.4.9) can be written as

$$
f(p+t q)=\frac{1}{d(d-1)}\langle p, p\rangle+\frac{1}{d-1}\langle p, q\rangle t+\frac{1}{2}\langle q, q\rangle t^{2}+O(3)
$$

(ii) A point $p$ is a smooth point of $C$ if and only if $\langle p, p\rangle=0$ but $\langle p, v\rangle$ is not identically zero.
proof. (i) This follows from Lemma 1.4.13.
(ii) $\langle p, v\rangle=(d-1) \nabla_{p} v$ isn't identically zero at a smooth point $p$ because $\nabla_{p}$ won't be zero.
1.4.15. Corollary. Let p be a smooth point of $C$, let $q$ be a point of $\mathbb{P}^{2}$ distinct from $p$. and let $L$ be the line through $p$ and $q$. Then
(i) $L$ is tangent to $C$ at $p$ if and only if $\langle p, p\rangle=\langle p, q\rangle=0$, and
(ii) $p$ is a flex point of $C$ with tangent line $L$ if and only if $\langle p, p\rangle=\langle p, q\rangle=\langle q, q\rangle=0$.
1.4.16. Theorem. A smooth point $p$ of the curve $C$ is a flex point if and only if the determinant $\operatorname{det} H_{p}$ of the Hessian matrix at p is zero.
proof. Let $p$ be a smooth point of $C$, so that $\langle p, p\rangle=0$. If $\operatorname{det} H_{p}=0$, the form $\langle u, v\rangle$ is degenerate. There is a nonzero null vector $q$, so that $\langle p, q\rangle=\langle q, q\rangle=0$. But because $\langle p, v\rangle$ isn't identically zero, $q$ is distinct from $p$. So $p$ is a flex point.

Conversely, suppose that $p$ is a flex point and let $q$ be a point on the tangent line at $p$ and distinct from $p$, so that $\langle p, p\rangle=\langle p, q\rangle=\langle q, q\rangle=0$. The restriction of the form to the two-dimensional space $W$ spanned by $p$ and $q$ is zero, and this implies that the form is degenerate. If $(p, q, v)$ is a basis of $V$ with $p, q$ in $W$, the matrix of the form will look like this:

$$
\left(\begin{array}{lll}
0 & 0 & * \\
0 & 0 & * \\
* & * & *
\end{array}\right)
$$

### 1.4.17. Proposition.

(i) Let $f(x, y, z)$ be an irreducible homogeneous polynomial of degree at least two. The Hessian determinant det $H$ isn't divisible by $f$. In particular, the Hessian determinant isn't identically zero.
(ii) A curve that isn't a line has finitely many flex points.
proof. (i) Let $C$ be the curve defined by $f$. If $f$ divides the Hessian determinant, every smooth point of $C$ will be a flex point. We set $z=1$ and look on the standard affine $\mathbb{U}^{2}$, choosing coordinates so that the origin $p$ is a smooth point of $C$, and $\frac{\partial f}{\partial y} \neq 0$ at $p$. The Implicit Function Theorem tells us that we can solve the equation $f(x, y, 1)=0$ for $y$ locally, say $y=\varphi(x)$. The graph $\Gamma:\{y=\varphi(x)\}$ will be equal to $C$ in a neighborhood of $p$ (see below). A point of $\Gamma$ is a flex point if and only if $\frac{d^{2} \varphi}{d x^{2}}$ is zero there. If this is true for all points near to $p$, then $\frac{d^{2} \varphi}{d x^{2}}$ will be identically zero, and this implies that $\varphi$ is linear: $y=a x$. Then $y=a x$ solves $f=0$, and therefore $y-a x$ divides $f(x, y, 1)$. But $f(x, y, z)$ is irreducible, and so is $f(x, y, 1)$. Therefore $f(x, y, 1)$ is linear, contrary to hypothesis.
(ii) This follows from (i) and 1.3.11). The irreducible polynomial $f$ and the Hessian determinant have finitely many common zeros.

### 1.4.18. Review. (about the Implicit Function Theorem)

Let $f(x, y)$ be a polynomial such that $f(0,0)=0$ and $\frac{d f}{d y}(0,0) \neq 0$. The Implicit Function Theorem asserts that there is a unique analytic function $\varphi(x)$, defined for small $x$, such that $\varphi(0)=0$ and $f(x, \varphi(x))$ is identically zero.

We make some further remarks. Let $\mathcal{R}$ be the ring of functions that are defined and analytic for small $x$. In the ring $\mathcal{R}[y]$ of polynomials in $y$ with coefficients in $\mathcal{R}$, the polynomial $y-\varphi(x)$, which is monic in $y$, divides $f(x, y)$. To see this, we do division with remainder of $f$ by $y-\varphi(x)$ :

$$
\begin{equation*}
f(x, y)=(y-\varphi(x)) q(x, y)+r(x) \tag{1.4.19}
\end{equation*}
$$

The quotient $q$ and remainder $r$ are in $\mathcal{R}[y]$, and $r(x)$ has degree zero in $y$, so it is in $\mathcal{R}$. Setting $y=\varphi(x)$ in the equation, one sees that $r(x)=0$.

Let $\Gamma$ be the graph of $\varphi$ in a suitable neighborhood $U$ of the origin in $x, y$-space. Since $f(x, y)=(y-$ $\varphi(x)) q(x, y)$, the locus $f(x, y)=0$ in $U$ has the form $\Gamma \cup \Delta$, where $\Gamma$, the zero lous of $y-\varphi(x)$, is the graph of $\varphi$ and $\Delta$ is the zero locus of $q(x, y)$. Differentiating, we find that $\frac{\partial f}{\partial y}(0,0)=q(0,0)$. So $q(0,0) \neq 0$. Then $\Delta$ doesn't contain the origin, while $\Gamma$ does. This implies that $\Delta$ is disjoint from $\Gamma$, locally. A sufficiently small neighborhood $U$ of the origin won't contain any of points $\Delta$. In such a neighborhood, the locus of zeros of $f$ will be $\Gamma$.

### 1.5 Transcendence Degree

Let $F \subset K$ be a field extension. A set $\alpha=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ of elements of $K$ is algebraically dependent over $F$ if there is a nonzero polynomial $f\left(x_{1}, \ldots, x_{n}\right)$ with coefficients in $F$, such that $f(\alpha)=0$. If there is no such polynomial, the set $\alpha$ is algebraically independent over $F$.

An infinite set is called algebraically independent if every finite subset is algebraically independent, in other words, if there is no polynomial relation among any finite set of its elements.

The set $\left\{\alpha_{1}\right\}$ consisting of a single element of $K$ will be algebraically dependent if $\alpha_{1}$ is algebraic over $F$. Otherwise, it will be algebraically independent, and then $\alpha_{1}$ is said to be transcendental over $F$.

An algebraically independent set $\alpha=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ that isn't contained in a larger algebraically independent set is a transcendence basis for $K$ over $F$. If there is a finite transcendence basis, its order is the transcendence degree of the field extension $K$ of $F$. Lemma 1.5 .2 below shows that all transcendence bases for $K$ over $F$ have the same order, so the transcendence degree is well-defined. If there is no finite transcendence basis, the transcendence degree of $K$ over $F$ is infinite.

For example, let $K=F\left(x_{1}, \ldots, x_{n}\right)$ be the field of rational functions in $n$ variables. The variables $x_{i}$ form a transcendence basis of $K$ over $F$, and the transcendence degree of $K$ over $F$ is $n$.

A domain is a nonzero ring with no zero divisors, and a domain that contains a field $F$ as a subring is called an $F$-algebra. We use the customary notation $F\left[\alpha_{1}, \ldots, \alpha_{n}\right]$ or $F[\alpha]$ for the $F$-algebra generated by a set $\alpha=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$, and we may denote its field of fractions by $F\left(\alpha_{1}, \ldots, \alpha_{n}\right)$ or by $F(\alpha)$. The set $\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ is algebraically independent over $F$ if and only if the surjective map from the polynomial algebra $F\left[x_{1}, \ldots, x_{n}\right]$ to $F\left[\alpha_{1}, \ldots, \alpha_{n}\right]$ that sends $x_{i}$ to $\alpha_{i}$ is bijective.
1.5.1. Lemma. Let $K / F$ be a field extension, let $\alpha=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ be a set of elements of $K$ that is algebraically independent over $F$, and let $F(\alpha)$ be the field of fractions of $F[\alpha]$.
(i) Every element of the field $F(\alpha)$ that isn't in $F$ is transcendental over $F$.
(ii) If $\beta$ is another element of $K$, the set $\left\{\alpha_{1}, \ldots, \alpha_{n}, \beta\right\}$ is algebraically dependent if and only if $\beta$ is algebraic over $F(\alpha)$.
(iii) The algebraically independent set $\alpha$ is a transcendence basis if and only if every element of $K$ is algebraic over $F(\alpha)$.
proof. (i) We an write an element $z$ of $F(\alpha)$ as a fraction $p / q=p(\alpha) / q(\alpha)$, where $p(x)$ and $q(x)$ are relatively prime polynomials. Suppose that $z$ satisfies a nontrivial polynomial relation $c_{0} z^{n}+c_{1} z^{n-1}+\cdots+c_{n}=0$ with $c_{i}$ in $F$. We may assume that $c_{0}=1$. Substituting $z=p / q$ and multiplying by $q^{n}$ gives us the equation

$$
p^{n}=-q\left(c_{1} p^{n-1}+\cdots+c_{n} q^{n-1}\right)
$$

By hypothesis, $\alpha$ is an algebraically independent set, so this equation is equivalent with a polynomial equation in $F[x]$. It shows that $q$ divides $p^{n}$, which contradicts the hypothesis that $p$ and $q$ are relatively prime. So $z$ satisfies no polynomial relation, and therefore it is transcendental.

The other assertions are left as an exercise.
1.5.2. Lemma.
(i) Let $K / F$ be a field extension. If $K$ has a finite transcendence basis, then all algebraically independent subsets of $K$ are finite, and all transcendence bases have the same order.
(ii) If $L \supset K \supset F$ are fields and if the degree $[L: K]$ of $L$ over $K$ is finite, then $K$ and $L$ have the same transcendence degree over $F$.
proof. (i) Let $\alpha=\left\{\alpha_{1}, \ldots, \alpha_{r}\right\}$ and $\beta=\left\{\beta_{1}, \ldots, \beta_{s}\right\}$. Assume that $K$ is algebraic over $F(\alpha)$ and that the set $\beta$ is algebraically independent. We show that $s \leq r$. The fact that all transcendence bases have the same order will follow: If both $\alpha$ and $\beta$ are transcendence bases, then $s \leq r$, and since we can interchange $\alpha$ and $\beta$, $r \leq s$.

The proof that $s \leq r$ proceeds by reducing to the trivial case that $\beta$ is a subset of $\alpha$. Suppose that some element of $\beta$, say $\beta_{s}$, isn't in the set $\alpha$. The set $\beta^{\prime}=\left\{\beta_{1}, \ldots, \beta_{s-1}\right\}$ is algebraically independent, but it isn't a transcendence basis. So $K$ isn't algebraic over $F\left(\beta^{\prime}\right)$. Since $K$ is algebraic over $F(\alpha)$, there is at least one element of $\alpha$, say $\alpha_{r}$, that isn't algebraic over $F\left(\beta^{\prime}\right)$. Then $\gamma=\beta^{\prime} \cup\left\{\alpha_{s}\right\}$ will be an algebraically independent set of order $s$, and it contains more elements of the set $\alpha$ than $\beta$ does. Induction shows that $s \leq r$.

### 1.6 The Dual Curve

## (1.6.1) the dual plane

Let $\mathbb{P}$ denote the projective plane with coordinates $x_{0}, x_{1}, x_{2}$, and let $L$ be the line in $\mathbb{P}$ with the equation

$$
\begin{equation*}
s_{0} x_{0}+s_{1} x_{1}+s_{2} x_{2}=0 \tag{1.6.2}
\end{equation*}
$$

The solutions of this equation determine the coefficients $s_{i}$ only up to a common nonzero scalar factor, so $L$ determines a point $\left(s_{0}, s_{1}, s_{2}\right)$ in another projective plane $\mathbb{P}^{*}$ called the dual plane. We denote that point by $L^{*}$. Moreover, a point $p=\left(x_{0}, x_{1}, x_{2}\right)$ in $\mathbb{P}$ determines a line in the dual plane, the line with the equation (1.6.2), when $s_{i}$ are regarded as the variables and $x_{i}$ as the scalar coefficients. We denote that line by $p^{*}$. The equation exhibits a duality between $\mathbb{P}$ and $\mathbb{P}^{*}$. A point $p$ of $\mathbb{P}$ lies on the line $L$ if and only if the equation is satisfied, and this means that, in $\mathbb{P}^{*}$, the point $L^{*}$ lies on the line $p^{*}$.

## (1.6.3) the dual curve

Let $C$ be a plane projective curve of degree at least two, and let $U$ be the set of its smooth points. Corollary 1.4.7) tells us that $U$ is the complement of a finite subset of $C$. We define a map

$$
U \xrightarrow{t} \mathbb{P}^{*}
$$

as follows: Let $p$ be a point of $U$ and let $L$ be the tangent line to $C$ at $p$. Then $t(p)=L^{*}$, where $L^{*}$ is the point of $\mathbb{P}^{*}$ that corresponds to the line $L$.

Denoting the partial derivative $\frac{\partial f}{\partial x_{i}}$ by $f_{i}$ as before, the tangent line $L$ at a smooth point $p=\left(x_{0}, x_{1}, x_{2}\right)$ of $C$ has the equation $f_{0} x_{0}+f_{1} x_{1}+f_{2} x_{2}=0$ 1.4.12. Therefore $L^{*}$ is the point

$$
\begin{equation*}
\left(s_{0}, s_{1}, s_{2}\right)=\left(f_{0}(x), f_{1}(x), f_{2}(x)\right) \tag{1.6.4}
\end{equation*}
$$

We'll drop some parentheses, denoting the image $t(U)$ of $U$ in $\mathbb{P}^{*}$ by $t U$. Points of $t U$ correspond to tangent lines at smooth points of $C$. We assume that $C$ has degree at least two because, if $C$ were a line, $t U$ would be a point. Since the partial derivatives have no common factor, the tangent lines aren't constant when the degree is two or more.

## ??figure??

1.6.5. Lemma. Let $\varphi\left(s_{0}, s_{1}, s_{2}\right)$ be a homogeneous polynomial, and let $g\left(x_{0}, x_{1}, x_{2}\right)=\varphi\left(f_{0}(x), f_{1}(x), f_{2}(x)\right)$. Then $\varphi(s)$ is identically zero on $t U$ if and only if $g(x)$ is identically zero on $U$. This is true if and only if $f$ divides $g$.
proof. The first assertion follows from the fact that $\left(s_{0}, s_{1}, s_{2}\right)$ and $\left(f_{0}(x), f_{1}(x), f_{2}(x)\right)$ represent the same point of $\mathbb{P}^{*}$, and the last one follows from Corollary 1.3 .15
1.6.6. Theorem. Let $C$ be the plane curve defined by an irreducible homogeneous polynomial $f$ of degree $d$ at least two. With notation as above, the image $t U$ is contained in a curve $C^{*}$ in the dual space $\mathbb{P}^{*}$.

The curve $C^{*}$ referred to in the theorem is the dual curve.
proof. If an irreducible homogeneous polynomial $\varphi(s)$ vanishes on $t U$, it will be unique up to scalar factor (Corollary 1.3.15).

Let's use vector notation: $x=\left(x_{0}, x_{1}, x_{2}\right), s=\left(s_{0}, s_{1}, s_{2}\right)$, and $\nabla f=\left(f_{0}, f_{1}, f_{2}\right)$. We show first that there is a nonzero polynomial $\varphi(s)$, not necessarily irreducible or homogeneous, that vanishes on $t U$. The field $\mathbb{C}\left(x_{0}, x_{1}, x_{2}\right)$ has transcendence degree three over $\mathbb{C}$. Therefore the four polynomials $f_{0}, f_{1}, f_{2}$, and $f$ are algebraically dependent. There is a nonzero polynomial $\psi\left(s_{0}, s_{1}, s_{2}, t\right)$ such that $\psi\left(f_{0}(x), f_{1}(x), f_{2}(x), f(x)\right)=$ $\psi(\nabla f(x), f(x))$ is the zero polynomial. We can cancel factors of $t$, so we may assume that $\psi$ isn't divisible by $t$. Let $\varphi(s)=\psi\left(s_{0}, s_{1}, s_{2}, 0\right)$. This isn't the zero polynomial because $t$ doesn't divide $\psi$.

Let $\bar{x}=\left(\bar{x}_{1}, \bar{x}_{2}, \bar{x}_{3}\right)$ be a vector that represents a point of $U$. Then $f(\bar{x})=0$, and therefore

$$
\psi(\nabla f(\bar{x}), f(\bar{x}))=\psi(\nabla f(\bar{x}), 0)=\varphi(\nabla f(\bar{x}))
$$

Since $\psi(\nabla f(x), f(x))$ is identically zero, $\varphi(\nabla f(\bar{x}))=0$ for all $\bar{x}$ in $U$.
exampledualone
exampledualtwo
equationdualthree

Next, because the vectors $\bar{x}$ and $\lambda \bar{x}$ represent the same point of $U, \varphi(\nabla f(\lambda \bar{x}))=0$. Since $f$ has degree $d$, the derivatives $f_{i}$ are homogeneous of degree $d-1$. Therefore $\varphi(\nabla f(\lambda \bar{x}))=\varphi\left(\lambda^{d-1} \nabla f(\bar{x})\right)=0$ for all $\lambda$. Because the scalar $\lambda^{d-1}$ can be any complex number, Lemma 1.3.2 tells us that the homogeneous parts of $\varphi(\nabla f(x))$ vanish for all $x \in U$. The homogeneous parts of degree $r$ of $\varphi(s)$ correspond to the homogenous parts of degree $r(d-1)$ of $\varphi(\nabla f(x))$. So the homogeneous parts of $\varphi(s)$ vanish on $t U$. This shows that there is a homogeneous polynomial $\varphi(s)$ that vanishes on $t U$. We choose such a polynomial $\varphi(s)$. Let its degree be $r$.

If $f$ has degree $d$, the polynomial $g(x)=\varphi(\nabla f(x))$ will be homogeneous, of degree $r(d-1)$. It will vanish on $U$, and therefore on $C$ 1.3.20. So $f$ will divide $g$. Finally, if $\varphi(s)$ factors, then $g(x)$ factors accordingly, and because $f$ is irreducible, it will divide one of the factors of $g$. The corresponding factor of $\varphi$ will vanish on $t U$ 1.6.5. So we may replace the homogeneous polynomial $\varphi$ by one of its irreducible factors.

In principle, the proof of Theorem 1.6.6 gives a method for finding a polynomial that vanishes on the dual curve. One looks for a polynomial relation among $f_{x}, f_{y}, f_{z}, f$, and then sets $f=0$. But it is usually painful to determine the defining polynomial of $C^{*}$ explicitly. Most often, the degrees of $C$ and $C^{*}$ will be different, and several points of the dual curve $C^{*}$ may correspond to a singular point of $C$, and vice versa.

However, computation is easy for a conic.

### 1.6.7. Examples.

(i) (the dual of a conic) Let $f=x_{0} x_{1}+x_{0} x_{2}+x_{1} x_{2}$ and let $C$ be the conic $f=0$. Let $\left(s_{0}, s_{1}, s_{2}\right)=$ $\left(f_{0}, f_{1}, f_{2}\right)=\left(x_{1}+x_{2}, x_{0}+x_{2}, x_{0}+x_{1}\right)$. Then
(1.6.8) $\quad s_{0}^{2}+s_{1}^{2}+s_{2}^{2}-2\left(x_{0}^{2}+x_{1}^{2}+x_{2}^{2}\right)=2 f \quad$ and $\quad s_{0} s_{1}+s_{1} s_{2}+s_{0} s_{2}-\left(x_{0}^{2}+x_{1}^{2}+x_{2}^{2}\right)=3 f$

We eliminate $\left(x_{0}^{2}+x_{1}^{2}+x_{2}^{2}\right)$ from the two equations.

$$
\begin{equation*}
\left(s_{0}^{2}+s_{1}^{2}+s_{2}^{2}\right)-2\left(s_{0} s_{1}+s_{1} s_{2}+s_{0} s_{2}\right)=-4 f \tag{1.6.9}
\end{equation*}
$$

Setting $f=0$ gives us the equation of the dual curve. It is another conic.
(ii) (the dual of a cuspidal cubic) It is too much work to computethe dual of a smooth cubic, which is a curve of degree 6 . We compute the dual of a cubic with a cusp instead. The curve $C$ defined by the irreducible polynomial $f=y^{2} z+x^{3}$ has a cusp at $(0,0,1)$. The Hessian matrix of $f$ is

$$
H=\left(\begin{array}{ccc}
6 x & 0 & 0 \\
0 & 2 z & 2 y \\
0 & 2 y & 0
\end{array}\right)
$$

and the Hessian determinant $h=\operatorname{det} H$ is $-24 x y^{2}$. The common zeros of $f$ and $h$ are the cusp point $(0,0,1)$ and a single flex point $(0,1,0)$.

We scale the partial derivatives of $f$ to simplify notation. Let $u=f_{x} / 3=x^{2}, v=f_{y} / 2=y z$, and $w=f_{z}=y^{2}$. Then

$$
v^{2} w-u^{3}=y^{4} z^{2}-x^{6}=\left(y^{2} z+x^{3}\right)\left(y^{2} z-x^{3}\right)=f\left(y^{2} z-x^{3}\right)
$$

The zero locus of the irreducible polynomial $v^{2} w-u^{3}$ is the dual curve. It is another cuspidal cubic.

## (1.6.10) a local equation for the dual curve

We label the coordinates in $\mathbb{P}$ and $\mathbb{P}^{*}$ as $x, y, z$ and $u, v, w$, respectively, and we work in a neighborhood of a smooth point $p_{0}$ of the curve $C$ defined by a homogeneous polynomial $f(x, y, z)$, choosing coordinates so that $p_{0}=(0,0,1)$, and that the tangent line at $p_{0}$ is the line $L_{0}:\{y=0\}$. The image of $p_{0}$ in the dual curve $C^{*}$ is $L_{0}^{*}:(u, v, w)=(0,1,0)$.

Let $\tilde{f}(x, y)=f(x, y, 1)$. In the affine $x, y$-plane, the point $p_{0}$ becomes the origin $p_{0}=(0,0)$. So $\widetilde{f}\left(p_{0}\right)=0$, and since the tangent line is $L_{0}, \frac{\partial \widetilde{f}}{\partial x}\left(p_{0}\right)=0$, while $\frac{\partial \widetilde{f}}{\partial y}\left(p_{0}\right) \neq 0$. We solve the equation $\widetilde{f}=0$ for $y$ as an analytic function $y(x)$ for small $x$, with $y(0)=0$. Let $y^{\prime}(x)$ denote the derivative $\frac{d y}{d x}$. Differentiating the equation $f(x, y(x))=0$ shows that $y^{\prime}(0)=0$.

Let $\widetilde{p}_{1}=\left(x_{1}, y_{1}\right)$ be a point of $C_{0}$ near to $\widetilde{p}_{0}$, so that $y_{1}=y\left(x_{1}\right)$, and let $y_{1}^{\prime}=y^{\prime}\left(x_{1}\right)$. The tangent line $L_{1}$ at $\widetilde{p}_{1}$ has the equation

$$
\begin{equation*}
y-y_{1}=y_{1}^{\prime}\left(x-x_{1}\right) \tag{1.6.11}
\end{equation*}
$$

Putting $z$ back, the homogeneous equation of the tangent line $L_{1}$ at the point $p_{1}=\left(x_{1}, y_{1}, 1\right)$ is

$$
-y_{1}^{\prime} x+y+\left(y_{1}^{\prime} x_{1}-y_{1}\right) z=0
$$

The point $L_{1}^{*}$ of the dual plane that corresponds to $L_{1}$ is $\left(-y_{1}^{\prime}, 1, y_{1}^{\prime} x_{1}-y_{1}\right)$.
Let's drop the subscript 1 . As $x$ varies, and writing $y=y(x)$ and $y^{\prime}=y^{\prime}(x)$,

$$
\begin{equation*}
(u, v, w)=\left(-y^{\prime}, 1, y^{\prime} x-y\right) \tag{1.6.12}
\end{equation*}
$$

There may be accidents: $L_{0}$ might be tangent to $C$ at distinct smooth points $q_{0}$ and $p_{0}$, or it might pass through a singular point of $C$. If either of these accidents occurs, we can't analyze the neighborhood of $L_{0}^{*}$ in $C^{*}$ by this method. But, provided that there are no accidents, the path $\sqrt{1.6 .12}$ will trace out the dual curve $C^{*}$ near to $L_{0}^{*}=(0,1,0)$. (See 1.4.18.)

## (1.6.13) the bidual

The bidual $C^{* *}$ of a curve $C$ is the dual of the curve $C^{*}$. It is a curve in the space $\mathbb{P}^{* *}$, which is $\mathbb{P}$.
1.6.14. Theorem. A plane curve of degree greater than one is equal to its bidual: $C^{* *}=C$.

Let $C$ be a plane curve. Weuse the following notation:

- $U$ is the set of smooth points of $C$, as above, and $U^{*}$ is the set of smooth points of $C^{*}$.
- $U^{*} \xrightarrow{t^{*}} \mathbb{P}^{* *}=\mathbb{P}$ is the map analogous to the map $U \xrightarrow{t} \mathbb{P}^{*}$.
- $V$ is the set of smooth points $p$ of $C$ such that $t(p)$ is a smooth point of $C^{*}$, and $V^{*}$ is the image $t(V)$ of $V$.

Thus $V \subset U \subset C$ and $V^{*} \subset U^{*} \subset C^{*}$.

### 1.6.15. Lemma.

(i) $V$ is the complement of a finite set of points of $C$.
(ii) Let $p_{1}$ be a point near to a smooth point $p$ of a curve $C$, let $L_{1}$ and $L$ be the tangent line to $C$ at $p_{1}$ and $p$, respectively, and let $q$ be intersection point $L_{1} \cap L$. Then $\lim _{p_{1} \rightarrow p_{0}} q=p$.
(iii) If $p$ is a point of $V$ with tangent line $L$, the tangent line to $C^{*}$ at $L^{*}$ is $p^{*}$.
(iv) If $L$ is the tangent line at a point $p$ of $V$, then $t(p)=L^{*} \in V^{*}$ and $t^{*}\left(L^{*}\right)=p$.
proof. (i) Let $S$ and $S^{*}$ denote the finite sets of singular points of $C$, and $C^{*}$, respectively, the set $U$ of smooth points of $C$ is the complement of $S$ in $C$, and $V$ is obtained from $U$ by deleting points whose images are in $S^{*}$. The fibre of $t$ over a point $L^{*}$ of $C^{*}$ is the set of smooth points $p$ of $C$ such that the tangent line at $p$ is $L$. Since $L$ meets $C$ in finitely many points, the fibre is finite. So the inverse image of $S^{*} \cap U$ will be a finite subset of $U$.
(ii) We work analytically in a neighborhood of $p$, choosing coordinates so that $p=(0,0,1)$ and that $L$ is the line $\{y=0\}$. Let $\left(x_{q}, y_{q}, 1\right)$ be the coordinates of $q=L \cap L_{1}$. Since $q$ is a point of $L, y_{q}=0$. The coordinate $x_{q}$ can be obtained by substituting $x=x_{q}$ and $y=0$ into the equation 1.6.11) of $L_{1}$ :

$$
x_{q}=x_{1}-y_{1} / y_{1}^{\prime} .
$$

When a function has an $n$th order zero at the point $x=0$, i.e, when it has the form $y=x^{n} h(x)$, where $n>0$ and $h(0) \neq 0$, the order of zero of its derivative at that point is $n-1$. This is verified by differentiating
$x^{n} h(x)$. Since the function $y(x)$ has a zero of positive order at $p, \lim _{p_{1} \rightarrow p_{0}} y_{1} / y_{1}^{\prime}=0$. We also have $\lim _{p_{1} \rightarrow p_{0}} x_{1}=$ 0 . So $\lim _{p_{1} \rightarrow p_{0}} x_{q}=0$ and $\lim _{p_{1} \rightarrow p_{0}} q=\lim _{p_{1} \rightarrow p_{0}}\left(x_{q}, y_{q}, 1\right)=(0,0,1)=p$.

## figure

(iii) Let $p_{1}$ be a point of $C$ near to $p$, and let $L_{1}$ be the tangent line to $C$ at $p_{1}$. The image of $p_{1}$ is $L_{1}^{*}=$ $\left(f_{0}\left(p_{1}\right), f_{1}\left(p_{1}\right), f_{2}\left(p_{1}\right)\right)$. Because the partial derivatives $f_{i}$ are continuous,

$$
\lim _{p_{1} \rightarrow p_{0}} L_{1}^{*}=\left(f_{0}(p), f_{1}(p), f_{2}(p)\right)=L^{*}
$$

Let $q=L \cap L_{1}$. Then $q^{*}$ is the line through the points $L^{*}$ and $L_{1}^{*}$. As $p_{1}$ approaches $p, L_{1}^{*}$ approaches $L^{*}$, and therefore $q^{*}$ approaches the tangent line to $C^{*}$ at $L^{*}$. On the other hand, (ii) tells us that $q^{*}$ approaches $p^{*}$. Therefore the tangent line at $L^{*}$ is $p^{*}$.
(iv) Since $V \subset U, t(p)=L^{*}$ by definition of $t$, and if $p \in V$, then $t(p) \in V^{*}$ by definition of $V^{*}$. Since $L^{*}$ is a point of $V^{*}$, and since the definition of $t^{*}$ is analogous to $t, t^{*}\left(L^{*}\right)$ is the tangent line to $C^{*}$ at $L^{*}$, which, by (iii), is $p^{*}$.
proof of theorem 1.6.14 Let $V$ be the subset of $C$ defined above, let $U^{*}$ be the set of smooth points of $C^{*}$, and let $U^{*} \xrightarrow{t^{*}} \mathbb{P}^{* *}=\mathbb{P}$ be the map analogous to the map $U \xrightarrow{t} \mathbb{P}^{*}$. For all points $p$ of $V$, the map $t^{*}$ is $t^{*}\left(L^{*}\right)=\left(p^{*}\right)^{*}=p$. Thus $t^{*} t(p)=p$. It follows that the restriction of $t$ to $V$ is injective, and that it defines a bijective map from $V$ to its image $t V$, whose inverse function is $t^{*}$. So $V$ is contained in the bidual $C^{* *}$. Since $V$ is dense in $C$ and $C^{* *}$ is a closed set, $C \subset C^{* *}$. Since $C$ and $C^{* *}$ are curves, $C=C^{* *}$.
1.6.16. Corollary. (i) Let $U$ be the set of smooth points of a plane curve $C$, and let denote the map from $U$ to the dual curve $C^{*}$. The image $t U$ of $U$ is the complement of a finite subset of $C^{*}$.
(ii) If $C$ is smooth, the map $C \xrightarrow{t} C^{*}$, which is defined at all points of $C$, is surjective.
proof. (i) Let $U$ and $V$ be as above, and let $U^{*}$ be the set of smooth points of $C^{*}$. The image $t V$ of $V$ is contained in $U^{*}$. Then $V=t^{*} t V \subset t^{*} U^{*} \subset C^{* *}=C$. Since $V$ is the complement of a finite subset of $C$, $t^{*} U^{*}$ is q a finite subset of $C$ too. The assertion to be proved follows when we switch $C$ and $C^{*}$.
(ii) The map $t$ is continuous, so its image $t C$ is a commpact subset of $C^{*}$, and by (i), its complement $S$ is a finite set. Therefore $S$ is both open and closed. It consists of isolated points of $C^{*}$. Since a plane curve has no isolated point (1.3.18), $S$ is empty.

### 1.7 Resultants and Discriminants

$$
\begin{equation*}
F(x)=x^{m}+a_{1} x^{m-1}+\cdots+a_{m} \quad \text { and } \quad G(x)=x^{n}+b_{1} x^{n-1}+\cdots+b_{n} \tag{1.7.1}
\end{equation*}
$$

The resultant $\operatorname{Res}(F, G)$ of $F$ and $G$ is a certain polynomial in the coefficients. Its important property is that, when the coefficients of are in a field, the resultant is zero if and only if $F$ and $G$ have a common factor.

As an example, suppose that the coefficients $a_{i}$ and $b_{i}$ in 1.7.1) are polynomials in $t$, so that $F$ and $G$ become polynomials in two variables. Let $C$ and $D$ be (possibly reducible) curves $F=0$ and $G=0$ in the affine plane $\mathbb{A}_{t, x}^{2}$, and let $S$ be the set of intersections: $S=C \cap D$. The resultant $\operatorname{Res}(F, G)$, computed regarding $x$ as the variable, will be a polynomial in $t$ whose roots are the $t$-coordinates of the set $S$.

## figure

The analogous statement is true when there are more variables. For example, if $F$ and $G$ are polynomials in $x, y, z$, the loci $C:\{F=0\}$ and $D:\{G=0\}$ in $\mathbb{A}^{3}$ will be surfaces, and $S=C \cap D$ will be a curve. The resultant $\operatorname{Res}_{z}(F, G)$, computed regarding $z$ as the variable, is a polynomial in $x, y$. Its zero locus in the plane $\mathbb{A}_{x y}^{2}$ is the projection of $S$ to the plane.

The formula for the resultant is nicest when one allows leading coefficients different from 1 . We work with homogeneous polynomials in two variables to prevent the degrees from dropping when a leading coefficient happens to be zero.

Let $f$ and $g$ be homogeneous polynomials in $x, y$ with complex coefficients:

$$
\begin{equation*}
f(x, y)=a_{0} x^{m}+a_{1} x^{m-1} y+\cdots+a_{m} y^{m}, \quad g(x, y)=b_{0} x^{n}+b_{1} x^{n-1} y+\cdots+b_{n} y^{n} \tag{1.7.2}
\end{equation*}
$$

Suppose that they have a common zero $(x, y)=(u, v)$ in $\mathbb{P}_{x y}^{1}$. Then $v x-u y$ divides both $g$ and $f$. The polynomial $h=f g /(v x-u y)$ has degree $r=m+n-1$, and it will be divisible by $f$ and by $g$, say $h=p f=q g$, where $p$ and $q$ are homogeneous polynomials of degrees $n-1$ and $m-1$, respectively. Then $h$ will be a linear combination $p f$ of the polynomials $x^{i} y^{j} f$, with $i+j=n-1$, and it will also be a linear combination $q g$ of the polynomials $x^{k} y^{\ell} g$, with $k+\ell=m-1$. The equation $p f=q g$ tells us that the $r+1$ polynomials of degree $r$,

$$
\begin{equation*}
x^{n-1} f, x^{n-2} y f, \ldots, y^{n-1} f ; x^{m-1} g, x^{m-2} y g, \ldots, y^{m-1} g \tag{1.7.3}
\end{equation*}
$$

will be dependent. For example, suppose that $f$ has degree 3 and $g$ has degree 2 . If $f$ and $g$ have a common zero, the polynomials

$$
\begin{array}{lr}
x f= & a_{0} x^{4}+a_{1} x^{3} y+a_{2} x^{2} y^{2}+a_{3} x y^{3} \\
y f= & a_{0} x^{3} y+a_{1} x^{2} y^{2}+a_{2} x y^{3}+a_{3} y^{4} \\
x^{2} g= & b_{0} x^{4}+b_{1} x^{3} y+b_{2} x^{2} y^{2} \\
x y g= & b_{0} x^{3} y+b_{1} x^{2} y^{2}+b_{2} x y^{3} \\
y^{2} g= & b x^{2} y^{2}+b_{1} x y^{3}+b_{2} y^{4}
\end{array}
$$

will be dependent. Conversely, if the polynomials 1.7.3) are dependent, there will be an equation of the form $p f=q g$, with $p$ of degree $n-1$ and $q$ of degree $m-1$. Then at least one zero of $g$ must also be a zero of $f$.

The polynomials 1.7 .3 have degree $r(=m+n-1)$. We form a square $(r+1) \times(r+1)$ matrix $\mathcal{R}$, the resultant matrix, whose columns are indexed by the monomials $x^{r}, x^{r-1} y, \ldots, y^{r}$ of degree $r$, and whose rows list the coefficients of the polynomials (1.7.3). The matrix is illustrated below for the cases $m, n=3,2$ and $m, n=1,2$, with dots representing entries that are zero:

$$
\mathcal{R}=\left(\begin{array}{ccccc}
a_{0} & a_{1} & a_{2} & a_{3} & \cdot  \tag{1.7.4}\\
\cdot & a_{0} & a_{1} & a_{2} & a_{3} \\
b_{0} & b_{1} & b_{2} & \cdot & \cdot \\
\cdot & b_{0} & b_{1} & b_{2} & \cdot \\
\cdot & \cdot & b_{0} & b_{1} & b_{2}
\end{array}\right) \quad \text { or } \quad \mathcal{R}=\left(\begin{array}{ccc}
a_{0} & a_{1} & \cdot \\
\cdot & a_{0} & a_{1} \\
b_{0} & b_{1} & b_{2}
\end{array}\right)
$$

The resultant of $f$ and $g$ is defined to be the determinant of $\mathcal{R}$.

$$
\begin{equation*}
\operatorname{Res}(f, g)=\operatorname{det} \mathcal{R} \tag{1.7.5}
\end{equation*}
$$

Here, the coefficients of $f$ and $g$ can be in any ring.
The resultant $\operatorname{Res}(F, G)$ of the monic, one-variable polynomials $F(x)=x^{m}+a_{1} x^{m-1}+\cdots+a_{m}$ and $G(x)=x^{n}+b_{1} x^{n-1}+\cdots+b_{n}$ is the determinant of the matrix $\mathcal{R}$, with $a_{0}=b_{0}=1$.
1.7.6. Corollary. Let $f$ and $g$ be homogeneous polynomials in two variables, or monic polynomials in one variable, of degrees $m$ and $n$, respectively, and with coefficients in a field. The resultant $\operatorname{Res}(f, g)$ is zero if and only if $f$ and $g$ have a common factor. If so, there will be polynomials $p$ and $q$ of degrees $n-1$ and $m-1$ respectively, such that $p f=q g$. If the coefficients are complex numbers, the resultant is zero if and only if $f$ and $g$ have a common root.

When the leading coefficients $a_{0}$ and $b_{0}$ of $f$ and $g$ are both zero, the point $(1,0)$ of $\mathbb{P}_{x y}^{1}$ will be a zero of $f$ and of $g$. In this case, one could say that $f$ and $g$ have a common zero at infinity.

When defining the degree of a polynomial, one may assign an integer called a weight to each variable. If one assigns weight $w_{i}$ to the variable $x_{i}$, the monomial $x_{1}^{e_{1}} \cdots x_{n}^{e_{n}}$ gets the weighted degree

$$
e_{1} w_{1}+\cdots+e_{n} w_{n}
$$

For instance, it is natural to assign weight $k$ to the coefficient $a_{k}$ of the polynomial $f(x)=x^{n}-a_{1} x^{n-1}+$ $a_{2} x^{n-2}-\cdots \pm a_{n}$ because, if $f$ factors into linear factors, $f(x)=\left(x-\alpha_{1}\right) \cdots\left(x-\alpha_{n}\right)$, then $a_{k}$ will be the $k$ th elementary symmetric function in $\alpha_{1}, \ldots, \alpha_{n}$. When written as a polynomial in $\alpha$, the degree of $a_{k}$ will be $k$.

We leave the proof of the next lemma as an exercise.
degresult
resroots
restriviali-
ties
discrim-
sect
discrdef
1.7.8. Lemma. Let $f(x, y)$ and $g(x, y)$ be homogeneous polynomials of degrees $m$ and $n$ respectively, with variable coefficients $a_{i}$ and $b_{i}$, as in (1.7.2). When one assigns weight $i$ to $a_{i}$ and to $b_{i}$, the resultant $\operatorname{Res}(f, g)$ becomes a weighted homogeneous polynomial of degree $m n$ in the variables $\left\{a_{i}, b_{j}\right\}$.
1.7.9. Proposition. Let $F$ and $G$ be products of monic linear polynomials, say $F=\prod_{i}\left(x-\alpha_{i}\right)$ and $G=$ $\prod_{j}\left(x-\beta_{j}\right)$. Then

$$
\operatorname{Res}(F, G)=\prod_{i, j}\left(\alpha_{i}-\beta_{j}\right)=\prod_{i} G\left(\alpha_{i}\right)
$$

Note. Since the resultant vanishes when $\alpha_{i}=\beta_{j}$, it must be divisible by $\alpha_{i}-\beta_{j}$. So its weighted degree, though rather large, is as small as it could be.
proof. The equality of the second and third terms is obtained by substituting $\alpha_{i}$ for $x$ into the formula $G=$ $\Pi\left(x-\beta_{j}\right)$. We prove that the first and second terms are equal.

Let the elements $\alpha_{i}$ and $\beta_{j}$ be variables, let $R$ denote the resultant $\operatorname{Res}(F, G)$ and let $\Pi$ denote the product $\prod_{i . j}\left(\alpha_{i}-\beta_{j}\right)$. When we write the coefficients of $F$ and $G$ as symmetric functions in the roots $\alpha_{i}$ and $\beta_{j}, R$ will be homogeneous. Its (unweighted) degree in $\left\{\alpha_{i}, \beta_{j}\right\}$ will be $m n$, the same as the degree of $\Pi$ (Lemma 1.7.8). To show that $R=\Pi$, we choose $i, j$ and divide $R$ by the polynomial $\alpha_{i}-\beta_{j}$, considered as a monic polynomial in $\alpha_{i}$ :

$$
R=\left(\alpha_{i}-\beta_{j}\right) q+r
$$

where $r$ has degree zero in $\alpha_{i}$. The resultant $R$ vanishes when we substitute $\alpha_{i}=\beta_{j}$. Looking at this equation, we see that the remainder $r$ also vanishes when $\alpha_{i}=\beta_{j}$. On the other hand, the remainder is independent of $\alpha_{i}$. It doesn't change when we set $\alpha_{i}=\beta_{j}$. Therefore the remainder is zero, and $\alpha_{i}-\beta_{j}$ divides $R$. This is true for all $i$ and all $j$, so $\Pi$ divides $R$, and since these two polynomials have the same degree, $R=c \Pi$ for some scalar $c$. To show that $c=1$, one computes $R$ and $\Pi$ for some particular polynomials. We suggest doing the computation with $F=x^{m}$ and $G=x^{n}-1$.
1.7.10. Corollary. Let $F, G, H$ be monic polynomials and let $c$ be a scalar. Then
(i) $\operatorname{Res}(F, G H)=\operatorname{Res}(F, G) \operatorname{Res}(F, H)$, and
(ii) $\operatorname{Res}(F(x-c), G(x-c))=\operatorname{Res}(F(x), G(x))$.

## (1.7.11) the discriminant

The discriminant $\operatorname{Discr}(F)$ of a polynomial $F=a_{0} x^{m}+a_{1} x^{n-1}+\cdots a_{m}$ is the resultant of $F$ and its derivative $F^{\prime}$ :

$$
\begin{equation*}
\operatorname{Discr}(F)=\operatorname{Res}\left(F, F^{\prime}\right) \tag{1.7.12}
\end{equation*}
$$

The computation of the discriminant is made using the formula for the resultant of a polynomial of degree $m$. It will be a weighted polynomial of degree $m(m-1)$. The definition makes sense when the leading coefficient $a_{0}$ is zero, but the discriminant will be zero in that case.

When the coefficients of $F$ are complex numbers, the discriminant is zero if and only if either $F$ has a multiple root, which happens when $F$ and $F^{\prime}$ have a common factor, or else $F$ has degree less than $m$.

Note. The formula for the discriminant is often normalized by a factor $\pm a_{0}^{k}$. We won't make this normalization, so our formula is slightly different from the usual one.

Suppose that the coefficients $a_{i}$ of $F$ are polynomials in $t$, so that $F$ becomes a polynomial in two variables. Let $C$ be the locus $F=0$ in the affine plane $\mathbb{A}_{t, x}^{2}$. The discriminant $\operatorname{Discr}_{x}(F)$, computed regarding $x$ as the variable, will be a polynomial in $t$. At a root $t_{0}$ of the discriminant, the line $L_{0}:\left\{t=t_{0}\right\}$ is tangent to $C$, or passes though a singular point of $C$.

The discriminant of the quadratic polynomial $F(x)=a x^{2}+b x+c$ is

$$
\operatorname{det}\left(\begin{array}{ccc}
a & b & c  \tag{1.7.13}\\
2 a & b & \cdot \\
\cdot & 2 a & b
\end{array}\right)=-a\left(b^{2}-4 a c\right)
$$

discrquadr

The discriminant of the monic cubic $x^{3}+p x+q$ whose quadratic coefficient is zero is

$$
\operatorname{det}\left(\begin{array}{ccccc}
1 & \cdot & p & q & \cdot  \tag{1.7.14}\\
\cdot & 1 & \cdot & p & q \\
3 & \cdot & p & \cdot & \cdot \\
\cdot & 3 & \cdot & p & \cdot \\
\cdot & \cdot & 3 & \cdot & p
\end{array}\right)=4 p^{3}+27 q^{2}
$$

discrcubic

These are the negatives of the usual formulas. The signs are artifacts of our definition. Though it conflicts with our definition, we'll follow tradition and continue writing the discriminant of the polynomial $a x^{2}+b x+c$ as $b^{2}-4 a c$.
1.7.15. Proposition. Let $K$ be a field of characteristic zero. The discriminant of an irreducible polynomial $F$ with coefficients in $K$ isn't zero. Therefore an irreducible polynomial $F$ with coefficients in $K$ has no multiple root.
proof. When $F$ is irreducible, it cannot have a factor in common with the derivative $F^{\prime}$, which has lower degree.

This proposition is false when the characteristic of $K$ isn't zero. In characteristic $p$, the derivative $F^{\prime}$ might be the zero polynomial.
1.7.16. Proposition. Let $F=\prod\left(x-\alpha_{i}\right)$ be a polynomial that is a product of monic linear factors. Then

$$
\operatorname{Discr}(F)=\prod_{i} F^{\prime}\left(\alpha_{i}\right)=\prod_{i \neq j}\left(\alpha_{i}-\alpha_{j}\right)= \pm \prod_{i<j}\left(\alpha_{i}-\alpha_{j}\right)^{2}
$$

proof. The fact that $\operatorname{Discr}(F)=\prod F^{\prime}\left(\alpha_{i}\right)$ follows from Proposition 1.7.9. We show that $F^{\prime}\left(\alpha_{i}\right)=\prod_{j, j \neq i}\left(\alpha_{i}-\right.$ $\alpha_{j}$ ), noting that

$$
\prod_{j, j \neq i}\left(\alpha_{i}-\alpha_{j}\right)=\prod_{i}\left(\alpha_{i}-\alpha_{1}\right) \cdots\left(\widehat{\alpha_{i}-\alpha_{i}}\right) \cdots\left(\alpha_{i}-\alpha_{n}\right)
$$

where the hat ${ }^{\wedge}$ indicates that that term is deleted. By the product rule for differentiation,

$$
F^{\prime}(x)=\sum_{k}\left(x-\alpha_{1}\right) \cdots\left(\widehat{x-\alpha_{k}}\right) \cdots\left(x-\alpha_{n}\right)
$$

Substituting $x=\alpha_{i}$, all terms in the sum, except the one with $i=k$, become zero.
1.7.17. Corollary. $\operatorname{Discr}(F(x))=\operatorname{Discr}(F(x-c))$.
translate-
discr
discrprop

$$
\operatorname{Discr}(F G)= \pm \operatorname{Discr}(F) \operatorname{Discr}(G) \operatorname{Res}(F, G)^{2}
$$

dis-
crnotzero
discrformulas
1.7.18. Proposition. Let $F(x)$ and $G(x)$ be monic polynomials. Then

## firrover $K$

coverline
poly-
proof. This proposition follows from Propositions 1.7 .9 and 1.7 .16 for polynomials with complex coefficients. It is true for polynomials with coefficients in any ring because it is an identity. For the same reason, Corollary 1.7.10 remains true with coefficients in any ring.

When $f$ and $g$ are polynomials in several variables including a variable $z, \operatorname{Res}_{z}(f, g)$ and $\operatorname{Discr}_{z}(f)$ denote the resultant and the discriminant, computed regarding $f, g$ as polynomials in $z$. They will be polynomials in the other variables.

$$
\begin{equation*}
f=c_{0} z^{d}+c_{1} z^{d-1}+\cdots+c_{d} \tag{1.7.21}
\end{equation*}
$$

with $c_{i}$ homogeneous, of degree $i$ in $x, y$. Then $c_{0}=f(0,0,1)$ will be a nonzero constant that we normalize to 1 , so that $f$ becomes a monic polynomial of degree $d$ in $z$.

The fibre of $C$ over a point $\widetilde{p}=\left(x_{0}, y_{0}\right)$ of $\mathbb{P}^{1}$ is the intersection of $C$ with the line $L_{p q}$ described above. It consists of the points $\left(x_{0}, y_{0}, \alpha\right)$ such that $\alpha$ is a root of the one-variable polynomial

$$
\begin{equation*}
\tilde{f}(z)=f\left(x_{0}, y_{0}, z\right) \tag{1.7.22}
\end{equation*}
$$

We call $C$ a branched covering of $\mathbb{P}^{1}$ of degree $d$. All but finitely many fibres of $C$ over $\mathbb{P}^{1}$ consist of $d$ points (Lemma 1.7.19). The fibres with fewer than $d$ points are those above the zeros of the discriminant. They are the branch points of the covering.
(1.7.23) the genus of a plane curve

We use the discriminant to describe the topological structure of smooth plane curves in the classical topology.
1.7.24. Theorem. A smooth projective plane curve of degree $d$ is a compact, orientable and connected manifold of dimension two.

The fact that a smooth curve is a two-dimensional manifold follows from the Implicit Function Theorem. (See the discussion at $\mathbf{1 . 4 . 4}$ ).
orientability: A two-dimensional manifold is orientable if one can choose one of its two sides in a continuous, consistent way. A smooth curve $C$ is orientable because its tangent space at a point is a one-dimensional complex vector space - the affine line with the equation 1.4.11. Multiplication by $i$ orients the tangent space by defining the counterclockwise rotation. Then the right-hand rule tells us which side of $C$ is "up".
compactness: A plane projective curve is compact because it is a closed subset of the compact space $\mathbb{P}^{2}$.
The connectedness of a plane curve is a subtle fact whose proof mixes topology and algebra. Unfortunately, I don't know a proof that fits into our discussion here. It will be proved later (see Theorem 8.3.11.

The topological Euler characteristic $e$ of a compact, orientable two-dimensional manifold $M$ is the alternating sum $b^{0}-b^{1}+b^{2}$ of its Betti numbers. It can be computed using a topological triangulation, a subdivision of $M$ into topological triangles, called faces, by the formula

$$
\begin{equation*}
e=\mid \text { vertices }|-| \text { edges }|+| \text { faces } \mid \tag{1.7.25}
\end{equation*}
$$

For example, a sphere is homeomorphic to a tetrahedron, which has four vertices, six edges, and four faces. Its Euler characteristic is $4-6+4=2$. Any other topological triangulation of a sphere, including the one given by the icosahedron, yields the same Euler characteristic.

Every compact, connected, orientable two-dimensional manifold is homeomorphic to a sphere with a finite number of "handles". Its genus is the number of handles. A torus has one handle. Its genus is one. The projective line $\mathbb{P}^{1}$, which is a two-dimensional sphere, has genus zero.

## Figure

The Euler characteristic and the genus are related by the formula

$$
\begin{equation*}
e=2-2 g \tag{1.7.26}
\end{equation*}
$$

The Euler characteristic of a torus is zero, and the Euler characteristic of $\mathbb{P}^{1}$ is two.
To compute the the Euler characteristic of a smooth curve $C$ of degree $d$, we analyze a generic projection to represent $C$ as a branched covering of the projective line: $C \xrightarrow{\pi} \mathbb{P}^{1}$.

## figure

We choose generic coordinates $x, y, z$ in $\mathbb{P}^{2}$ and project form the point $q=(0,0,1)$. When the defining equation of $C$ is written as a monic polynomial in $z: \quad f=z^{d}+c_{1} z^{d-1}+\cdots+c_{d}$ where $c_{i}$ is a homogeneous polynomial of degree $i$ in the variables $x, y$, the discriminant $\operatorname{Discr}_{z}(f)$ with respect to $z$ will be a homogeneous polynomial of degree $d(d-1)=d^{2}-d$ in $x, y$.

Let $\widetilde{p}$ be the image in $\mathbb{P}^{1}$ of a point $p$ of $C$. The covering $C \xrightarrow{\pi} \mathbb{P}^{1}$ will be branched at $\widetilde{p}$ when the tangent line at $p$ is the line $L_{p q}$ through $p$ and the center of projection $q$. When $q$ is generic, such a point $p$ will not be a flex point, and then $C$ and $L_{p q}$ will have one intersection $p$ of multiplicity two, and $d-2$ intersections of multiplicity one $\sqrt[1.9]{ }$. It is intuitively plausible that the discriminant $\operatorname{Discr}_{z}(f)$ will have a simple zero at the image $\widetilde{p}$ of $p$. This will be proved below, in Proposition 1.9.13 Leet's assume that this is known. Since the discriminant has degree $d^{2}-d$, there will be $d^{2}-d$ points $\widetilde{p}$ in $\mathbb{P}^{1}$ at which the discriminant vanishes and the fibre contains $d-1$ points. They are the branch points of the covering. All other fibres consist of $d$ points.

We triangulate the sphere $\mathbb{P}^{1}$ in such a way that the branch points are among the vertices, and we use the inverse images of the vertices, edges, and faces to triangulate $C$. Then $C$ will have $d$ faces and $d$ edges lying over each face and each edge of $\mathbb{P}^{1}$, respectively. There will also be $d$ vertices of $C$ lying over a vertex of $\mathbb{P}^{1}$, except when it is one of the $d^{2}-d$ branch points. In that case the the fibre will contain only $d-1$ vertices. The Euler characteristic of $C$ is obtained by multiplying the Euler characteristic of $\mathbb{P}^{1}$ by $d$ and subtracting the number of branch points.

$$
\begin{equation*}
e(C)=d e\left(\mathbb{P}^{1}\right)-\left(d^{2}-d\right)=2 d-\left(d^{2}-d\right)=3 d-d^{2} \tag{1.7.27}
\end{equation*}
$$

eulercover
This is the Euler characteristic of any smooth curve of degree $d$, so we denote it by $e_{d}$ :

$$
\begin{equation*}
e_{d}=3 d-d^{2} \tag{1.7.28}
\end{equation*}
$$

equatione
equationg

$$
\begin{equation*}
g_{d}=\frac{1}{2}\left(d^{2}-3 d+2\right)=\binom{d-1}{2} \tag{1.7.29}
\end{equation*}
$$

Thus smooth curves of degrees $1,2,3,4,5,6, \ldots$ have genus $0,0,1,3,6,10, \ldots$, respectively. A smooth plane curve cannot have genus two.

The generic projection to $\mathbb{P}^{1}$ also computes the degree of the dual curve $C^{*}$ of a smooth curve $C$ of degree $d$. The degree of $C^{*}$ is the number of intersections of $C^{*}$ with the generic line $q^{*}$ in $\mathbb{P}^{*}$. The intersection points have the form $L^{*}$, where $q$ is a point of $L$, and $L$ is tangent to $C$ at some point $p$. As we have seen, there are $d(d-1)$ such points.
degdual
1.7.30. Corollary. Let $C$ be a plane curve of degree $d$.
(i) The degree $d^{*}$ of the dual curve $C^{*}$ is equal to the number of tangent lines at smooth points of $C$ that pass through a generic point $q$ of the plane.
(ii) If $C$ smooth, the degree $d^{*}$ of the dual curve $C^{*}$ is $d(d-1)$.

The formula $d^{*}=d(d-1)$ is incorrect when $C$ is singular. If $C$ is a smooth curve of degree $3, C^{*}$ will have degree 6 , and if $C^{*}$ were smooth its dual curve $C^{* *}$, would have degree 30 . Since $C^{* *}=C$, the dual curve is singular.

### 1.8 Nodes and Cusps

nodes
seriesf
singmult
multr

a Singular Point, with its Special Lines

To analyze a singularity at the origin $p$, we blow up the plane. The blowup is the map $W \xrightarrow{\pi} X$ from the $(x, w)$-plane $W$ to the $(x, y)$-plane $X$ defined by $\pi(x, w)=(x, x w)$. It is called a "blowup" of $X$ because the fibre over the origin in $X$ is the $w$-axis $\{x=0\}$ in $W$. The map $\pi$ is bijective at points at which $x \neq 0$, and points $(x, 0)$ of $X$ with $x \neq 0$ aren't in its image. (It might seem more appropriate to call the inverse of $\pi$ the blowup, but the inverse isn't a map.)

Suppose that the origin $p$ is a double point, a point of multiplicity 2 . Let the quadratic part of $f$ be

$$
\begin{equation*}
f_{2}=a x^{2}+b x y+c y^{2} \tag{1.8.6}
\end{equation*}
$$

We may adjust coordinates so that $c$ isn't zero, and we normalize $c$ to 1 . Writing $f(x, y)=a x^{2}+b x y+y^{2}+$ $d x^{3}+\cdots$, we make the substitution $y=x w$ and cancel $x^{2}$. This gives us a polynomial

$$
g(x, w)=f(x, x w) / x^{2}=a+b w+w^{2}+d x+\cdots
$$

in which the terms represented by $\cdots$ are divisible by $x$. Let $D$ be the locus $\{g=0\}$ in $W$. The map $\pi$ restricts to a map $D \xrightarrow{\bar{\pi}} C$. Since $\pi$ is bijective at points at which $x \neq 0$, so is $\bar{\pi}$.

Suppose first that the quadratic polynomial $y^{2}+b y+a$ has distinct roots $\alpha, \beta$, so that $a x^{2}+b x y+y^{2}=$ $(y-\alpha x)(y-\beta x)$ and $g(x, w)=(w-\alpha)(w-\beta)+d x+\cdots$. In this case, the fibre of $D$ over the origin $p$ in $X$ consists of the two points $p_{1}=(0, \alpha)$ and $p_{2}=(0, \beta)$. The partial derivative $\frac{\partial g}{\partial t}$ is nonzero at $p_{1}$ and $p_{2}$, so those are smooth points of $D$. We can solve $g(x, w)=0$ for $w$ as analytic functions of $x$ near zero, say $w=u(x)$ and $w=v(x)$ with $u(0)=\alpha$ and $v(0)=\beta$. The image of $\pi(D)$ is $C$, so $C$ has two analytic branches $y=x u(x)$ and $y=x v(x)$ through the origin with distinct tangent directions $\alpha$ and $\beta$. This singularity is called a node. A node is the simplest singularity that a curve can have.

When the discriminant $b^{2}-4 a c$ is zero, $f_{2}$ will be a square, and we will have

$$
f(x, y)=(y-\alpha x)^{2}+d x^{3}+\cdots
$$

The singularity of $C$ at the origin is called a cusp. The standard cusp is the locus $y^{2}=x^{3}$.
The blowup substitution $y=x w$ gives $g(x, w)=(w-\alpha)^{2}+d x+\cdots$. Here the fibre over $(x, y)=(0,0)$ is the point $(x, w)=(0, \alpha)$, and $g_{w}(0, \alpha)=0$. However, if $d \neq 0$, then $g_{x}(0, \alpha) \neq 0$. In this case, $D$ is smooth at $(0,0)$, and the equation of $C$ has the form $(y-\alpha x)^{2}=d x^{3}+\cdots$. All cusps are analytically equivalent with the standard cusp.

Cusps have an interesting geometry. The intersection of the standard cusp $X:\left\{y^{2}=x^{3}\right\}$ with a small 3 -sphere $S:\left\{\bar{x} x+\bar{y} y=\epsilon\right.$ in $\mathbb{C}^{2}$ is a trefoil knot.

To explain this, we parametrize $X$ by $(x(t), y(t))=\left(t^{3}, t^{2}\right)$, and we restrict to the unit circle $t=e^{i \theta}$. The locus of points of $X$ of absolute value $\sqrt{2}$ is $(x(t), y(t))=\left(e^{3 i \theta}, e^{2 i \theta}\right)$. To visualize this locus, we embed it in the product of the unit $x$-circle and the unit $y$-circle, a torus, and we distort that torus, representing it as the usual torus $T$ in $\mathbb{R}^{3}$. Let the circumference of $T$ represent the $x$-coordinate, and let the loop through the hole represent $y$. Then, as $\theta$ runs from 0 to $2 \pi,(x(t), y(t))$ goes around the circumference twice, and it loops through the hole three times, as is illustrated below.

> trefoil2.png

## figure

1.8.7. Corollary. A double point p of a curve $C$ is a node or a cusp if and only if the blowup of $C$ is smooth at the points that lie over $p$.

## nodeor-

 cuspThe simplest example of a double point that isn't a node or cusp is a tacnode, a point at which two smooth branches of a curve intersect with the same tangent direction.

1.8.8.
a Node, a Cusp, and a Tacnode (real locus)

A note about figures. In algebraic geometry, the dimensions are too big to allow realistic figures. Even with an affine plane curve, one is dealing with a locus in the space $\mathbb{A}^{2}$, whose dimension as a real vector space is four. In some cases, such as in the figures above, depicting the real locus can be helpful, but in most cases, even the real locus is too big, and one must make do with a schematic diagram. The figure below is an example of such a diagram. My students have told me that all of my figures look more or less like this:
1.8.9.


## A Typical Schematic Figure

### 1.9 Hensel's Lemma

hensel
multiplypolys
prodeqns

The resultant matrix 1.7 .4 arises in a second context that we explain here.
Suppose given a product $P=F G$ of two polynomials, say
(1.9.1) $\left(c_{0} x^{m+n}+c_{1} x^{m+n-1}+\cdots+c_{m+n}\right)=\left(a_{0} x^{m}+a_{1} x^{m-1}+\cdots+a_{m}\right)\left(b_{0} x^{n}+b_{1} x^{n-1}+\cdots+b_{n}\right)$

We call the relations among the coefficients implied by this polynomial equation the product equations. They are

$$
c_{i}=a_{i} b_{0}+a_{i-1} b_{1}+\cdots+a_{0} b_{i}
$$

for $i=0, \ldots, m+n$. For instance, when $m=3$ and $n=2$, they are
1.9.2.

$$
\begin{aligned}
& c_{0}=a_{0} b_{0} \\
& c_{1}=a_{1} b_{0}+a_{0} b_{1} \\
& c_{2}=a_{2} b_{0}+a_{1} b_{1}+a_{0} b_{2} \\
& c_{3}=a_{3} b_{0}+a_{2} b_{1}+a_{1} b_{2} \\
& c_{4}=r \\
& c_{5}=r a_{2} b_{1}+a_{2} b_{2} \\
& a_{3} b_{2}
\end{aligned}
$$

Let $J$ denote the Jacobian matrix of partial derivatives of $c_{1}, \ldots, c_{m+n}$ with respect to the variables $b_{1}, \ldots, b_{n}$ and $a_{1}, \ldots, a_{m}$, treating $a_{0}, b_{0}$ and $c_{0}$ as constants. When $m, n=3,2$,

$$
J=\frac{\partial\left(c_{i}\right)}{\partial\left(b_{j}, a_{k}\right)}=\left(\begin{array}{ccccc}
a_{0} & . & b_{0} & . & .  \tag{1.9.3}\\
a_{1} & a_{0} & b_{1} & b_{0} & \cdot \\
a_{2} & a_{1} & b_{2} & b_{1} & b_{0} \\
a_{3} & a_{2} & \cdot & b_{2} & b_{1} \\
\cdot & a_{3} & \cdot & . & b_{2}
\end{array}\right)
$$

1.9.4. Lemma. The Jacobian matrix $J$ is the transpose of the resultant matrix $\mathcal{R}$ 1.7.4.

## jacres

1.9.5. Corollary. Let $F$ and $G$ be polynomials with complex coefficients. The Jacobian matrix is singular if and only if $F$ and $G$ have a common root, or else $a_{0}=b_{0}=0$.

This corollary has an application to polynomials with analytic coefficients. Let

$$
\begin{equation*}
P(t, x)=c_{0}(t) x^{d}+c_{1}(t) x^{d-1}+\cdots+c_{d}(t) \tag{1.9.6}
\end{equation*}
$$

be a polynomial in $x$ whose coefficients $c_{i}$ are analytic functions, defined for small values of $t$, and let $\bar{P}=$ $P(0, x)=\bar{c}_{0} x^{d}+\bar{c}_{1} x^{d-1}+\cdots+\bar{c}_{d}$ be the evaluation of $P$ at $t=0$, so that $\bar{c}_{i}=c_{i}(0)$. Suppose given a factorization $\bar{P}=\bar{F} \bar{G}$, where $\bar{G}=\bar{b}_{0} x^{n}+\bar{b}_{1} x^{n-1}+\cdots+\bar{b}_{n}$ is a polynomial and $\bar{F}=x^{m}+\bar{a}_{1} x^{m-1}+\cdots+\bar{a}_{m}$ is a monic polynomial, both with complex coefficients. Are there polynomials $F(t, x)=x^{m}+a_{1} x^{m-1}+$ $\cdots+a_{m}$ and $G(t, x)=b_{0} x^{n}+b_{1} x^{n-1}+\cdots+b_{n}$, with $F$ monic, whose coefficients $a_{i}$ and $b_{i}$ are analytic functions defined for small $t$, such that $F(0, x)=\bar{F}, G(0, x)=\bar{G}$ and $P=F G$ ?
1.9.7. Hensel's Lemma. With notation as above, suppose that $\bar{F}$ and $\bar{G}$ have no common root. Then $P$ factors, as above.
proof. Since $F$ is supposed to be monic, we set $a_{0}(t)=1$. The first product equation tells us that $b_{0}(t)=c_{0}(t)$. Corollary 1.9.5 tells us that the Jacobian matrix for the remaining product equations is nonsingular at $t=0$, so according to the Implicit Function Theorem, the product equations have a unique solution in analytic functions $a_{i}(t), b_{j}(t)$ for small $t$.

Note that $P$ isn't assumed to be monic. If $\bar{c}_{0}=0$, the degree of $\bar{P}$ will be less than the degree of $P$. In that case, $\bar{G}$ will have lower degree than $G$.

## figure

1.9.8. Example. Let $P=c_{0}(t) x^{2}+c_{1}(t) x+c_{2}(t)$. The product equations for factoring $P$ as a product $F G=\left(x+a_{1}\right)\left(b_{0} x+b_{1}\right)$ of linear polynomials, with $F$ monic, are

$$
c_{0}=b_{0}, \quad c_{1}=a_{1} b_{0}+b_{1}, \quad c_{2}=a_{1} b_{1}
$$

and the Jacobian matrix is

$$
\frac{\partial\left(c_{1}, c_{2}\right)}{\partial\left(b_{1}, a_{1}\right)}=\left(\begin{array}{cc}
1 & b_{0} \\
a_{1} & b_{1}
\end{array}\right)
$$

Suppose that $\bar{P}=P(0, x)$ factors: $\bar{c}_{0} x^{2}+\bar{c}_{1} x+\bar{c}_{2}=\left(x+\bar{a}_{1}\right)\left(\bar{b}_{0} x+\bar{b}_{1}\right)=\bar{F} \bar{G}$. The determinant of the Jacobian matrix at $t=0$ is $\bar{b}_{1}-\bar{a}_{1} \bar{b}_{0}$. It is nonzero if and only if the two factors are relatively prime, in which case $P$ factors too.

On the other hand, the one-variable Jacobian criterion allows us to solve the equation $P(t, x)=0$ for $x$ as function of $t$ with $x(0)=-\bar{a}_{1}$, provided that $\frac{\partial P}{\partial x}=2 c_{0} x+c_{1}$ isn't zero at the point $(t, x)=\left(0,-\bar{a}_{1}\right)$. Evaluating, $\frac{\partial P}{\partial x}\left(0,-\bar{a}_{1}\right)=-2 \bar{c}_{0} \bar{a}_{1}+\bar{c}_{1}$. Substituting $\bar{c}_{0}=\bar{b}_{0}$ and $\bar{c}_{1}=\bar{a}_{1} \bar{b}_{0}+\bar{b}_{1}$, shows that $\frac{\partial P}{\partial x}\left(0,-\bar{a}_{1}\right)=$ $-2 \bar{c}_{0} \bar{a}_{1}+\bar{c}_{1}=\bar{b}_{1}-\bar{a}_{1} \bar{b}_{0}$. Not surprisingly, the two conditions for factoring are the same.

## (1.9.9) general position

In algebraic geometry, the phrases general position and generic indicate an object, such as a point, has no special 'bad' properties. Typically, the object will be parametrized somehow, and the word generic indicates

## generic-

 condfinlines vandisc
discrimvanishing
that the parameter representing that particular object avoids a proper closed subset of the parameter space that may be described explicitly or not. Proposition 1.9 .13 below refers to a generic point $q$. To be precise about the requirement in this case, $q$ shall not lie on any of these lines:
flex tangent lines and bitangent lines,
lines that contain more than one singular point,
special lines through singular points (see $\mathbf{1 . 8 . 2}$ ),
tangent lines that contain a singular point of $C$.
1.9.11. Lemma. This is a list of finitely many lines that $q$ must avoid.
beginning of the proof. Proposition 1.4.17 shows that there are finitely many flex tangents. Since there are finitely many singular points, there are finitely many special lines and finitely many lines through pairs of singular points. To show that there are finitely many tangent lines that pass through singular points, we project $C$ from a singular point $p$ and apply Lemma 1.7.19. The discriminant isn't identically zero, so it vanishes finitely often. The proof that there are finitely many bitangents will be given later, in Corollary 1.10.15.

## (1.9.12) order of vanishing of the discriminant

Let $f(x, y, z)$ be a homogeneous polynomial with no multiple factors, and let $C$ be the (possibly reducible) plane curve $\{f=0\}$. Suppose that $q=(0,0,1)$ is in general position, in the sense described above. Let $L_{p q}$ be the line through a point $p=s\left(x_{0}, y_{0}, 0\right)$ and $q$, the set of points $\left(x_{0}, y_{0}, z_{0}\right)$, n , as before.
1.9.13. Proposition. (i) If $p$ is a smooth point of $C$ with tangent line $L_{p q}$, the discriminant $\operatorname{Discr}_{z}(f)$ has a simple zero at $\widetilde{p}$.
(ii) If $p$ is a node of $C, \operatorname{Discr}_{z}(f)$ has a double zero at $\widetilde{p}$.
(iii) If $p$ is a cusp, $\operatorname{Discr}_{z}(f)$ has a triple zero at $\widetilde{p}$.
(iv) If $p$ is a an ordinary flex point of $C(\sqrt{\mathbf{1 . 4 . 8}})$ with tangent line $L_{p q}$, $\operatorname{Discr}_{z}(f)$ has a double zero at $\widetilde{p}$.
proof. (i)-(iii) There are several ways to proceed, none especially simple. We'll use Hensel's Lemma. We set $x=1$, to work in the standard affine open set $\mathbb{U}$ with coordinates $y, z$. In affine coordinates, the projection $\pi$ is the map $(y, z) \rightarrow y$. We may suppose that $p$ is the origin in $\mathbb{U}$. Its image $\widetilde{p}$ will be the point $y=0$ of the affine $y$-line, and the intersection of the line $L_{p q}$ with $\mathbb{U}$ will be the line $\widetilde{L}:\{y=0\}$. We'll denote the defining polynomial of the curve $C$, restricted to $\mathbb{U}$, by $f(y, z)$ instead of $f(1, y, z)$. Let $\widetilde{f}(z)=f(0, z)$.

In each of the cases (i),(ii),(iii), the polynomial $\bar{f}(z)=f(0, z)$ will have a double zero at $z=0$, so we will have $\bar{f}(z)=z^{2} \bar{h}(z)$, with $\bar{h}(0) \neq 0$. Then $z^{2}$ and $\bar{h}(z)$ have no common root, so we may apply Hensel's Lemma to write $f(y, z)=g(y, z) h(y, z)$, where $g$ and $h$ are polynomials in $z$ whose coefficients are analytic functions of $y$, defined for small $y, g$ is monic, $g(0, z)=z^{2}$, and $h(0, z)=\bar{h}$. Then

$$
\begin{equation*}
\operatorname{Discr}_{z}(f)= \pm \operatorname{Discr}_{z}(g) \operatorname{Discr}_{z}(h) \operatorname{Res}_{z}(g, h)^{2} \tag{1.9.14}
\end{equation*}
$$

1.7.18). Since $q$ is in general position, $\bar{h}$ will have simple zeros. Then $\operatorname{Discr}_{z}(h)$ doesn't vanish at $y=0$. Neither does $\operatorname{Res}_{z}(g, h)$. So the orders of vanishing of $\operatorname{Discr}_{z}(f)$ and $\operatorname{Discr}_{z}(g)$ are equal. We replace $f$ by $g$.

Since $g$ is a monic quadratic polynomial, it will have the form

$$
g(y, z)=z^{2}+b(y) z+c(y)
$$

The coefficients $b$ and $c$ are analytic functions of $y$, and $g(0, z)=z^{2}$. The discriminant $\operatorname{Discr}_{z}(g)=b^{2}-4 c$ is unchanged when we complete the square by the substitution of $z-\frac{1}{2} b$ for $z$, and if $\widetilde{p}$ is a node or a cusp, that property isn't affected by this change of coordinates. So we may assume that $g$ has the form $z^{2}+c(y)$. The discriminant is $D=4 c(y)$.

We write $c(y)$ as a series in $y$ :

$$
c(y)=c_{0}+c_{1} y+c_{2} y^{2}+c_{3} y^{3}+\cdots
$$

The constant coefficient $c_{0}$ is zero because $\widetilde{p}$ is a point of $C$. If $c_{1} \neq 0, \widetilde{p}$ is a smooth point with tangent line $\widetilde{L}:\{y=0\}$, and $D$ has a simple zero. If $\widetilde{p}$ is a node, $c_{0}=c_{1}=0$ and $c_{2} \neq 0$. Then $D$ has a double zero. If $\widetilde{p}$ is a cusp, $c_{0}=c_{1}=c_{2}=0$, and $c_{3} \neq 0$. Then $D$ has a triple zero at $\widetilde{p}$.
(iv) In this case, the polynomial $\tilde{f}(z)=f(0, z)$ will have a triple zero at $z=0$. Proceding as above, we may factor: $f=g h$ where $g$ and $h$ are polynomials in $z$ with analyic coefficients in $y$, and $g(y, z)=z^{3}+a(y) z^{2}+$ $b(y) z+c(y)$. We eliminate the quadratic coefficient $a$ by substituting $z-\frac{1}{3 a}$ for $z$. With $g=z^{3}+b z+c$ in the new coordinates, the discriminant $\operatorname{Discr}_{z}(g)$ is $4 b^{3}+27 c^{2}$ 1.7.14. We write $c(y)=c_{0}+c_{1} y+\cdots$ and $b(y)=b_{0}+b_{1} y+\cdots$. Since $p$ is a point of $C$ with tangent line $\{y=0\}, c_{0}=0$ and $c_{1} \neq 0$. Since the intersection multiplicity of $C$ with the line $\{y=0\}$ at $\widetilde{p}$ is three, $b_{0}=0$. The discriminant has a zero of order two.
1.9.15. Corollary. Let $C:\{g=0\}$ and $D:\{h=0\}$ be plane curves that intersect transversally at a point $p=\left(x_{0}, y_{0}, z_{0}\right)$. With coordinates in general position, $\operatorname{Res}_{z}(g, h)$ has a simple zero at $\left(x_{0}, y_{0}\right)$.

Two curves are said to intersect transversally at a point $p$ if they are smooth at $p$ and their tangent lines there are distinct.
proof. Proposition 1.9 .13 (ii) applies to the product $g h$, whose zero locus is the union $C \cup D$. It shows that the discriminant $\operatorname{Discr}_{z}(g h)$ has a double zero at $\tilde{p}$. We also have the formula 1.9.14 with $f=g h$. When coordinates are in general position, $\operatorname{Discr}_{z}(g)$ and $\operatorname{Discr}_{z}(h)$ will not be zero at $\widetilde{p}$. Then $\operatorname{Res}_{z}(g, h)$ has a simple zero there.

### 1.10 Bézout's Theorem

Bézout's Theorem counts intersections of plane curves. We state it here in a form that is ambiguous because it contains a term "multiplicity" that hasn't yet been defined.
1.10.1. Bézout's Theorem. Let $C$ and $D$ be distinct curves of degrees $m$ and $n$, respectively. When intersections are counted with the appropriate multiplicity, the number of intersections is equal to $m n$. Moreover, the multiplicity at a point is 1 at a transversal intersection.

As before, $C$ and $D$ intersect transversally at $p$ if they are smooth at $p$ and their tangent lines there are distinct.
1.10.2. Corollary. Bézout's Theorem is true when one of the curves is a line.

See Corollary 1.3.9. The multiplicity of intersection of a curve and a line is the one that was defined there.
The proof in the general case requires some algebra that we would rather defer. It will be given later (Theorem7.8.1). It is possible to determine the intersections by counting the zeros of the resultant with respect to one of the variables. To do this, one chooses generic coordinates $x, y, z$, Then neither $C$ nor $D$ contains the point $(0,0,1)$. One writes their defining polynomials $f$ and $g$ as polynomials in $z$ with coefficients in $\mathbb{C}[x, y]$. The resultant $R$ with respect to $z$ will be a homogeneous polynomial in $x, y$, of degree $m n$. It will have $m n$ zeros in $\mathbb{P}_{x, y}^{1}$, counted with multiplicity. If $\widetilde{p}=\left(x_{0}, y_{0}\right)$ is a zero of $R, f\left(x_{0}, y_{0}, z\right)$ and $g\left(x_{0}, y_{0}, z\right)$, which are polynomials in $z$, have a common root $z=z_{0}$, and then $p=\left(x_{0}, y_{0}, z_{0}\right)$ will be a point of $C \cap D$. It is a fact that the multiplicity of the zero of the resultant $R$ at the image $\widetilde{p}$ is the (as yet undefined) intersection multiplicity of $C$ and $D$ at $p$. Unfortunately, this won't be obvious, even when multiplicity is defined. However, one can prove the next proposition using this approach.
1.10.3. Proposition. Let $C$ and $D$ be distinct plane curves of degrees $m$ and $n$, respectively.
(i) The curves $C$ and $D$ have at least one point of intersection, and the number of intersections is at most $m n$.
(ii) If all intersections are transversal, the number of intersections is precisely mn.

It isn't obvious that two curves in the projective plane intersect. If two curves in the affine plane have no intersection, if they are parallel lines, for instance, their closures in the projective plane meet on the line at infinity.
1.10.4. Lemma. Let $f$ and $g$ be homogeneous polynomials in $x, y, z$ of degrees $m$ and $n$, respectively, and suppose that the point $(0,0,1)$ isn't a zero of $f$ or $g$. If the resultant $\operatorname{Res}_{z}(f, g)$ with respect to $z$ is identically zero, then $f$ and $g$ have a common factor.
proof. Let the degrees of $f$ and $g$ be $m$ and $n$, respectively, and let $F$ denote the field of rational functions $\mathbb{C}(x, y)$. If the resultant is zero, $f$ and $g$ have a common factor in $F[z]$ (Corollary 1.7.6. There will be polynomials $p$ and $q$ in $F[z]$, of degrees at most $n-1$ and $m-1$ in $z$, respectively, such that $p f=q g$ 1.7.2 . We may clear denominators, so we may assume that the coefficients of $p$ and $q$ are in $\mathbb{C}[x, y]$. Then $p f=q g$ is an equation in $\mathbb{C}[x, y, z]$. Since $p$ has degree at most $n-1$ in $z$, it isn't divisible by $g$, which has degree $n$ in $z$. Since $\mathbb{C}[x, y, z]$ is a unique factorization domain, $f$ and $g$ have a common factor.
proof of Proposition 1.10 .3 (i) Let $f$ and $g$ be irreducible polynomials whose zero sets $C$ and $D$, are distinct. Proposition 1.3 .11 shows that there are finitely many intersections. We project to $\mathbb{P}^{1}$ from a point $q$ that doesn't lie on any of the finitely many lines through pairs of intersection points. Then a line through $q$ passes through at most one intersection, and the zeros of the resultant $\operatorname{Res}_{z}(f, g)$ that correspond to the intersection points will be distinct. Since the resultant has degree $m n(1.7 .8)$, it has at least one zero, and at most $m n$ of them. Therefore $C$ and $D$ have at least one and at most $m n$ intersections.
(ii) Every zero of the resultant will be the image of an intersection of $C$ and $D$. To show that there are $m n$ intersections if all intersections are transversal, it suffices to show that the resultant has simple zeros. This is Corollary 1.9.15
1.10.5. Corollary. If the curve $X$ defined by a homogeneous polynomial $f(x, y, z)$ is smooth, then $f$ is irreducible, and therefore $X$ is a smooth curve.
proof. Suppose that $f=g h$, and let $p$ be a point of intersection of the loci $\{g=0\}$ and $\{h=0\}$. The previous proposition shows that such a point exists. All partial derivatives of $f$ vanish at $p$, so $p$ is a singular point of $X$.
1.10.6. Corollary. (i) Let $d$ be an integer $\geq 3$. A smooth plane curve of degree $d$ has at least one flex point, and the number of flex points is at most $3 d(d-2)$.
(ii) If all flex points are ordinary, the number of flex points is equal to $3 d(d-2)$.

Thus smooth curves of degrees $2,3,4,5, \ldots$ have at most $0,9,24,45, \ldots$ flex points, respectively. proof. (i) The flex points are intersections of a smooth curve $C$ with its Hessian divisor $D:\{\operatorname{det} H=0\}$. (We use the definition of divisor that is given in 1.3.12, ) Let $C:\left\{f\left(x_{0}, x_{1}, x_{2}\right)=0\right\}$ be a smooth curve of degree $d$. The entries of the $3 \times 3$ Hessian matrix $H$ are the second partial derivatives $\frac{\partial^{2} f}{\partial x_{i} \partial x_{j}}$. They are homogeneous polynomials of degree $d-2$, so the Hessian determinant is homogeneous, of degree $3(d-2)$. Propositions 1.4.17 and 1.10.3 tell us that there are at most $3 d(d-2)$ intersections.
(ii) Recall that a flex point is ordinary if the multiplicity of intersection of the curve and its tangent line is 3 . Bézout's Theorem asserts that the number of flex points is equal to $3 d(d-2)$ if the intersections of $C$ with its Hessian divisor $D$ are transversal, and therefore have multiplicity 1. So the next lemma completes the proof.
1.10.7. Lemma. A curve $C:\{f=0\}$ intersects its Hessian divisor $D$ transversally at a point $p$ if and only $p$ is an ordinary flex point of $C$.
proof. We prove this by computation. There may be a conceptual proof, but I don't know one.
Let $L$ be the tangent line to $C$ at the flex point $p$, and let $h$ denote the restriction of the Hessian determinant to $L$. The Hessian divisor $D$ will be transversal to $C$ at $p$ if and only if it is transversal to $L$, and this will be true if and only if the order of vanishing of $h$ at $p$ is 1 .

We adjust coordinates $x, y, z$ so that $p=(0,0,1)$ and $L$ is the line $\{y=0\}$, and we write the polynomial $f$ of degree $d$ as

$$
\begin{equation*}
f(x, y, z)=\sum_{i+j+k=d} a_{i j} x^{i} y^{j} z^{k}, \tag{1.10.8}
\end{equation*}
$$

We set $y=0$ and $z=1$, to restrict $f$ to $L$. The restriction of the polynomial $f$ is

$$
f(x, 0,1)=\sum_{i \leq d} a_{i 0} x^{i}
$$

Since $p$ is a flex point with tangent line $L$, the coefficients $a_{00}, a_{10}$, and $a_{20}$ are zero, and $p$ is an ordinary flex point if and only if the coefficient $a_{30}$ is nonzero.

Let $h$ be the restiction of $\operatorname{det} H$ to $L: h=\operatorname{det} H(x, 0,1)$. We must show that $p$ is an ordinary flex point if and only if $h$ has a simple zero at $x=0$.

To evaluate the restriction $f_{x x}(x, 0,1)$ of the partial derivative to $L$, the relevant terms in the sum 1.10.8) have $j=0$. Since $a_{00}=a_{10}=0$,

$$
f_{x x}(x, 0,1)=6 a_{30}+12 a_{40} x^{2}+\cdots=6 a_{30} x+O(2)
$$

Similarly,

$$
\begin{aligned}
& f_{x z}(x, 0,1)=0+O(2) \\
& f_{z z}(x, 0,1)=0+O(2)
\end{aligned}
$$

For the restriction of $f_{y z}$, the relevant terms are those with $j=1$ :

$$
f_{y z}(x, 0,1)=(d-1) a_{01}+(d-2) a_{11} x+O(2)
$$

We don't need $f_{x y}$ or $f_{y y}$.
Let $v=6 a_{30} x$ and $w=(d-1) a_{01}+(d-2) a_{11} x$. The restricted Hessian matrix has the form

$$
H(x, 0,1)=\left(\begin{array}{ccc}
v & * & 0  \tag{1.10.9}\\
* & * & w \\
0 & w & 0
\end{array}\right)+O(2)
$$

where $*$ are entries that don't affect terms of degree at most one in the determinant. The determinant is

$$
h=-v w^{2}+O(2)=-6(d-1)^{2} a_{30} a_{01}^{2} x+O(2)
$$

It has a zero of order 1 at $x=0$ if and only if $a_{30}$ and $a_{01}$ aren't zero. Since $C$ is smooth at $p$ and $a_{10}=0$, the coefficient $a_{01}$ isn't zero. Thus the curve $C$ and its Hessian divisor $D$ intersect transversally, and $C$ and $L$ intersect with multiplicity 3 , if and only if $a_{30}$ is nonzero, which is true if and only if $p$ is an ordinary flex.
1.10.10. Corollary. A smooth cubic curve contains exactly 9 flex points.
proof. Let $f$ be the irreducible cubic polynomial whose zero locus is a smooth cubic $C$. The degree of the Hessian divisor $D$ is also 3 , so Bézout predicts at most 9 intersections of $D$ with $C$. To derive the corollary, we show that $C$ intersects $D$ transversally. According to Proposition 1.10.7, a nontransversal intersection would correspond to a point at which the curve and its tangent line intersect with multiplicity greater than 3 . This is impossible when the curve is a cubic.

## (1.10.11) singularities of the dual curve

Let $C$ be a plane curve. As before, an ordinary flex point is a smooth point $p$ such that the intersection multiplicity of the curve and its tangent line $L$ at $p$ is precisely 3 . A bitangent to $C$ is a line $L$ that is tangent to $C$ at distinct smooth points $p$ and $q$, and an ordinary bitangent is one such that neither $p$ nor $q$ is a flex point. A tangent line $L$ at a smooth point $p$ of $C$ is an ordinary tangent if it isn't a flex point or a bitangent.

The line $L$ will have other intersections with $C$. Most often, these other intersections will be transversal. However, it may happen that $L$ is tangent to $C$ at such a point, or that it is a singular point of $C$. Let's call such occurences accidents.
1.10.12. Proposition. Let $p$ be a smooth point of a curve $C$, and let $L$ be the tangent line at $p$. Suppose that
there are no accidents.
(i) If $L$ is an ordinary tangent at $p$, then $L^{*}$ is a smooth point of $C^{*}$.
(ii) If $L$ is an ordinary bitangent, then $L^{*}$ is a node of $C^{*}$.
(iii) If $p$ is an ordinary flex point, then $L^{*}$ is a cusp of $C^{*}$.
proof. We refer to the map $U \xrightarrow{t} C^{*}$ 1.6.3 from the set of smooth points of $C$ to the dual curve. We set $z=1$ and choose affine coordinates so that $p$ is the origin, and the tangent line $L$ at $p$ is the line $\{y=0\}$.

### 1.10.15. Corollary. A plane curve has finitely many bitangents.

This corollary is true whether or not the bitangents are ordinary. It follows from the fact that the dual curve $C^{*}$ has finitely many singular points 1.4.7. If $L$ is a bitangent, ordinary or not, $L^{*}$ will be a singular point of $C^{*}$.

### 1.11 The Plücker Formulas

plucker A plane curve $C$ is ordinary if it is smooth, all of its bitangents and flex points are ordinary (see (1.10.11), and there are no accidents. The Plücker formulas compute the number of flexes and bitangents of an ordinary plane curve.

For the next proposition, we refer back to the notation of Section 1.7.20. With coordinates in general position, let $\pi: C \rightarrow X$ be the projection of a plane curve $C$ to the projective line $X$ from $q=(0,0,1)$. If $\widetilde{p}=\left(x_{0}, y_{0}\right)$ is a point of $X$, we denote by $L_{\widetilde{p}}$ the line in $\mathbb{P}^{2}$ such that the fibre of $\pi$ over $\widetilde{p}$ is the complement of $q$ in $L_{\widetilde{p}}$.

The covering $\pi$ will be branched at the points $\widetilde{p}=\left(x_{0}, y_{0}\right)$ of $X$ such that $L_{\widetilde{p}}$ tangent line to $C$ at some point. It will also be branched the images of singular points of $C$.
1.11.1. Theorem: Plücker Formulas. Let $C$ be an ordinary curve of degree d at least two, and let $C^{*}$ be its dual curve. Let $f$ and $b$ denote the numbers of flex points and bitangents of $C$, and let $f^{*}, \delta^{*}$ and $\kappa^{*}$ denote the degree, the numbers of nodes, and the number of cusps of $C^{*}$, respectively. Then:
(i) The dual curve $C^{*}$ has no flexes or bitangents. Its singularities are nodes and cusps.
(ii) $d^{*}=d^{2}-2, \quad f=\kappa^{*}=3 d(d-2)$, and $\quad b=\delta^{*}=\frac{1}{2} d(d-2)\left(d^{2}-9\right)$.
proof. (i) A bitangent or a flex on $C^{*}$ would produce a singularity on the bidual $C^{* *}$, which is the smooth curve $C$.
(ii) The degree $d^{*}$ was computed in Corollary 1.7.30 Bézout's Theorem counts the flex points (see (1.10.6)). The facts that $\kappa^{*}=f$ and $\delta^{*}=b$ are dealt with in Proposition 1.10.12. Thus $\kappa^{*}=f=3 d(d-2)$.

We project $C^{*}$ to $\mathbb{P}^{1}$ from a generic point $s$ of $\mathbb{P}^{*}$. The number of branch points that correspond to tangent lines through $s$ at smooth points of $C^{*}$ is the degree of $C^{* *}=C$ 1.7.30, which is $d$.

Next, let $F$ be the defining polynomial for $C^{*}$. The discriminant $\operatorname{Discr}_{z}(F)$ has degree $d^{* 2}-d^{*}$. Proposition 1.9.13 describes the order of vanishing of the discriminant at the images of the $d$ tangent lines through $s$, the $\delta$ nodes of $C^{*}$, and the $\kappa$ cusps of $C^{*}$. It tells us that

$$
d^{* 2}-d^{*}=d+2 \delta^{*}+3 \kappa^{*}
$$

Substituting the known values $d^{*}=d^{2}-d$, and $\kappa^{*}=3 d(d-2)$ into this formula gives us

$$
\left(d^{2}-d\right)^{2}-\left(d^{2}-d\right)=d+2 \delta^{*}+9 d(d-2) \quad \text { or } \quad 2 \delta^{*}=d^{4}-2 d^{3}-9 d^{2}+18 d
$$

Note. It isn't easy to count the number of bitangents directly.

### 1.11.2. Examples.

(i) All curves of degree 2 and all smooth curves of degree 3 are ordinary.
(ii) A curve of degree 2 has no flexes and no bitangents. Its dual curve has degree 2 .
(iii) A smooth curve of degree 3 has 9 flexes and no bitangents. Its dual curve has degree 6 .
(iv) An ordinary curve $C$ of degree 4 has 24 flexes and 28 bitangents. Its dual curve has degree 12 .

We will make use of the fact that a quartic curve has 28 bitangents in Chapter 4 (see 4.8.14). The Plücker Formulas are rarely used for curves of degree greater than four.

### 1.11.3. Example.

We describe the dual of a plane cubic curve $C$ with a cusp again.
Projecting generically to $X=\mathbb{P}^{1}, C$ becomes a triple cover of $X$. The discriminant has degree 6 , and it has a triple zero at the image of the cusp (1.9.13), and it will also have three simple zeros. The degree of $C^{*}$ is three.

## Chapter 2 AFFINE ALGEBRAIC GEOMETRY

affine
july 17
1 Rings and Modules
2.2 The Zariski Topology
2.3 Some Affine Varieties
2.4 The Nullstellensatz
2.5 The Spectrum
2.6 Localization
2.7 Morphisms of Affine Varieties
2.8 Finite Group Actions

In the next chapters, we study varieties of arbitrary dimension. We will use some of the basic terminology that was introduced in Chapter 1 including the concepts of discriminant and transcendence degree, but most of the results of Chapter 1 won't be used again until Chapter 8

To begin, we review some basic facts about rings and modules, omitting proofs. Please look up information on the concepts that aren't familiar, as needed.

### 2.1 Rings and Modules

By the word 'ring', we mean 'commutative ring', $a b=b a$, unless the contrary is stated explicitly. A domain is a ring that has no zero divisors and isn't the zero ring.

An algebra is a ring that contains the field $\mathbb{C}$ of complex numbers as a subring. A set of elements $\alpha=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ generates an algebra $A$ if every element of $A$ can be expressed (usually not uniquely) as a polynomial in $\alpha_{1}, \ldots, \alpha_{n}$, with complex coefficients. Another way to state this is that $\alpha$ generates $A$ if the homomorphism $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] \xrightarrow{\tau} A$ that evaluates a polynomial at $x=\alpha$ is surjective. If $\alpha$ generates $A$, then $A$ will be isomorphic to the quotient $\mathbb{C}[x] / I$ of the polynomial algebra $\mathbb{C}[x]$, where $I$ is the kernel of $\tau$. A finite-type algebra is one that can be generated by a finite set of elements.

If $I$ and $J$ are ideals of a ring $R$, the product ideal, which is denoted by $I J$, is the ideal whose elements are finite sums of products $\sum a_{i} b_{i}$, with $a_{i} \in I$ and $b_{i} \in J$. (The product ideal is not the product set, whose elements are the products $a b$, with $a \in I$ and $b \in J$.) The power $I^{k}$ of $I$ is the product of $k$ copies of $I$, the ideal spanned by products of $k$ elements of $I$. The intersection $I \cap J$ is also an ideal, and

$$
\begin{equation*}
(I \cap J)^{2} \subset I J \subset I \cap J \tag{2.1.1}
\end{equation*}
$$

An ideal $M$ of a ring $R$ is maximal if it isn't the unit ideal $R$, and there is no ideal $I$ such that $M<I<R$. This is true if and only if the quotient ring $R / M$ is a field. An ideal $P$ of a ring $R$ is a prime ideal if the quotient $R / P$ is a domain. A maximal ideal is a prime ideal.
2.1.2. Lemma. Let $P$ be an ideal of a ring $R$, not the unit ideal. The following conditions are equivalent.
(i) $P$ is a prime ideal.
(ii) If $a$ and $b$ are elements of $R$ and if $a b \in P$, then $a \in P$ or $b \in P$.
(iii) If $A$ and $B$ are ideals of $R$, and if the product ideal $A B$ is contained in $P$, then $A \subset P$ or $B \subset P$.

It is sometimes convenient to state (iii) this way:
(iii') If $A$ and $B$ are ideals that contain $P$, and if the product ideal $A B$ is contained in $P$, then $A=P$ or $B=P$.
2.1.3. Mapping Property of Quotient Rings. Let $R$ and $S$ be rings, let $K$ be an ideal of $R$, and let $R \xrightarrow{\tau} \bar{R}$ denote the canonical map from $R$ to the quotient ring $\bar{R}=R / K$. Homomorphisms $\bar{R} \xrightarrow{\bar{\varphi}} S$ correspond bijectively to homomorphisms $R \xrightarrow{\varphi} S$ whose kernels contain $K$, the correspondence being $\varphi=\bar{\varphi} \circ \tau$ :


If $\operatorname{ker} \varphi=I$, then $\operatorname{ker} \bar{\varphi}=I / K$.

## (2.1.4) commutative diagrams

In the diagram displayed above, the maps $\bar{\varphi} \tau$ and $\varphi$ from $R$ to $S$ are equal. This is referred to by saying that the diagram is commutative. A commutative diagram is one in which every map that can be obtained by composing its arrows depends only on the domain and range of the map. In these notes, all diagrams of maps are commutative. We won't mention commutativity most of the time.

### 2.1.5. Correspondence Theorem.

(i) Let $R \stackrel{\varphi}{\longrightarrow} S$ be a surjective ring homomorphism with kernel $K$. For instance, $\varphi$ might be the canonical map from $R$ to the quotient ring $R / K$. (In any case, $S$ will be isomoprhic to $R / K$.) There is a bijective correspondence

$$
\{\text { ideals of } R \text { that contain } K\} \quad \longleftrightarrow\{\text { ideals of } S\}
$$

This correspondence associates an ideal $I$ of $R$ that contains $K$ with its image $\varphi(I)$ in $S$ and it associates an ideal $J$ of $S$ with its inverse image $\varphi^{-1}(J)$ in $R$.

If an ideal $I$ of $R$ that contains $K$ corresponds to an ideal $J$ of $S$, then $\varphi$ induces an isomorphism of quotient rings $R / I \rightarrow S / J$. If one of the ideals, I or $J$, is prime or maximal, they both are.
(ii) Let $R$ be a ring, and let $M \xrightarrow{\varphi} N$ be a surjective homomorphism of $R$-modules with kernel $L$. There is a bijective correspondence

$$
\{\text { submodules of } M \text { that contain } L\} \longleftrightarrow\{\text { submodules of } N\}
$$

This correspondence associates a submodule $S$ of $M$ that contains $L$ with its image $\varphi(S)$ in $N$ and it associates a submodule $T$ of $N$ with its inverse image $\varphi^{-1}(T)$ in $M$.

Ideals $I_{1}, \ldots, I_{k}$ of a ring $R$ are said to be comaximal if the sum of any two of them is the unit ideal.
2.1.6. Chinese Remainder Theorem. Let $I_{1}, \ldots, I_{k}$ be comaximal ideals of a ring $R$.
(i) The product ideal $I_{1} \cdots I_{k}$ is equal to the intersection $I_{1} \cap \cdots \cap I_{k}$.
(ii) The map $R \longrightarrow R / I_{1} \times \cdots \times R / I_{k}$ that sends an element a of $R$ to its vector of residues is a surjective homomorphism whose kernel is $I_{1} \cap \cdots \cap I_{k} \quad\left(=I_{1} \cdots I_{k}\right)$.
(iii) Let $M$ be an $R$-module. The canonical homomorphism $M \rightarrow\left(M / I_{1} M\right) \times \cdots \times\left(M / I_{k} M\right)$ is surjective.
2.1.7. Proposition. Let $R$ be a product of rings, $R=R_{1} \times \cdots \times R_{k}$, let $I$ be an ideal of $R$, and let $\bar{R}=R / I$ be the quotient ring. There are ideals $I_{j}$ of $R_{j}$ such that $I=I_{1} \times \cdots \times I_{k}$ and $\bar{R}=R_{1} / I_{1} \times \cdots \times R_{k} / I_{k}$.

## (2.1.8) Noetherian rings

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Let $M$ and $N$ be modules over a ring $R$. By an $R$-linear map $M \rightarrow N$ we simply mean a homomorphism of $R$-modules. When we refer to a map as being linear without mentioning a ring, we mean a $\mathbb{C}$-linear map.

A finite module $M$ over a ring $R$ is a module that is spanned, or generated, by a finite set $\left\{m_{1}, \ldots, m_{k}\right\}$ of elements. To say that the set generates means that every element of $M$ can be obtained as a combination $r_{1} m_{1}+\cdots+r_{k} m_{k}$ with coefficients $r_{i}$ in $R$, or that the homomorphism from the free $R$-module $R^{k}$ to $M$ that sends a vector $\left(r_{1}, \ldots, r_{k}\right)$ to the combination $r_{1} m_{1}+\cdots+r_{k} m_{k}$ is surjective.

An ideal of a ring $R$ is finitely generated if, when regarded as an $R$-module, it is a finite module.
A ring $R$ is noetherian if all of its ideals are finitely generated. The ring $\mathbb{Z}$ of integers is noetherian. Fields are notherian. If $I$ is an ideal of a noetherian ring $R$, the quotient ring $R / I$ is noetherian.
2.1.9. Hilbert Basis Theorem. Let $R$ be a noetherian ring. The ring $R\left[x_{1}, \ldots, x_{n}\right]$ of polynomials with coefficients in $R$ is noetherian.

Thus $\mathbb{Z}\left[x_{1}, \ldots, x_{n}\right]$ and $F\left[x_{a}, \ldots, x_{n}\right], F$ a field, are noetherian rings.

### 2.1.10. Corollary. Every finite-type algebra is noetherian.

Note. It is important not to confuse the concept of a finite-type algebra with that of a finite module. A finite $R$-module $M$ is a module in which every element can be written as a (linear) combination $r_{1} m_{1}+\cdots+r_{k} m_{k}$ of some finite set $\left\{m_{1}, \ldots, m_{k}\right\}$ of elements of $M$, with coefficients in $R$. A finite-type algebra $A$ is an algebra which contains a finite set of elements $\left\{\alpha_{1}, \ldots, \alpha_{k}\right\}$, such that every element can be written as a polynomial $f\left(\alpha_{1}, \ldots, \alpha_{k}\right)$, with complex coefficients.

## (2.1.11) the ascending chain condition

The condition that a ring $R$ be noetherian can be rewritten in several ways that we review here.
Our convention is that if $X^{\prime}$ and $X$ are sets, the notation $X^{\prime} \subset X$ means that $X^{\prime}$ is a subset of $X$, while $X^{\prime}<X$ means that $X^{\prime}$ is a subset that is different from $X$. A proper subset $X^{\prime}$ of a set $X$ is a nonempty subset different from $X$ - a set such that $\emptyset<X^{\prime}<X$.

A sequence $X_{1}, X_{2}, \ldots$, finite or infinite, of subsets of a set $Z$ forms an increasing chain if $X_{n} \subset X_{n+1}$ for all $n$, equality $X_{n}=X_{n+1}$ being permitted. If $X_{n}<X_{n+1}$ for all $n$, the chain is strictly increasing.

Let $\mathcal{S}$ be a set whose elements are subsets of a set $Z$. A member $M$ of $\mathcal{S}$ is a maximal member if there is no member $M^{\prime}$ of $\mathcal{S}$ such that $M<M^{\prime}$. For example, the set of proper subsets of a set of five elements contains five maximal members, the subsets of order four. The set of finite subsets of the set of integers contains no maximal member.

A maximal ideal of a ring $R$ is a maximal member of the set of ideals of $R$ different from the unit ideal.
2.1.12. Proposition. The following conditions on a ring $R$ are equivalent:
(i) $R$ is noetherian: Every ideal of $R$ is finitely generated.
(ii) The ascending chain condition: Every strictly increasing chain $I_{1}<I_{2}<\cdots$ of ideals of $R$ is finite.
(iii) Every nonempty set of ideals of $R$ contains a maximal member.

The next corollary follows from the ascending chain condition, but the conclusions are true whether or not $R$ is noetherian.
2.1.13. Corollary. Let $R$ be a noetherian ring.
(i) If $R$ isn't the zero ring, every ideal of $R$ except the unit ideal is contained in a maximal ideal.
(ii) A nonzero ring $R$ contains at least one maximal ideal.
(iii) An element of a ring $R$ that isn't in any maximal ideal is a unit - an invertible element of $R$.
2.1.14. Corollary. Let $s_{1}, \ldots, s_{k}$ be elements that generate the unit ideal of a noetherian ring $R$. For any positive integer $n$, the powers $s_{1}^{n}, \ldots, s_{k}^{n}$ generate the unit ideal.
proof. When $s_{1}, \ldots, s_{k}$ generate the unit ideal, there will be an equation of the form $1=\sum r_{i} s_{i}$, and for any $N, 1=1^{N}=\left(\sum r_{i} s_{i}\right)^{N}$. If $N \geq n k$, then when the right side is expanded, every term will be divisible by $s_{i}^{n}$ for some $n$.
2.1.15. Proposition. Let $R$ be a noetherian ring, and let $M$ be a finite $R$-module.
noetheri-
anmodule
(ii) The set of submodules of $M$ satisfies the ascending chain condition.
(iii) Every nonempty set of submodules of $M$ contains a maximal member.

This concludes our review of rings and modules.

### 2.2 The Zariski Topology

As before, the affine space $\mathbb{A}^{n}$ is the space of $n$-tuples $\left(a_{1}, \ldots, a_{n}\right)$ of complex numbers. Algebraic geometry studies polynomial equations in terms of their solutions in affine space. Let $f_{1}, \ldots, f_{k}$ be polynomials in $x_{1}, \ldots, x_{n}$. The set of points of $\mathbb{A}^{n}$ that solve the system of equations

$$
\begin{equation*}
f_{1}=0, \ldots, f_{k}=0 \tag{2.2.1}
\end{equation*}
$$

(the locus of zeros of $f$ ) is a Zariski closed subset of $\mathbb{A}^{n}$. A Zariski open subset $U$ is a subset whose complement, which is the set of points not in $U$, is Zariski closed.

When it seems unlikely to cause confusion, we may abbreviate the notation for an indexed set, using a single letter. The polynomial algebra $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ may be denoted by $\mathbb{C}[x]$, and the system of equations 2.2.1 by $f=0$. The locus of solutions of the equations $f=0$ may be denoted by $V\left(f_{1}, \ldots, f_{k}\right)$ or by $V(f)$. Its points are called the zeros of the polynomials $f$.

We use analogous notation for infinite sets. If $\mathcal{F}$ is any set of polynomials, $V(\mathcal{F})$ denotes the set of points of affine space at which all elements of $\mathcal{F}$ are zero. In particular, if $I$ is an ideal of the polynomial ring, $V(I)$ denotes the set of points at which all elements of $I$ vanish.

The ideal $I$ of $\mathbb{C}[x]$ that is generated by the polynomials $f_{1}, \ldots, f_{k}$ is the set of combinations $r_{1} f_{1}+\cdots+r_{k} f_{k}$ with polynomial coefficients $r_{i}$. Some notations for this ideal are $\left(f_{1}, \ldots, f_{k}\right)$ and $(f)$. All elements of $I$ vanish on the zero set $V(f)$, so $V(f)=V(I)$. So the Zariski closed subsets of $\mathbb{A}^{n}$ can be described as the sets $V(I)$, where $I$ is an ideal.

We note a few simple relations among ideals and their zero sets here. To begin with, we note that an ideal $I$ isn't determined by its zero locus $V(I)$. For any $k>0$, the power $f^{k}$ has the same zeros as $f$.

The radical of an ideal $I$ of a ring $R$, which will be denoted by $\operatorname{rad} I$, is the set of elements $\alpha$ of $R$ such that some power $\alpha^{r}$ is in $I$.

$$
\begin{equation*}
\operatorname{rad} I=\left\{\alpha \in R \mid \alpha^{r} \in I \text { for some } r>0\right\} \tag{2.2.2}
\end{equation*}
$$

The radical of $I$ is an ideal that contains $I$. An ideal that is equal to its radical is a radical ideal. A prime ideal is a radical ideal.
2.2.3. Lemma. If $I$ is an ideal of the polynomial ring $\mathbb{C}[x]$, then $V(I)=V(\operatorname{rad} I)$.

Consequently, if $I$ and $J$ are ideals and if $\operatorname{rad} I=\operatorname{rad} J$, then $V(I)=V(J)$. The converse of this statement is also true: If $V(I)=V(J)$, then $\operatorname{rad} I=\operatorname{rad} J$. This is a consequence of the Strong Nullstellensatz that will be proved later in this chapter. (See (??).)

Because $(I \cap J)^{2} \subset I J \subset I \cap J$,

$$
\begin{equation*}
\operatorname{rad}(I J)=\operatorname{rad}(I \cap J) \tag{2.2.4}
\end{equation*}
$$

Also, $\operatorname{rad}(I \cap J)=(\operatorname{rad} I) \cap(\operatorname{rad} J)$.
2.2.5. Lemma. Let $I$ and $J$ be ideals of the polynomial ring $\mathbb{C}[x]$.
(i) If $I \subset J$, then $V(I) \supset V(J)$.
(ii) $V\left(I^{k}\right)=V(I)$.
(iii) $V(I \cap J)=V(I J)=V(I) \cup V(J)$.
(iv) If $I_{\nu}$ are ideals, then $V\left(\sum I_{\nu}\right)=\bigcap V\left(I_{\nu}\right)$.
proof. (iii) $V(I \cap J)=V(I J)$ because the two ideals have the same radical, and because $I$ and $J$ contain $I J$, $V(I J) \supset V(I) \cup V(J)$. To prove that $V(I J) \subset V(I) \cup V(J)$, we note that $V(I J)$ is the locus of common zeros of the products $f g$ with $f$ in $I$ and $g$ in $J$. Suppose that a point $p$ is a common zero: $f(p) g(p)=0$ for all $f$ in $I$ and all $g$ in $J$. If $f(p) \neq 0$ for some $f$ in $I$, we must have $g(p)=0$ for every $g$ in $J$, and then $p$ is a point of $V(J)$. If $f(p)=0$ for all $f$ in $I$, then $p$ is a point of $V(I)$. In either case, $p$ is a point of $V(I) \cup V(J)$.

Zariski closed sets are the closed sets in the Zariski topology on $\mathbb{A}^{n}$. The Zariski topology is very useful in algebraic geometry, though it is very different from the classical topology.

To verify that the Zariski closed sets are the closed sets of a topology, one must show that

- the empty set and the whole space are Zariski closed,
- the intersection $\bigcap C_{\nu}$ of an arbitrary family of Zariski closed sets is Zariski closed, and
- the union $C \cup D$ of two Zariski closed sets is Zariski closed.

The empty set and the whole space are the zero sets of the elements 1 and 0 , respectively. The other conditions follow from Lemma 2.2.5
ztopdimone
2.2.6. Example. The proper Zariski closed subsets of the affine line, or of a plane affine curve, are finite sets. The proper Zariski closed subsets of the affine plane are finite unions of points and curves. Let's omit the proofs of these facts. The corresponding facts for loci in the projective line and the projective plane have been noted before. (See 1.3.4) and 1.3.14).)
figure
(Caption: A Zariski closed subset of the affine plane (real locus).)
A subset $S$ of a topological space $X$ becomes a topological space with the induced topology. The closed (or open) subsets of $S$ in the induced topology are intersections $S \cap Y$, where $Y$ is closed (or open) in $X$.

The induced topology on a subset $S$ of $\mathbb{A}^{n}$ is called the Zariski topology too. A subset of $S$ is closed in the Zariski topology if it has the form $S \cap Y$ for some Zariski closed subset $Y$ of $\mathbb{A}^{n}$. If $S$ itself is a Zariski closed subset of $\mathbb{A}^{n}$, a closed subset of $S$ will be a closed subset of $\mathbb{A}^{n}$ that is contained in $S$.

Affine space also has a classical topology. A subset $U$ of $\mathbb{A}^{n}$ is open in the classical topology if, whenever a point $p$ is in $U$, all points sufficently near to $p$ are in $U$. Since polynomial functions are continuous, their zero sets are closed in the classical topology. Therefore Zariski closed sets are closed in the classical topology too.

When two topologies $T$ and $T^{\prime}$ on a set $X$ are given, $T^{\prime}$ is said to be coarser than $T$ if $T^{\prime}$ contains fewer closed sets or fewer open sets than $T$, and $T^{\prime}$ finer than $T$ if it contains more closed sets or more open sets than $T$. The Zariski topology is coarser than the classical topology. As the next proposition shows, it is much coarser.
2.2.7. Proposition. Every nonempty Zariski open subset of $\mathbb{A}^{n}$ is dense and path connected in the classical topology.
proof. The (complex) line $L$ through distinct points $p$ and $q$ of $\mathbb{A}^{n}$ is a Zariski closed set whose points can be written as $p+t(q-p)$, with $t$ in $\mathbb{C}$. It corresponds bijectively to the one-dimensional affine $t$-space $\mathbb{A}^{1}$, and the Zariski closed subsets of $L$ correspond to Zariski closed subsets of $\mathbb{A}^{1}$. They are the finite subsets, and $L$ itself.

Let $U$ be a nonempty Zariski open set, and let $C$ be its Zariski closed complement. To show that $U$ is dense in the classical topology, we choose distinct points $p$ and $q$ of $\mathbb{A}^{n}$, with $p$ in $U$. If $L$ is the line through $p$ and $q, C \cap L$ will be a Zariski closed subset of $L$, a finite set, that doesn't contain $p$. Its complement is $U \times L$. In the classical topology, the closure of $U \cap L$, will be the whole line $L$. It contains $q$. Then the closure of $U$ contains $q$, and since $q$ was arbitrary, the closure of $U$ is $\mathbb{A}^{n}$.

Next, let $L$ be the line through two points $p$ and $q$ of $U$. As before, $C \cap L$ will be a finite set of pints. In the classical topology, $L$ is a complex plane. The points $p$ and $q$ can be joined by a path in $L$ that avoids this finite set.

Though we will use the classical topology from time to time, the Zariski topology will appear more often. For this reason, we will refer to a Zariski closed subset simply as a closed set. Similarly, by an open set we mean a Zariski open set. We will mention the adjective "Zariski" only for emphasis.

## (2.2.8) irreducible closed sets

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The fact that the polynomial algebra is a noetherian ring has important consequences for the Zariski topology that we discuss here.

A topological space $X$ satisfies the descending chain condition on closed subsets if there is no infinite, strictly descending chain $C_{1}>C_{2}>\cdots$ of closed subsets of $X$. The descending chain condition on closed subsets is equivalent with the ascending chain condition on open sets.

A topological space that satisfies the descending chain condition on closed sets is called a noetherian space. In a noetherian space, every nonempty family $\mathcal{S}$ of closed subsets has a minimal member, one that doesn't contain any other member of $\mathcal{S}$, and every nonempty family of open sets has a maximal member. (See 2.1.11.)
2.2.9. Lemma. A noetherian topological space is quasicompact, that is, every open covering has a finite subcovering.
2.2.10. Proposition. With its Zariski topology, $\mathbb{A}^{n}$ is a noetherian space.
proof. Suppose that a strictly descending chain $C_{1}>C_{2}>\cdots$ of closed subsets of $\mathbb{A}^{n}$ is given. Let $I_{j}$ be the ideal of all elements of the polynomial ring $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ that are identically zero on $C_{j}$. Then $C_{j}=V\left(I_{j}\right)$. Since $C_{j}>C_{j+1}, \quad V\left(I_{j}\right)>V\left(I_{j+1}\right)$. Therefore $I_{j}<I_{j+1}$. The ideals $I_{j}$ form a strictly increasing chain. Since $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ is noetherian, this chain is finite. Therefore the strictly decreasing chain $C_{j}$ is finite too
2.2.11. Definition. A topological space $X$ is irreducible if it isn't the union of two proper closed subsets.

Another way to say that a topological space $X$ is irreducible is this:
2.2.12. If $C$ and $D$ are closed subsets of $X$, and if $X=C \cup D$, then $X=C$ or $X=D$.

The concept of irreducibility is useful primarily for noetherian spaces. The only irreducible subsets of a Hausdorff space are its points. In particular, with the classical topology, the only irreducible subsets of affine space are points.

Irreducibility is somewhat analogous to connectedness. A topological space is connected if it isn't the union $C \cup D$ of two proper disjoint closed subsets. However, the condition that a space be irreducible is much more restrictive because, in Definition 2.2.11, the closed sets $C$ and $D$ aren't required to be disjoint. In the Zariski topology on the affine plane, lines are irreducible closed sets. The union of two intersecting lines is connected, but not irreducible.
2.2.13. Lemma. The following conditions on topological space $X$ are equivalent.

- X is irreducible.
- The intersection $U \cap V$ of two nonempty open subsets $U$ and $V$ of $X$ is nonempty.
- Every nonempty open subset $U$ of $X$ is dense (its closure is $X$ ).
2.2.14. Theorem. In a noetherian topological space, every closed subset is the union of finitely many irreducible closed sets.
proof. If a closed subset $C_{0}$ of a topological space $X$ isn't a union of finitely many irreducible closed sets, then it isn't irreducible, so it is a union $C_{1} \cup D_{1}$, where $C_{1}$ and $D_{1}$ are proper closed subsets of $C_{0}$, and therefore closed subsets of $X$. Since $C_{0}$ isn't a finite union of irreducible closed sets, $C_{1}$ and $D_{1}$ cannot both be finite unions of irreducible closed sets. Say that $C_{1}$ isn't such a union. We have the beginning $C_{0}>C_{1}$ of a chain of closed subsets. We repeat the argument, replacing $C_{0}$ by $C_{1}$, and we continue in this way, to construct an infinite, strictly descending chain $C_{0}>C_{1}>C_{2}>\cdots$. So $X$ isn't a noetherian space.
2.2.15. Definition. An affine variety is an irreducible closed subset of affine space $\mathbb{A}^{n}$.
noethq-
comp
deschain
de-
firrspace
ir-
redspacetwo
ir-
redlemma
unionirred
defadffvar

Theorem 2.2.14 tells us that every closed subset of $\mathbb{A}^{n}$ is a finite union of affine varieties. Since an affine variety is irreducible, it is connected in the Zariski topology. It is also connected in the classical topology, but this isn't very easy to prove. We may not get to the proof.

The closure of a subset $S$ of a topological space $X$ is the smallest closed subset that contains $S$. The closure exists because it is the intersection of all closed subsets that contain $S$.
2.2.16. Lemma. (i) Let $Z$ be a subspace of a topological space $X$, let $S$ be a subset of $Z$, and let $\bar{S}$ denote the closure of $S$ in $X$. The closure of $S$ in $Z$ is the intersection $\bar{S} \cap Z$.
(ii) Let $\bar{Z}$ be the closure of a subspace $Z$ of a topological space $X$. Then $\bar{Z}$ is irreducible if and only if $Z$ is irreducible.
(iii) A nonempty open subspace $W$ of an irreducible space $X$ is irreducible.
proof. (ii) Let $Z$ be an irreducible subset of $X$, and suppose that its closure $\bar{Z}$ is the union $\bar{C} \cup \bar{D}$ of two closed sets $\bar{C}$ and $\bar{D}$. Then $Z$ is the union of the sets $C=\bar{C} \cap Z$ and $D=\bar{D} \cap Z$, and they are closed in $\underline{Z}$. Therefore $Z$ is one of those two sets; say $Z=C$. Then $Z \subset \bar{C}$, and since $\bar{C}$ is closed, $\bar{Z} \subset \bar{C}$. Because $\bar{C} \subset \bar{Z}$ as well, $\bar{C}=\bar{Z}$. Conversely, suppose that the closure $\bar{Z}$ of a subset $Z$ of $X$ is irreducible, and that $Z$ is a union $C \cup D$ of closed subsets. Then $\bar{Z}=\bar{C} \cup \bar{D}$, and therefore $\bar{Z}=\bar{C}$ or $\bar{Z}=\bar{D}$, say $\bar{Z}=\bar{C}$ So $Z=\bar{C} \cap Z=C$. Then $C$ is not a proper subset.
(iii) The closure of $W$ is the irreducible space $X$.

## (2.2.17) noetherian induction

In a noetherian space $X$ one may be able to use noetherian induction in proofs. Suppose that a statement $\mathcal{S}$ is to be proved for every nonempty closed subset of $X$. Then it suffices to prove it for $X$ under the assumption that $\mathcal{S}$ is true for every proper closed subset of $X$.

Or, to prove a statement $\mathcal{S}$ for every affine variety $X$, it is permissible prove it for $X$ under the assumption that $\mathcal{S}$ is true for every proper closed subvariety of $X$.

The justification of noetherian induction is similar to the justification of complete induction.

## (2.2.18) the coordinate algebra of a variety

2.2.19. Proposition. The affine varieties in $\mathbb{A}^{n}$ are the sets $V(P)$, where $P$ is a prime ideal of the polynomial algebra $\mathbb{C}[x]=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$. If $P$ is a radical ideal, then $V(P)$ is an affine variety if and only if $P$ is a prime ideal.

We will use this proposition in the next section, where we give a few examples of varieties, but we defer the proof to Section 2.5, where the proposition will be proved in a more general form. (See Proposition 2.5.13).)

As before, an algebra is a ring that contains the complex numbers.
2.2.20. Definition. Let $P$ be a prime ideal of the polynomial ring $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$, and let $V$ be the affine variety $V(P)$ in $\mathbb{A}^{n}$. The coordinate algebra of $V$ is the quotient algebra $A=\mathbb{C}[x] / P$.

Geometric properties of the variety are reflected in algebraic properties of its coordinate algebra and vice versa.
In a primitive sense, one can regard the geometry of an affine variety $V$ as given by closed subsets and incidence relations - the inclusion of one closed set into another, as when a point lies on a line. A finer study of the geometry takes into account other things tangency for instance, but it is reasonable to begin by studying incidences $C^{\prime} \subset C$ among closed subvarieties. Such incidences translate into inclusions $P^{\prime} \supset P$ in the opposite direction among prime ideals.

### 2.3 Some affine varieties

This section contains a few simple examples of varieties.
2.3.1. A point $p=\left(a_{1}, \ldots, a_{n}\right)$ of affine space $\mathbb{A}^{n}$ is the set of solutions of the $n$ equations $x_{i}-a_{i}=0, i=$ $1, \ldots, n$. A point is a variety because the polynomials $x_{i}-a_{i}$ generate a maximal ideal in the polynomial algebra $\mathbb{C}[x]$, and a maximal ideal is a prime ideal. We denote the maximal ideal that corresponds to the point $p$ by $\mathfrak{m}_{p}$. It is the kernel of the substitution homomorphism $\pi_{p}: \mathbb{C}[x] \rightarrow \mathbb{C}$ that evaluates a polynomial $g\left(x_{1}, \ldots, x_{n}\right)$ at $p: \pi_{p}(g(x))=g\left(a_{1}, \ldots, a_{n}\right)=g(p)$. As here, we denote the homomorphism that evaluates a polynomial at $p$ by $\pi_{p}$.

The coordinate algebra of a point $p$ is the quotient algebra $\mathbb{C}[x] / \mathfrak{m}_{p}$. It is also called the residue field at $p$, and it will be denoted by $k(p)$. The residue field $k(p)$ is isomorphic to the image of $\pi_{p}$, which is the field $\mathbb{C}$ of complex numbers, but it is a particular quotient of the polynomial ring.
2.3.2. The varieties in the affine line $\mathbb{A}^{1}$ are points, and the whole line $\mathbb{A}^{1}$. The varieties in the affine plane $\mathbb{A}^{2}$ are points, plane affine curves, and the whole plane.

This is true because the varieties correspond to the prime ideals of the polynomial ring. The prime ideals of $\mathbb{C}\left[x_{1}, x_{2}\right]$ are the maximal ideals, the principal ideals generated by irreducible polynomials, and the zero ideal. The proof of this is an exercise.
2.3.3. The set $X$ of solutions of a single irreducible polynomial equation $f_{1}\left(x_{1}, \ldots, x_{n}\right)=0$ in $\mathbb{A}^{n}$ is a variety that is called an affine hypersurface.

For instance, the special linear group $S L_{2}$, the group of complex $2 \times 2$ matrices with determinant 1 , is a hypersurface in $\mathbb{A}^{4}$. It is the locus of zeros of the irreducible polynomial $x_{11} x_{22}-x_{12} x_{21}-1$.

The reason that an affine hypersurface is a variety is that an irreducible element of a unique factorization domain is a prime element, and a prime element generates a prime ideal. The polynomial ring $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ is a unique factorization domain.
2.3.4. A hypersurface in the affine plane $\mathbb{A}^{2}$ is a plane affine curve.

A line in the plane, the locus of a linear equation $a x+b y-c=0$, is a plane affine curve. Its coordinate algebra is isomorphic to a polynomial ring in one variable. Every line is isomorphic to the affine line $\mathbb{A}^{1}$.
2.3.5. Let $p=\left(a_{1}, \ldots, a_{n}\right)$ and $q=\left(b_{1}, \ldots, b_{n}\right)$ be distinct points of $\mathbb{A}^{n}$. The point pair $(p, q)$ is the closed set defined by the system of $n^{2}$ equations $\left(x_{i}-a_{i}\right)\left(x_{j}-b_{j}\right)=0,1 \leq i, j \leq n$. A point pair isn't a variety because the ideal $I$ generated by the polynomials $\left(x_{i}-a_{i}\right)\left(x_{j}-b_{j}\right)$ isn't a prime ideal. The next proposition, which follows from the Chinese Remainder Theorem[2.1.6, describes the ideal $I$.
2.3.6. Proposition. The ideal of polynomials that vanish on a point pair $p, q$ is the product $\mathfrak{m}_{p} \mathfrak{m}_{q}$ of the maximal ideals at the points, and the quotient algebra $\mathbb{C}[x] / I$ is isomorphic to the product algebra $\mathbb{C} \times \mathbb{C}$.

### 2.4 Hilbert's Nullstellensatz

2.4.1. Nullstellensatz (version 1). Let $\mathbb{C}[x]$ be the polynomial algebra in the variables $x_{1}, \ldots, x_{n}$. There are bijective correspondences between the following sets.

- points $p$ of the affine space $\mathbb{A}^{n}$,
- algebra homomorphisms $\pi_{p}: \mathbb{C}[x] \rightarrow \mathbb{C}$,
- maximal ideals $\mathfrak{m}_{p}$ of $\mathbb{C}[x]$.

If $p=\left(a_{1}, \ldots, a_{n}\right)$ is a point of $\mathbb{A}^{n}$, the homomorphism $\pi_{p}$ evaluates a polynomial at $p: \pi_{p}(g)=g\left(a_{1}, \ldots ., a_{n}\right)=$ $g(p)$. The maximal ideal $\mathfrak{m}_{p}$ is the kernel of $\pi_{p}$, which is the ideal generated by the linear polynomials $x_{1}-a_{1}, \ldots, x_{n}-a_{n}$.

It is obvious that every algebra homomorphism $\mathbb{C}[x] \rightarrow \mathbb{C}$ is surjective and that its kernel is a maximal ideal. It isn't obvious that every maximal ideal of $\mathbb{C}[x]$ is the kernel of such a homomorphism. The proof can be found manywhere ${ }^{1}$

The Nullstellensatz gives us a way to describe the closed set $V(I)$ of zeros of an ideal $I$ in affine space in terms of maximal ideals. The points of $V(I)$ are those at which all elements of $I$ vanish. They are the points $p$ such that the ideal $I$ is contained in $\mathfrak{m}_{p}$.

[^0]Thus
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$$
\begin{equation*}
V(I)=\left\{p \in \mathbb{A}^{n} \mid I \subset \mathfrak{m}_{p}\right\} \tag{2.4.2}
\end{equation*}
$$

2.4.3. Proposition. Let I be an ideal of the polynomial ring $\mathbb{C}[x]$. If the zero locus $V(I)$ is empty, then $I$ is the unit ideal.
proof. Every ideal $I$ that is not the unit ideal is contained in a maximal ideal (Corollary 2.1.13).
2.4.4. Nullstellensatz (version 2). Let A be a finite-type algebra. There are bijective correspondences between the following sets:

- algebra homomorphisms $\bar{\pi}: A \rightarrow \mathbb{C}$,
- maximal ideals $\overline{\mathfrak{m}}$ of $A$.

The maximal ideal $\overline{\mathfrak{m}}$ that corresponds to a homomorphism $\bar{\pi}$ is the kernel of $\bar{\pi}$.
If $A$ is presented as a quotient of a polynomial ring, say $A \approx \mathbb{C}\left[x_{1}, \ldots, x_{n}\right] / I$, then these sets also correspond bijectively to points of the set $V(I)$ of zeros of $I$ in $\mathbb{A}^{n}$.

The symbol $\approx$ indicates an isomorphism here. As before, a finite-type algebra is an algebra that can be generated by a finite set of elements.
proof. We choose a presentation of $A$ as a quotient of a polynomial ring to identify $A$ with a quotient $\mathbb{C}[x] / I$. The Correspondence Theorem tells us that maximal ideals of $A$ correspond to maximal ideals of $\mathbb{C}[x]$ that contain $I$. Those maximal ideals correspond to points of $V(I)$.

Let $\tau$ denote the canonical homomorphism $\mathbb{C}[x] \rightarrow A$. The Mapping Property 2.1.3, applied to $\tau$, tells us that homomorphisms $A \xrightarrow{\bar{\pi}} \mathbb{C}$ correspond to homomorphisms $\mathbb{C}[x] \xrightarrow{\pi} \mathbb{C}$ whose kernels contain $I$. Those homomorphisms also correspond to points of $V(I)$.

2.4.6. Strong Nullstellensatz. Let I be an ideal of the polynomial algebra $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$, and let $V$ be the locus of zeros of $I$ in $\mathbb{A}^{n}: V=V(I)$. If a polynomial $g(x)$ vanishes at every point of $V$, then $I$ contains $a$ power of $g$.
proof. This beautiful proof is due to Rainich. Let $g(x)$ be a polynomial that is identically zero on $V$. We are to show that $I$ contains a power of $g$. If $g$ is the zero polynomial, it is in $I$, so we may assume that $g$ isn't zero.

The Hilbert Basis Theorem tells us that $I$ is a finitely generated ideal. Let $f=\left\{f_{1}, \ldots, f_{k}\right\}$ be a set of generators. We introduce an new variable $y$. Let $W$ be the locus of solutions of the $k+1$ equations

$$
\begin{equation*}
f_{1}(x)=\cdots=f_{k}(x)=0 \quad \text { and } \quad g(x) y-1=0 \tag{2.4.7}
\end{equation*}
$$

in the $n+1$-dimensional affine space with coordinates $\left(x_{1}, \ldots, x_{n}, y\right)$. Suppose that we have a solution $x$ of the equations $f(x)=0$, say $\left(x_{1}, \ldots, x_{n}\right)=\left(a_{1}, \ldots, a_{n}\right)$. Then $a$ is a point of $V$, and our hypothesis tells us that $g(a)=0$ too. So there can be no $b$ such that $g(a) b=1$. There is no point $\left(a_{1}, \ldots, a_{n}, b\right)$ that solves the equations 2.4.7): The locus $W$ is empty. Proposition 2.4.3 tells us that the polynomials $f_{1}, \ldots, f_{k}, g y-1$ generate the unit ideal of $\mathbb{C}\left[x_{1}, \ldots, x_{n}, y\right]$. There are polynomials $p_{1}(x, y), \ldots, p_{k}(x, y)$ and $q(x, y)$ such that

$$
\begin{equation*}
p_{1} f_{1}+\cdots+p_{k} f_{k}+q(g y-1)=1 \tag{2.4.8}
\end{equation*}
$$

The ring $R=\mathbb{C}[x, y] /(g y-1)$ can be described as the one obtained by adjoining an inverse of $g$ to the polynomial ring $\mathbb{C}[x]$. The residue of $y$ is the inverse of $g$. Since $g$ isn't zero, $\mathbb{C}[x]$ is a subring of $R$. In $R$, $g y-1=0$. So the equation 2.4.8) becomes $p_{1} f_{1}+\cdots+p_{k} f_{k}=1$. When we multiply both sides of this
equation by a large power $g^{N}$ of $g$, we can use the equation $g y=1$, which is true in $R$, to cancel all occurences of $y$ in the polynomials $p_{i}(x, y)$. Let $h_{i}(x)$ denote the polynomial in $x$ that is obtained by cancelling $y$ in $g^{N} p_{i}$. Then

$$
h_{1}(x) f_{1}(x)+\cdots+h_{k}(x) f_{k}(x)=g^{N}(x)
$$

is a polynomial equation that is true in $R$ and in its subring $\mathbb{C}[x]$. Since $f_{1}, \ldots, f_{k}$ are in $I$, this equation shows that $g^{N}$ is in $I$.

### 2.4.9. Corollary.et $\mathbb{C}[x]$ denote the polynomial ring $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$.

(i) Let $P$ be a prime ideal of $\mathbb{C}[x]$, and let $V=V(P)$ be the variety of zeros of $P$ in $\mathbb{A}^{n}$. If a polynomial $g$ vanishes at every point of $V$, then $g$ is an element of $P$.
(ii) Let $f$ be an irreducible polynomial in $\mathbb{C}[x]$. If a polynomial $g$ vanishes at every point of $V(f)$, then $f$ divides $g$.
(iii) Let $I$ and $J$ be ideals of $\mathbb{C}[x]$. Then $V(I) \supset V(J)$ if and only if $\operatorname{rad} I \subset \operatorname{rad} J$, and $V(I)>V(J)$ if and only if $\operatorname{rad} I>\operatorname{rad} J$. (See (2.2.2).)

### 2.4.10. Examples.

(i) Let $I$ be the ideal generated by $y^{5}$ and $y^{2}-x^{3}$ in the polynomial algebra $\mathbb{C}[x, y]$ in two variables. The origin $y=x=0$ is the only common zero of these polynomials in the affine plane, and the polynomial $x$ also vanishes at the origin. The Strong Nullstellensatz predicts that $I$ contains a power of $x$. This is verified by the following equation:

$$
y y^{5}-\left(y^{4}+y^{2} x^{3}+x^{6}\right)\left(y^{2}-x^{3}\right)=x^{9}
$$

(ii) We may regard pairs $A, B$ of $n \times n$ matrices as points of an affine space $\mathbb{A}^{2 n^{2}}$ with coordinates $a_{i j}, b_{i j}$, $1 \leq i, j \leq n$. The pairs of commuting matrices $(A B=B A)$ form a closed subset of $\mathbb{A}^{2 n^{2}}$, the locus of common zeros of the $n^{2}$ polynomials $c_{i j}$ that compute the entries of the matrix $A B-B A$ :

$$
\begin{equation*}
c_{i j}(a, b)=\sum_{\nu} a_{i \nu} b_{\nu j}-b_{i \nu} a_{\nu j} \tag{2.4.11}
\end{equation*}
$$

Let $I$ denote the ideal of the polynomial algebra $\mathbb{C}[a, b]$ generated by the polynomials $c_{i j}$. Then $V(I)$ is the set of pairs of commuting complex matrices. The Strong Nullstellensatz asserts that if a polynomial $g(a, b)$ vanishes on every pair of commuting matrices, some power of $g$ is in $I$. Is $g$ itself in $I$ ? It is a famous conjecture that $I$ is a prime ideal. If so, $g$ would be in $I$. Proving the conjecture would establish your reputation as a mathematician, but I don't recommend spending very much time on it right now.

### 2.5 The Spectrum

The Nullstellensatz allows us to associate a set of points to a finite-type domain $A$ without reference to a presentation. We can do this because the maximal ideals of $A$ and the homomorphisms $A \rightarrow \mathbb{C}$ don't depend on the presentation. When a finite-type domain $A$ is presented as a quotient $\mathbb{C}[x] / P$ of a polynomial ring, where $P$ is a prime ideal, $A$ becomes the coordinate algebra of the variety $V(P)$ in affine space. Then the points of $V(P)$ correspond to maximal ideals of $A$ and also to homomorphisms $A \rightarrow \mathbb{C}$.

When a finite-type domain $A$ is given without a presentation, we replace the variety $V(P)$ by an abstract set of points, the spectrum of $A$, that we denote by $\operatorname{Spec} A$ and call an affine variety. We put one point $p$ into the spectrum for every maximal ideal of $A$, and then we turn around and denote the maximal ideal that corresponds to a point $p$ by $\overline{\mathfrak{m}}_{p}$. The Nullstellensatz tells us that $p$ also corresponds to a homomorphism $A \rightarrow \mathbb{C}$ whose kernel is $\overline{\mathfrak{m}}_{p}$. We denote that homomorphism by $\bar{\pi}_{p}$. In analogy with 2.2.20, the domain $A$ is called the coordinate algebra of the affine variety $\operatorname{Spec} A$. To work with $\operatorname{Spec} A$, we may interpret its points as maximal ideals or as homomorphisms to $\mathbb{C}$, whichever is convenient.

When defined in this way, the variety $\operatorname{Spec} A$ isn't embedded into affine space, but because $A$ is a finitetype algebra, it can be presented as a quotient $\mathbb{C}[x] / P$. When this is done, points of Spec $A$ correspond to points of the subset $V(P)$ in $\mathbb{A}^{n}$.

Even when the coordinate ring $A$ of an affine variety is presented as $\mathbb{C}[x] / P$, we may denote the variety by Spec $A$ rather than by $V(P)$.

Let $X=\operatorname{Spec} A$. The elements of $A$ define (complex-valued) functions on $X$ as follows: A point $p$ of $X$ corresponds to a homomorphism $A \xrightarrow{\bar{\pi}_{p}} \mathbb{C}$. If $\alpha$ is an element of $A$, the value of the function $\alpha$ at $p$ is defined to be $\bar{\pi}_{p}(\alpha)$ :

$$
\begin{equation*}
\alpha(p) \stackrel{\text { def }}{=} \bar{\pi}_{p}(\alpha) \tag{2.5.1}
\end{equation*}
$$

The kernel $\overline{\mathfrak{m}}_{p}$ of $\bar{\pi}_{p}$ is the set of elements $\alpha$ of the coordinate algebra $A$ such that $\alpha(p)=0$ :

$$
\overline{\mathfrak{m}}_{p}=\{\alpha \in A \mid \alpha(p)=0\}
$$

The functions defined by the elements of $A$ are called regular functions on $X$. (See Proposition 2.7.2 below.)

### 2.5.2. Lemma. The regular function determined by distinct elements $\alpha$ and $\beta$ of $A$ are distinct.

proof. We replace $\alpha$ by $\alpha-\beta$. Then what is to be shown is that, if the function determined by an element $\alpha$ is the zero function, then $\alpha$ is the zero element.

We present $A$ as $\mathbb{C}[x] / P, x=x_{1}, \ldots, x_{n}$, where $P$ is a prime ideal. Then $X$ is the locus of zeros of $P$ in $\mathbb{A}^{n}$, and Corollary ?? (iii) tells us that $P$ is the ideal of all elements that are zero on $X$.

Let $g(x)$ be a polynomial that represents $\alpha$. If $p$ is a point of $X$ at which $\alpha$ is zero, then $g(p)=0$. So if $\alpha$ is the zero function, then $g$ is in $P$, and therefore $\alpha=0$.

For example, the spectrum $\operatorname{Spec} \mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ of the polynomial algebra is the affine space $\mathbb{A}^{n}$. The homomorphism $\pi_{p}: \mathbb{C}[x] \rightarrow \mathbb{C}$ that corresponds to a point $p=\left(a_{1}, \ldots, a_{n}\right)$ of $\mathbb{A}^{n}$ is evaluation at $p$. So $\pi_{p}(g)=g\left(a_{1}, \ldots, a_{n}\right)=g(p)$. The function defined by a polynomial $g(x)$ is simply the polynomial function.

Note. In modern terminology, the word "spectrum" is usually used to denote the set of prime ideals of a ring. This becomes important when one studies rings that aren't finite-type algebras. When working with finite-type algebras, there are enough maximal ideals. The other prime ideals aren't needed, so we have eliminated them.
2.5.3. Lemma. Let $A$ be a quotient $\mathbb{C}[x] / P$ of the polynomial ring $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ modulo a prime ideal $P$, so that $\operatorname{Spec} A$ becomes the closed subset $V(P)$ of $\mathbb{A}^{n}$. Then a point $p$ of $\operatorname{Spec} A$ becomes a point of $\mathbb{A}^{n}$ : $p=\left(a_{1}, \ldots, a_{n}\right)$. When an element $\alpha$ of $A$ is represented by a polynomial $g(x)$, the value of $\alpha$ at $p$ is $\alpha(p)=g(p)=g\left(a_{1}, \ldots, a_{n}\right)$.
proof. The point $p$ of $\operatorname{Spec} A$ gives us a diagram 2.4.5, with $\pi=\pi_{p}$ and $\bar{\pi}=\bar{\pi}_{p}$, and where $\tau$ is the canonical map $\mathbb{C}[x] \rightarrow A$. Then $\alpha=\tau(g)$, and

$$
\begin{equation*}
g(p) \stackrel{\text { defn } n}{=} \pi_{p}(g)=\bar{\pi}_{p} \tau(g)=\bar{\pi}_{p}(\alpha) \stackrel{\text { def } n}{=} \alpha(p) . \tag{2.5.4}
\end{equation*}
$$

Thus the value $\alpha(p)$ at a point $p$ of $\operatorname{Spec} A$ can be obtained by evaluating a polynomial $g$, though the polynomial $g$ isn't unique. The polynomial $g$ that represents the regular function $\alpha$ won't be unique unless $P$ is the zero ideal.

## (2.5.5) the Zariski topology on an affine variety

Let $X=\operatorname{Spec} A$ be an affine variety with coordinate algebra $A$. An ideal $\bar{J}$ of $A$ defines a locus in $X$, a closed subset, that we denote by $V_{X}(\bar{J})$. The points of $V_{X}(\bar{J})$ are the points of $X$ at which all elements of $\bar{J}$ vanish. This is analogous to 2.4.2):

$$
\begin{equation*}
V_{X}(\bar{J})=\left\{p \in \operatorname{Spec} A \mid \bar{J} \subset \overline{\mathfrak{m}}_{p}\right\} \tag{2.5.6}
\end{equation*}
$$

2.5.7. Lemma. Let $A$ be a finite-type domain, presented as $A=\mathbb{C}[x] / P \approx A$. An ideal $\bar{J}$ of $A$ corresponds to an ideal $J$ of $\mathbb{C}[x]$ that contains $P$, and if $V_{\mathbb{A}^{n}}(J)$ denotes the zero locus of $J$ in $\mathbb{A}^{n}$, then $V_{X}(\bar{J})=V_{\mathbb{A}^{n}}(J)$.

The properties of closed sets in affine space that are given in Lemmas 2.2.3 and 2.2.5 are true for closed subsets of an affine variety. In particular, $V_{X}(\bar{J})=V_{X}(\operatorname{rad} \bar{J})$, and $V_{X}(\overline{I J})=V_{X}(\bar{I} \cap \bar{J})=V_{X}(\bar{I}) \cup V_{X}(\bar{J})$.
2.5.8. Proposition. Let $\bar{J}$ be an ideal of a finite-type domain $A$, and let $X=\operatorname{Spec} A$. The zero set $V_{X}(\bar{J})$ is empty if and only if $\bar{J}$ is the unit ideal of $A$. If $X$ is empty, then $A$ is the zero ring.
proof. The only ideal that isn't contained in a maximal ideal is the unit ideal.
2.5.9. Note. In this section, we have put bars on the symbols $\overline{\mathfrak{m}}$ and $\bar{\pi}$ here in order to distinguish maximal ideals of $A$ from maximal ideals of $\mathbb{C}[x]$ and homomorphisms $A \rightarrow \mathbb{C}$ from homomorphisms $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] \rightarrow$ $\mathbb{C}$. In the future, we will put bars over the letters only when there is a danger of confusion. Most of the time, we will denote an ideal of $A$ by an ordinari letter such as $J$.

## (2.5.10) ideals whose zero sets are equal

2.5.11. Lemma. An ideal I of a noetherian ring $R$ contains a power of its radical.
proof. Since $R$ is noetherian, the ideal $\operatorname{rad} I$ is generated by a finite set of elements $\alpha=\left\{\alpha_{1}, \ldots, \alpha_{k}\right\}$, and for large $r, \alpha_{i}^{r}$ is in $I$. We can use the same large integer $r$ for every $i$. A monomial $\beta=\alpha_{1}^{e_{1}} \cdots \alpha_{k}^{e_{k}}$ of sufficiently large degree $n$ in $\alpha$ will be divisible $\alpha_{i}^{r}$ for at least one $i$, and therefore it will be in $I$. The monomials of degree $n$ generate $(\operatorname{rad} I)^{n}$, so $(\operatorname{rad} I)^{n} \subset I$. And as has been remarked, $I \subset \operatorname{rad} I$.
2.5.12. Corollary. Let $I$ and $J$ be ideals of a finite-type domain $A$, and let $X=\operatorname{Spec} A$. Then $V_{X}(I) \supset$ $V_{X}(J)$ if and only if $\operatorname{rad} I \subset \operatorname{rad} J$.

This follows from Corollary ?? bf (iii) and Lemma 2.5.7.
The next proposition includes Proposition 2.2.19 as a special case.
2.5.13. Proposition. Let $X=\operatorname{Spec} A$, where $A$ is a finite-type domain. The closed subset $V_{X}(P)$ defined by a radical ideal $P$ is irreducible if and only if $P$ is a prime ideal.
proof. Let $P$ be a radical ideal of $A$, and let $Y=V_{X}(P)$. Let $C$ and $D$ be closed subsets of $X$ such that $Y=C \cup D$. Say $C=V_{X}(I), D=V_{X}(J)$. We may suppose that $I$ and $J$ are radical ideals. Then the inclusion $C \subset Y$ implies that $I \supset P$. Similarly, $J \supset P$. Because $Y=C \cup D$, we also have $Y=V_{X}(I \cap J)=V_{X}(I J)$. So $I J \subset P$ (Corollary 2.5.12). If $P$ is a prime ideal, then $I=P$ or $J=P$, and therefore $C=Y$ or $D=Y$ (Lemma2.1.2iii'), and $Y$ is irreducible. Conversely, if $P$ isn't a prime ideal, there are ideals $I, J$ strictly larger than the radical ideal $P$, such that $I J \subset P(2.1 .2)$. Then $Y$ will be the union of the two proper closed subsets $V_{X}(I)$ and $V_{X}(J)$, so $Y$ isn't irreducible 2.5.12).

## (2.5.14) the nilradical

The nilradical of a ring is the set of its nilpotent elements. It is the radical of the zero ideal. The nilradical of a domain is the zero ideal. If a ring $R$ is noetherian, its nilradical will be nilpotent : some power of will be the zero ideal (Lemma 2.5.11.
2.5.15. Proposition. The nilradical of a noetherian ring $R$ is the intersection of the prime ideals of $R$.
proof. Let $x$ be an element of the nilradical $N$. Some power of $x$ is zero. Since the zero element is in every prime ideal, $x$ is in every prime ideal. Therefore $N$ is contained in every prime ideal. Conversely, let $x$ be an element not in $N$, i.e., not nilpotent. We show that there is a prime ideal that doesn't contain $x$. Let $\mathcal{S}$ be the set of ideals that don't contain any power of $x$. The zero ideal is one such ideal, so $\mathcal{S}$ isn't empty. Since $R$ is noetherian, $\mathcal{S}$ contains a maximal member $P$ 2.1.11. We show that $P$ is a prime ideal by showing that, if two ideals $A$ and $B$ are strictly larger than $P$, their product $A B$ isn't contained in $P$. Since $P$ is a maximal member of $\mathcal{S}, A$ and $B$ aren't in $\mathcal{S}$. They contain powers of $x$, say $x^{k} \in A$ and $x^{\ell} \in B$. Then $x^{k+\ell}$ is in $A B$ but not in $P$. Therefore $A B \not \subset P$.

Note. The conclusion of this proposition is true whether or not the ring $R$ is noetherian.

### 2.5.16. Corollary.

(i) Let $A$ be a finite-type algebra. An element that is in every maximal ideal of $A$ is nilpotent.
(ii) Let $A$ be a finite-type domain. The intersection of the maximal ideals of $A$ is the zero ideal.

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 onlocalizationproof. (i) Say that $A$ is presented as $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] / I$. Let $\alpha$ be an element of $A$ that is in every maximal ideal, and let $g(x)$ be a polynomial whose residue in $A$ is $\alpha$. Then $\alpha$ is in every maximal ideal of $A$ if and only if $g=0$ at all points of the variety $V_{\mathbb{A}}(I)$ in $\mathbb{A}^{n}$. If so, the Strong Nullstellensatz asserts that some power $g^{N}$ is in $I$. Then $\alpha^{N}=0$.
2.5.17. Corollary. An element $\alpha$ of a finite-type domain $A$ is determined by the function that it defines on $X=\operatorname{Spec} A$.
proof. It is enough to show that an element $\alpha$ that defines the zero function is the zero element. Such an element is in every maximal ideal (2.5.8), so $\alpha$ is nilpotent, and since $A$ is a domain, $\alpha=0$.

### 2.6 Localization

Let $s$ be a nonzero element of a domain $A$. The ring $A\left[s^{-1}\right]$, obtained by adjoining an inverse of $s$ to $A$ is called a localization of $A$. It is isomorphic to the quotient $A[z] /(s z-1)$ of the polynomial ring $A[z]$ by the principal ideal generated by $s z-1$. We will often denote this localization by $A_{s}$. If $A$ is a finite-type domain, so is $A_{s}$. In that case, if $X$ denotes the variety $=\operatorname{Spec} A, X_{s}$ will denote the variety $\operatorname{Spec} A_{s}$, and $X_{s}$ will be called a localization of $X$ too.
2.6.1. Proposition. (i) With terminology as above, points of the localization $X_{s}=\operatorname{Spec} A_{s}$ correspond bijectively to the open subset of $X$ of points at which the function defined by s is nonzero.
(ii) When we identify a localization $X_{s}$ with a subset of $X$, the Zariski topology on $X_{s}$ is the induced topology from $X$. So $X_{s}$ is an open subspace of $X$.
proof. (i) Let $p$ be a point of $X$, let $A \xrightarrow{\pi_{p}} \mathbb{C}$ be the corresponding homomorphism. If $s$ isn't zero at $p$, i.e., $c=s(p) \neq 0$, then $\pi_{p}$ extends uniquely to a homomorphism $A_{s} \rightarrow \mathbb{C}$ that sends $s^{-1}$ to $c^{-1}$. This gives us a unique point of $X_{s}$ whose image in $X$ is $p$. If $c=0$, then $\pi_{p}$ doesn't extend to $A_{s}$.
(ii) Let $C$ be a closed subset of $X$, say the zero set of a set of elements $a_{1}, \ldots, a_{k}$ of $A$. Then $C \cap X_{s}$ is the zero set in $X_{s}$ of those same elements, so it is closed in $X_{s}$. Conversely, let $D$ be a closed subset of $X_{s}$, say the zero set in $X_{s}$ of some elements $\beta_{1}, \ldots, \beta_{k}$, where $\beta_{i}=b_{i} s^{-n}$ with $b_{i}$ in $A$. We can use the same exponent $n$ for each $i$. Since $s^{-1}$ doesn't vanish on $X_{s}$, the elements $b_{i}$ and $\beta_{i}$ have the same zeros in $X_{s}$. If we let $C$ be the zero set of $b_{1}, \ldots, b_{k}$ in $X$, we will have $C \cap X_{s}=D$.

We usually identify $X_{s}$ as an open subset of $X$. Then the effect of adjoining the inverse is to throw out the points of $X$ at which $s$ vanishes. For example, the spectrum of the Laurent polynomial ring $\mathbb{C}\left[t, t^{-1}\right]$ becomes the complement of the origin in the affine line $\mathbb{A}^{1}=\operatorname{Spec} \mathbb{C}[t]$.

This illustrates one benefit of working with an affine variety without a fixing an embedding into affine space. If $X$ is embedded into $\mathbb{A}^{n}$, the localization $X_{s}$ wants to be in $\mathbb{A}^{n+1}$.

As is true for many open sets, the complement $X^{\prime}$ of the origin in the affine plane $\operatorname{Spec} \mathbb{C}\left[x_{1}, x_{2}\right]$ isn't a localization. Every polynomial that vanishes at the origin vanishes on an affine curve, which has points distinct from the origin. So the inverse of such a polynomial doesn't define a function on $X^{\prime}$. This is rather obvious, but in other situations, it may be hard to tell whether or not an open set that is given is a localization.

Localizations are important for two reasons:

- The relation between an algebra $A$ and a localization $A_{s}$ is easy to understand, and
- The localizations $X_{s}$ of an affine variety $X$ form a basis for the Zariski topology on $X$.

A basis for the topology on a topological space $X$ is a family $\mathcal{B}$ of open sets such that every open set is a union of open sets that are members of $\mathcal{B}$.

We verify the second statement marked with a bullet. We must show that if $U$ is an open subset of $X=$ $\operatorname{Spec} A$, then $U$ can be covered by sets of the form $X_{s}$. The complement of $X_{s}$ in $X$, the set $X-X_{s}$, is the zero set $V_{X}(s)$ of $s$. Let $C$ be the closed complement $X-U$ of $U$ in $X$. Since $C$ is closed, it is the set of common zeros of some elements $s_{1}, \ldots, s_{k}$ of $A$. Then $C$ is the intersection $\bigcap V_{X}\left(s_{i}\right)$ of zero sets, and $U$ is the union of those sets $X_{s_{i}}$.
2.6.2. Lemma. Let $U$ and $V$ be affine open subsets of an affine variety $X$.
(i) If $V$ is a localization of $U$ and $U$ is a localization of $X$, then $V$ is a localization of $X$.
(ii) If $V \subset U$ and $V$ is a localization of $X$, then $V$ is a localization of $U$.
(iii) Let $p$ be a point of $U \cap V$. There is an open set $Z$ containing $p$ that is a localization of $U$ and also $a$ localization of $V$.
proof. (i) Say that $X=\operatorname{Spec} A, \quad U=X_{s}=\operatorname{Spec} A_{s}$ and $V=U_{t}=\operatorname{Spec}\left(A_{s}\right)_{t}$. Then $t$ is an element of $A_{s}$, say $t=r s^{-k}$ with $r$ in $A$. The localizations $\left(A_{s}\right)_{t}$ and $\left(A_{s}\right)_{r}$ are equal, and $\left(A_{s}\right)_{r}=A_{s r}$. So $V=X_{s r}$.
(ii) Say that $X=\operatorname{Spec} A, U=\operatorname{Spec} B$, and $V=\operatorname{Spec} A_{t}$, where $t$ is a nonzero element of $A$. A regular function on $X$ restricts to a regular on $U$, and a regular functon on $U$ to a regular function on $V$. The restrictions define homomorphisms $A \rightarrow B \rightarrow A_{t}$, which are injective because $A \subset A_{t}$. Since $t$ is in $B, B_{t} \subset A_{t}$, and therefore $R_{t}=A_{t}$.
(iii) Since localizations form a basis for the topology, $U \cap V$ contains a localization $X_{s}$ of $X$ that contains $p$. By (ii), $X_{s}$ is also a localization of $U$ and a localization of $V$.

## (2.6.3) extension and contraction of ideals

Let $A \subset B$ be the inclusion of a ring $A$ as a subring of a ring $B$. The extension of an ideal $I$ of $A$ is the ideal $I B$ of $B$ generated by $I$. Its elements are finite sums $\sum_{i} z_{i} b_{i}$ with $z_{i}$ in $I$ and $b_{i}$ in $B$. The contraction of an ideal $J$ of $B$ is the intersection $J \cap A$. It is an ideal of $A$.

If $A_{s}$ is a localization of $A$ and $I$ is an ideal of $A$, the elements of the extended ideal $I A_{s}$ are fractions of the form $z s^{-k}$, with $z$ in $I$. We denote this extended ideal by $I_{s}$.

### 2.6.4. Lemma.

(i) Let $A_{s}$ be a localization of a domain $A$, let $J$ be an ideal of $A_{s}$ and let $I=J \cap A$. Then $J=I_{s}$. Every ideal of $A_{s}$ is the extension of an ideal of $A$.
(ii) Let $P$ be a prime ideal of $A$. If $s$ is an element of $A$ that isn't in $P$, the extended ideal $P_{s}$ is a prime ideal of $A_{s}$. If $s$ is in $P$, the extended ideal $P_{s}$ is the unit ideal.

## (2.6.5) localizing a module

Let $A$ be a domain, and let $M$ be an $A$-module. torsion element of $M$ is an element that is annihilated by some nonzero element $s$ of $A: s m=0$. A nonzero element $m$ such that $s m=0$ is called an $s$-torsion element.

The set of all torsion elements of $M$ is its torsion submodule, and a module whose torsion submodule is zero is torsion-free.

Let $s$ be a nonzero element of a domain $A$. The localization $M_{s}$ of an $A$-module $M$ is defined in the natural way, as the $A_{s}$-module whose elements are equivalence classes of fractions $m / s^{r}=m s^{-r}$, with $m$ in $M$ and $r \geq 0$. An alternate notation for the localization $M_{s}$ is $M\left[s^{-1}\right]$. The only complication comes from the fact that there may be $s$-torsion elements in $M$. If $m s=0$ and $s$ is in $S$, then $m$ must map to zero in $M_{s}$, because we will have $\mathrm{mss}^{-1}=m$ in $M_{s}$.

To define $M_{s}$, it suffices to modify the equivalence relation. Two fractions $m_{1} s^{-r_{1}}$ and $m_{2} s^{-r_{2}}$ are defined to be equivalent $m_{1} s^{r_{2}+n}=m_{2} s^{r_{1}+n}$ when $n$ is sufficiently large. This takes care of torsion, and $M_{s}$ becomes an $A_{s}$-module. There will be a homomorphism $M \rightarrow M_{s}$ that sends an element $m$ to the fraction $m / 1$.

This is also how one localizes a ring that isn't a domain.

## multiplicative systems

To work with the inverses of finitely many nonzero elements, one may simply adjoin the inverse of their product. For working with an infinite set of inverses, the concept of a multiplicative system is useful.
inverseexamples
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A multiplicative system $S$ in a domain $A$ is a subset that consists of nonzero elements, is closed under multiplication, and contains 1 . If $S$ is a multiplicative system, the ring of $S$-fractions $A S^{-1}$ is the ring obtained by adjoining inverses of all elements of $S$. Its elements are equivalence classes of fractions $a s^{-1}$ with $a$ in $A$ and $s$ in $S$, the equivalence relation and the laws of composition being the usual ones for fractions. The ring $A S^{-1}$ will be called a localization too. For emphasis, the ring obtained by inverting a single element $s$ may be called a simple localization.
2.6.7. Examples. (i) The set consisting of the powers of a nonzero element $s$ of a domain $A$ is a multiplicative system. Its ring of fractions is the simple localization $A_{s}=A\left[s^{-1}\right]$.
(ii) The set $S$ of all nonzero elements of a domain $A$ is a multiplicative system. Its ring of fractions is the field of fractions of $A$.
(iii) An ideal $P$ of a domain $A$ is a prime ideal if and only if its complement, the set of elements of $A$ not in $P$, is a multiplicative system.
2.6.8. Proposition. Let $S$ be a multiplicative system in a domain $A$, and let $A^{\prime}$ be the localization $A S^{-1}$.
(i) Let $I$ be an ideal of $A$. The extended ideal $I A^{\prime}$ 2.6.3 is the set $I S^{-1}$ whose elements are classes of fractions $x s^{-1}$, with $x$ in I and sin $S$. The extended ideal is the unit ideal if and only if I contains an element of $S$.
(ii) Let $J$ be an ideal of the localization $A^{\prime}$ and let I denote its contraction $J \cap A$. The extended ideal $I A^{\prime}$ is equal to $J: J=(J \cap A) A^{\prime}$.
(iii) If $Q$ is a prime ideal of $A$ and if $Q \cap S$ is empty, the extended ideal $Q^{\prime}=Q A^{\prime}$ is a prime ideal of $A^{\prime}$, and the contraction $Q^{\prime} \cap A$ is equal to $Q$. If $Q \cap S$ isn't empty, the extended ideal is the unit ideal. Thus prime ideals of $A S^{-1}$ correspond bijectively to prime ideals of $A$ that don't meet $S$.
2.6.9. Corollary. Every localization $A S^{-1}$ of a noetherian domain $A$ is noetherian.

## (2.6.10) a general principle

An important, though elementary, principle for working with fractions is that any finite sequence of computations in a localization $A S^{-1}$ will involve only finitely many denominators, and can therefore be done in a simple localization $A_{s}$, where $s$ is a common denominator for the fractions that occur. (This has been mentioned before.)

## (2.6.11) localizing a module

Let $S$ be a multiplicative system in a domain $A$. The localization $M S^{-1}$ of an $A$-module $M$ is defined in a way analogous to the one used for simple localizations, as the $A S^{-1}$-module whose elements are equivalence classes of fractions $m s^{-1}$ with $m$ in $M$ and $s$ in $S$. To take care of possible torsion, two fractions $m_{1} s_{1}^{-1}$ and $m_{2} s_{2}^{-1}$ are defined to be equivalent if there is an element $\widetilde{s} \in S$ such that $m_{1} s_{2} \widetilde{s}=m_{2} s_{1} \widetilde{s}$. Then $m s_{1}^{-1}=0$ if and only if $m \widetilde{s}=0$ for some $\widetilde{s}$ in $S$. As with simple localizations, there will be a homomorphism $M \rightarrow M S^{-1}$ that sends an element $m$ to the fraction $m / 1$.
2.6.12. Proposition. Let $S$ be a multiplicative system in a domain $A$.
(i) Localization is an exact functor: A homomorphism $M \xrightarrow{\varphi} N$ of A-modules induces a homomorphism $M S^{-1} \xrightarrow{\varphi^{\prime}} N S^{-1}$ of $A S^{-1}$-modules, and if $M \xrightarrow{\varphi} N \xrightarrow{\psi} P$ is an exact sequence of $A$-modules, the localized sequence $M S^{-1} \xrightarrow{\varphi^{\prime}} N S^{-1} \xrightarrow{\psi^{\prime}} P S^{-1}$ is exact.
(ii) Let $M$ be an $A$-module. and let $N$ be an $A S^{-1}$-module. Homomorphisms of $A S^{-1}$-modules $M S^{-1} \rightarrow N$ correspond bijectively to homomorphisms of $A$-modules $M \rightarrow N$, when $N$ is made into an $A$-module by restriction of scalars.
(iii) If multiplication by $s$ is an injective map $M \rightarrow M$ for every $s$ in $S$, then $M \subset M S^{-1}$. If multiplication by every $s$ is a bijective map $M \rightarrow M$, then $M \approx M S^{-1}$.

### 2.7 Morphisms of Affine Varieties

Morphisms are the allowed maps between varieties. Morphisms between affine varieties are defined below. They correspond to algebra homomorphisms in the opposite direction between their coordinate algebras.

Morphisms of projective varieties require more thought. They will be defined in the next chapter.

## regular functions

The function field $K$ of an affine variety $X=\operatorname{Spec} A$ is the field of fractions of $A$. A rational function on $X$ is a nonzero element of the function field $K$.

As we have seen, 2.5.1) elements of the coordinate algebra $A$ define functions on $X$. The value of the function $\alpha$ at a point $p$ is $\alpha(p)=\pi_{p}(\alpha)$, where $\pi_{p}$ is the homomorphism $A \rightarrow \mathbb{C}$ that corresponds to $p$. A rational function $f=a / s$ with $a$ and $s$ in $A$ is an element of $A_{s}$, and it defines a function on the open subset $X_{s}$. A rational function $f$ is regular at a point $p$ of $X$ if it can be written as a fraction $a / s$ such that $s(p) \neq 0$. A rational function is regular on $X$ if it is regular at every point of $X$.
2.7.2. Proposition. The regular functions on an affine variety $X=\operatorname{Spec} A$ are the elements of the coordinate algebra $A$.
proof. Let $f$ be a rational function that is regular on $X$. So for every point $p$ of $X$, there is a localization $X_{s}=\operatorname{Spec} A_{s}$ that contains $p$, such that $f$ is an element of $A_{s}$. Because $X$ is quasicompact, a finite set of these localizations, say $X_{s_{1}}, \ldots, X_{s_{k}}$, will cover $X$. Then $s_{1}, \ldots, s_{k}$ have no common zeros on $X$, so they generate the unit ideal of $A$ 2.5.8. Since $f$ is in $A_{s_{i}}$, we can write $f=s_{i}^{-n} b_{i}$, or $s_{i}^{n} f=b_{i}$, with $b_{i}$ in $A$, and we can use the same exponent $n$ for each $i$. Since the elements $s_{i}$ generate the unit ideal of $A$, so do the powers $s_{i}^{n}$. Say that $\sum s_{i}^{n} c_{i}=1$, with $c_{i}$ in $A$. Then $f=\sum s_{i}^{n} c_{i} f=\sum c_{i} b_{i}$ is an element of $A$.

## morphisms

Let $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$ be affine varieties, and let $A \xrightarrow{\varphi} B$ be an algebra homomorphism. A point $q$ of $Y$ corresponds to an algebra homomorphism $B \xrightarrow{\pi_{q}} \mathbb{C}$. When we compose $\pi_{q}$ with $\varphi$, we obtain a homomorphism $A \xrightarrow{\pi_{q} \varphi} \mathbb{C}$. By definition, the points of Spec $A$ correspond to homomorphisms $A \xrightarrow{\pi_{p}} \mathbb{C}$. So there is a unique point $p$ of $X=\operatorname{Spec} A$ such that $\pi_{p}=\pi_{q} \varphi$ :

2.7.5. Definition. Let $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$. A morphism $Y \xrightarrow{u} X$ is a map that can be defined by an algebra homomorphism $A \xrightarrow{\varphi} B$, as follows: If $q$ is a point of $Y$, then $u q$ is the point $p$ of $X$ such that $\pi_{p}=\pi_{q} \varphi$.

Let $\alpha$ be an element of $A$ and let $\beta=\varphi(\alpha)$. If $p$ is the image in $X$ of a point $q$ of $Y$, then $\beta(q)=\alpha(p)$ :

$$
\begin{equation*}
\beta(q)=\pi_{q}(\beta)=\pi_{q}(\varphi \alpha)=\pi_{p}(\alpha)=\alpha(p) \tag{2.7.6}
\end{equation*}
$$

The morphism $Y \xrightarrow{u} X$ is an isomorphism if and only if there is an inverse morphism. This will be true if and only if $A \xrightarrow{\varphi} B$ is an isomorphism of algebras.

The next formula sums up the relation between a homomorphism $A \xrightarrow{\varphi} B$ and the associated morphism $Y \xrightarrow{u} X$. If $q$ is a point of $Y$ and $\alpha$ is an element of $A$, then

$$
\begin{equation*}
\alpha[u(q)]=[\varphi \alpha](q) \tag{2.7.7}
\end{equation*}
$$

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Thus the homomorphism $\varphi$ is determined by the map $u$, and vice-versa. But just as most maps $A \rightarrow B$ aren't homomorphisms, most maps $Y \rightarrow X$ aren't morphisms.

The description of a morphism can be confusing because the direction of the arrow is reversed. It will become clearer as we expand the discussion.

## Morphisms to the affine line.

A morphism $Y \xrightarrow{u} \mathbb{A}^{1}$ from a variety $Y=\operatorname{Spec} B$ to the affine line $\operatorname{Spec} \mathbb{C}[x]$ is defined by an algebra homomorphism $\mathbb{C}[x] \xrightarrow{\varphi} B$, which substitutes an element $\beta$ of $B$ for $x$. The corresponding morphism $u$ sends a point $q$ of $Y$ to the point of the $x$-line at which $x=\beta(q)$.

For example, let $Y$ be the space of $2 \times 2$ matrices, so that $B=\mathbb{C}\left[y_{i j}\right], 1 \leq i, j \leq 2$, and let $\mathbb{A}^{1}$ be the affine line Spec $\mathbb{C}[x]$. The determinant defines a morphism $Y \rightarrow \mathbb{A}^{1}$. The corresponding algebra homomorphism $\mathbb{C}[x] \xrightarrow{\varphi} \mathbb{C}\left[y_{i j}\right]$ substitutes $y_{11} y_{22}-y_{12} y_{21}$ for $x$. It sends a polynomial $f(x)$ to $f\left(y_{11} y_{22}-y_{12} y_{21}\right)$.

In the other direction, a morphism from $\mathbb{A}^{1}$ to a variety $Y$ may be called a (complex) polynomial path in $Y$. For example, if $Y$ is the space of matrices, a morphism $\mathbb{A}^{1} \rightarrow Y$ corresponds to a homomorphism $\mathbb{C}\left[y_{i j}\right] \rightarrow \mathbb{C}[x]$, and such a homomorphism substitutes polynomials in $x$ for the variables $y_{i j}$.

## Morphisms to affine space.

A morphism from an affine variety $Y=\operatorname{Spec} B$ to affine space $\mathbb{A}^{n}$ will be defined by a homomorphism $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] \xrightarrow{\Phi} B$, which substitutes elements $\beta_{i}$ of $B$ for $x_{i}: \Phi(f(x))=f(\beta)$. (We are using an upper case $\Phi$ here, in order to keep $\varphi$ in reserve.) The corresponding morphism $Y \xrightarrow{u} \mathbb{A}^{n}$ sends a point $q$ of $Y$ to the point $\left(\beta_{1}(q), \ldots, \beta_{n}(q)\right)$ of $\mathbb{A}^{n}$.

## Morphisms to affine varieties.

Let $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$ be affine varieties. Say that we have chosen a presentation $A=$ $\mathbb{C}\left[x_{1}, \ldots, x_{m}\right] /\left(f_{1}, \ldots, f_{k}\right)$ of $A$, so that $X$ becomes the closed subvariety $V(f)$ of affine space $\mathbb{A}^{m}$. There is no need to choose a presentation of $B$. A natural way to define a morphism from a variety $Y$ to $X$ is as a morphism $Y \xrightarrow{u} \mathbb{A}^{m}$ to affine space, whose image is contained in $X$. We check that this agrees with Definition 2.7.5.

As above, a morphism $Y \xrightarrow{u} \mathbb{A}^{m}$ corresponds to a homomorphism $\mathbb{C}\left[x_{1}, \ldots, x_{m}\right] \xrightarrow{\Phi} B$. It will be determined by the set $\left(\beta_{1}, \ldots, \beta_{m}\right)$ of elements of $B$ such that $\beta_{i}=\Phi\left(x_{i}\right)$. Since $X$ is the locus of zeros of the polynomials $f$, the image of $Y$ will be contained in $X$ if and only if $f_{i}\left(\beta_{1}(q), \ldots, \beta_{m}(q)\right)=0$ for every point $q$ of $Y$ and every $i$, i.e., if and only if $f_{i}(\beta)$ is in every maximal ideal of $B$, in which case $f_{i}(\beta)=0$ for every $i$ 2.5.16(i). Another way to say this is:

The image of $Y$ is contained in $X$ if and only if $\beta=\left(\beta_{1}, \ldots, \beta_{m}\right)$ solves the equations $f(x)=0$.
And, if $\beta$ is a solution, the map $\Phi$ defines a map $A \xrightarrow{\varphi} B$.


There is an elementary, but important, principle at work here:

- Homomorphisms from an algebra $A=\mathbb{C}[x] /(f)$ to an algebra $B$ correspond to solutions of the equations $f=0$ in $B$.
mor-phandho-
2.7.8. Example. Let $B=\mathbb{C}[x]$, and let $A$ be the coordinate algebra $\mathbb{C}[u, v] /\left(v^{2}-u^{3}\right)$ of a cusp curve. A homomorphism $A \rightarrow B$ is determined by a solution of the equation $v^{2}=u^{3}$ in $b b c[x]$. For example, $u=x^{3}$ and $v=x^{3}$ is a solution. Every solution will have the form $u=f^{3}, v=f^{2}$ with $f$ in $\mathbb{C}[x]$.
2.7.9. Corollary. Let $X=\operatorname{Spec} A$ and let $Y=\operatorname{Spec} B$ be affine varieties. Suppose that $A$ is presented as the quotient $\mathbb{C}\left[x_{1}, \ldots, x_{m}\right] /\left(f_{1}, \ldots, f_{k}\right)$ of a polynomial ring. There are bijective correspondences between the following sets:
- algebra homomorphisms $A \rightarrow B$, or morphisms $Y \rightarrow X$,
- morphisms $Y \rightarrow \mathbb{A}^{n}$ whose images are contained in $X$,
- solutions of the equations $f_{i}(x)=0$ in $B$,

The second and third sets refer to an embedding of the variety $X$ into affine space, but the first one does not. It shows that a morphism depends only on the varieties $X$ and $Y$, not on an embedding of $X$,

The geometry of a morphism will be described more completely in Chapters ?? and ??. We note a few more facts about them here.
2.7.10. Proposition. Let $X \stackrel{u}{\longleftarrow} Y$ be the morphism of affine varieties that corresponds to a homomorphism of coordinate algebras $A \xrightarrow{\varphi} B$.
(i) Let $Y \stackrel{v}{\longleftarrow} Z$ be another morphism, corresponding to another homomorphism $B \xrightarrow{\psi} R$ of finite-type domains. The the composition $Z \xrightarrow{u v} X$. is the morphism that corresponds to the composed homomorphism $A \xrightarrow{\psi \varphi} R$.
(ii) Suppose that $B=A / P$, where $P$ is a prime ideal of $A$, and that $\varphi$ is the canonical homomorphism $A \rightarrow A / P$. Then $u$ is the inclusion of the closed subvariety $Y=V_{X}(P)$ into $X$.
(iii) $\varphi$ is surjective if and only if u maps $Y$ isomorphically to a closed subvariety of $X$.

It is useful to rephrase the definition of the morphism $Y \xrightarrow{u} X$ that corresponds to a homomorphism $A \xrightarrow{\varphi} B$ in terms of maximal ideals. Let $\mathfrak{m}_{q}$ be the maximal ideal of $B$ at a point $q$ of $Y$. The inverse image of $\mathfrak{m}_{q}$ in $A$ is the kernel of the composed homomorphism $A \xrightarrow{\varphi} B \xrightarrow{\pi_{q}} \mathbb{C}$, so it is a maximal ideal of $A$ : $\varphi^{-1} \mathfrak{m}_{q}=\mathfrak{m}_{p}$ for some point $p$ of $X$. That point is the image of $q: p=u q$.

In the other direction, let $\mathfrak{m}_{p}$ be the maximal ideal at a point $p$ of $X$, and let $J$ be the ideal generated by the image of $\mathfrak{m}_{p}$ in $B$. This ideal is called the extension of $\mathfrak{m}_{p}$ to $B$. Its elements are finite sums $\sum \varphi\left(z_{i}\right) b_{i}$ with $z_{i}$ in $\mathfrak{m}_{p}$ and $b_{i}$ in $B$. If $q$ is is a point of $Y$, then $u q=p$ if and only if $\mathfrak{m}_{p}=\varphi^{-1} \mathfrak{m}_{q}$, and this will be true if and only if $J$ is contained in $\mathfrak{m}_{q}$.

Recall that, if $Y \xrightarrow{u} X$ is a map of sets, the fibre of $Y$ over a point $p$ of $X$ is the set of points $q$ of $Y$ that map to $p$.
2.7.11. Corollary. Let $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$, and let $Y \xrightarrow{u} X$ be the morphism corresponding to a homomorphism $A \xrightarrow{\varphi} B$. Let $\mathfrak{m}_{p}$ be the maximal ideal at a point $p$ of $X$, and let $J=\mathfrak{m}_{p} B$ be the extended ideal.
(i) The fibre of $Y$ over $p$ is the set $V_{Y}(J)$ of points $q$ such that $J \subset \mathfrak{m}_{q}$.
(ii) The fibre of $Y$ over $p$ is empty if and only if $J$ is the unit ideal of $B$.

### 2.7.12. Example. (blowing up the plane)

Let $Z$ and $Y$ be the affine planes with coordinates $x, z$ and $x, y$, respectively. The map $Z \xrightarrow{\pi} Y$ defined by $y=x z$. The morphism that corresponds to the algebra homomorphism $\mathbb{C}[x, y] \xrightarrow{\varphi} \mathbb{C}[x, z]$ defined by $\varphi(x)=x, \varphi(y)=x z$.

This morphism $\pi$ is bijective at points $(x, y) \mathrm{t}$ which $x \neq 0$. At such a point, $y=z x^{-1}$. The fibre of $Z$ over a point of $Y$ of the form $(0, y)$ is empty unless $y=0$, and the fibre over the origin $(0,0)$ in $Y$ is the $z$-axis $\{(0, z)\}$ in the plane $Z$. Because the origin in $Y$ is replaced by a line in $Z$, this morphism $\pi$ is called a blowup of the affine plane $Y$.

## figure

2.7.13. Proposition. A morphism $Y \xrightarrow{u} X$ of affine varieties is a continuous map in the Zariski topology and also in the classical topology.
proof. The Zariski topology: Let $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$, so that $u$ corresponds to an algebra homomorphism $A \xrightarrow{\varphi} B$. A closed subset $C$ of $X$ will be the zero locus of a set $\alpha=\left\{\alpha_{1}, \ldots, \alpha_{k}\right\}$ of elements of $A$. Let $\beta_{i}=\varphi \alpha_{i}$. The inverse image $u^{-1} C$ is the set of points $q$ such that $p=u q$ is in $C$, i.e., such that $\alpha_{i}(u q)=\beta_{i}(q)=02.7 .5$. So $u^{-1} C$ is the zero locus in $Y$ of the elements $\beta_{i}=\varphi\left(\alpha_{i}\right)$ of $B$. It is a closed set.

The classical topology: We use the fact that polynomials are continuous functions. First, a morphism if affine spaces $\mathbb{A}_{y}^{n} \xrightarrow{\widetilde{u}} \mathbb{A}_{x}^{m}$ is defined by an algebra homomorphism $\mathbb{C}\left[x_{1}, \ldots, x_{m}\right] \xrightarrow{\Phi} \mathbb{C}\left[y_{1}, \ldots, y_{n}\right]$, and this
homomorphism is determined by the polynomials $h_{1}(y), \ldots, h_{m}(y)$ that are the images of the variables $x$. The morphism $\widetilde{u}$ sends the point $\left(y_{1}, \ldots, y_{n}\right)$ of $\mathbb{A}^{n}$ to the point $\left(h_{1}(y), \ldots, h_{m}(y)\right)$ of $\mathbb{A}^{m}$. It is continuous.

Next, a morphism $Y \xrightarrow{u} X$ is defined by a homomorphism $A \xrightarrow{\varphi} B$. We choose presentations $A=\mathbb{C}[x] / I$ and $B=\mathbb{C}[y] / J$, and we form a diagram of homomorphisms and the associated diagram of morphisms:


Here $\alpha$ and $\beta$ are the canonical maps of a ring to a quotient ring. The map $\alpha$ sends $x_{1}, \ldots, x_{n}$ to $\alpha_{1}, \ldots, \alpha_{n}$. Then $\Phi$ is obtained by choosing elements $h_{i}$ of $\mathbb{C}[y]$, such that $\beta\left(h_{i}\right)=\varphi\left(\alpha_{i}\right)$.

In the diagram of morphisms, $\widetilde{u}$ is continuous, and the vertical arrows are the embeddings of $X$ and $Y$ into their affine spaces. Since the topologies on $X$ and $Y$ are induced from their embeddings, $u$ is continuous.

As we see here, every morphism of affine varieties can be obtained by restriction from a morphism of affine spaces. However, in the diagram above, the morphism $\widetilde{u}$ isn't unique. It depends on the choice of the polynomials $h_{i}$.

### 2.8 Finite Group Actions

Let $G$ be a finite group of automorphisms of a finite-type domain $B$. An invariant element of $B$ is an element that is sent to itself by every element $\sigma$ of $G$. For example, the product and the sum

$$
\begin{equation*}
\prod_{\sigma \in G} \sigma b \quad, \quad \sum_{\sigma \in G} \sigma b \tag{2.8.1}
\end{equation*}
$$

are invariant elements. The invariant elements form a subalgebra of $B$ that is often denoted by $B^{G}$. Theorem 2.8 .5 below asserts that $B^{G}$ is a finite-type algebra, and that points of $\operatorname{Spec} B^{G}$ correspond bijectively to $G$-orbits in Spec $B$.

### 2.8.2. Examples.

(i) The symmetric group $G=S_{n}$ operates on the polynomial ring $R=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ by permuting the variables, and the Symmetric Functions Theorem asserts that the elementary symmetric functions

$$
s_{1}=\sum_{i} x_{i}, \quad s_{2}=\sum_{i<j} x_{i} x_{j}, \ldots, s_{n}=x_{1} x_{2} \cdots x_{n}
$$

generate the ring $R^{G}$ of invariant polynomials. Moreover, $s_{1}, \ldots, s_{n}$ are algebraically independent, so $R^{G}$ is the polynomial algebra $\mathbb{C}\left[s_{1}, \ldots, s_{n}\right]$. The inclusion of $R^{G}$ into $R$ gives us a morphism from affine $x$-space $\mathbb{A}_{x}^{n}$ to affine $s$-space $\mathbb{A}_{s}^{n}=\operatorname{Spec} R^{G}$. If $a=\left(a_{1}, \ldots, a_{n}\right)$ is a point of $\mathbb{A}_{s}^{n}$, the points $b=\left(b_{1}, \ldots, b_{n}\right)$ of $\mathbb{A}_{x}^{n}$ that map to $a$ are those such that $s_{i}(b)=a_{i}$. They are the roots of the polynomial $x^{n}-a_{1} x^{n-1}+\cdots \pm a_{n}$. Since the roots form a $G$-orbit, the set of $G$-orbits of points of $\mathbb{A}_{x}^{n}$ maps bijectively to $\mathbb{A}_{s}^{n}$.
(ii) Let $\zeta=e^{2 \pi i / n}$, let $\sigma$ be the automorphism of the polynomial ring $B=\mathbb{C}\left[y_{1}, y_{2}\right]$ defined by $\sigma y_{1}=\zeta y_{1}$ and $\sigma y_{2}=\zeta^{-1} y_{2}$. Let $G$ be the cyclic group of order $n$ generated by $\sigma$, and let $A$ denote the algebra $B^{G}$ of invariant elements. A monomial $m=y_{1}^{i} y_{2}^{j}$ is invariant if and only if $n$ divides $i-j$, and an invariant polynomial is a linear combination of invariant monomials. You will be able to show that the three monomials

$$
\begin{equation*}
u_{1}=y_{1}^{n}, u_{2}=y_{2}^{n}, \text { and } w=y_{1} y_{2} \tag{2.8.3}
\end{equation*}
$$

generate $A$. We'll use the same symbols $u_{1}, u_{2}, w$ to denote variables in the polynomial ring $\mathbb{C}\left[u_{1}, u_{2}, w\right]$. Let $J$ be the kernel of the canonical homomorphism $\mathbb{C}\left[u_{1}, u_{2}, w\right] \xrightarrow{\tau} A$ that sends $u_{1}, u_{2}, w$ to $y_{1}^{n}, y_{2}^{n}, y_{1} y_{2}$.
2.8.4. Lemma. With notation as above, the kernel $J$ of $\tau$ is the principal ideal of $\mathbb{C}\left[u_{1}, u_{2}, w\right]$ generated by the polynomial $f=w^{n}-u_{1} u_{2}$.
proof. First, $f$ is an element of $J$. Next, let $g\left(u_{1}, u_{2}, w\right)$ be an element of $J$. So $g\left(y_{1}^{n}, y_{2}^{n}, y_{1} y_{2}\right)=0$. We divide $g$ by $f$, considered as a monic polynomial in $w$, say $g=f q+r$, where the remainder $r$ has degree $<n$ in $w$. The remainder will be in $J$ too: $r\left(y_{1}^{n}, y_{2}^{n}, y_{1} y_{2}\right)=0$. We write $r\left(u_{1}, u_{2}, w\right)$ as a polynomial of degree at most $n$ in $w$ : $\quad r=r_{0}\left(u_{1}, u_{2}\right)+r_{1}\left(u_{1}, u_{2}\right) w+\cdots+r_{n-1}\left(u_{1}, u_{2}\right) w^{n-1}$. When we substitute $y_{1}^{n}, y_{2}^{n}, y_{1} y_{2}$, the term $r_{i}\left(u_{1}, u_{2}\right) w^{i}$ becomes $r_{i}\left(y_{1}^{n}, y_{2}^{n}\right)\left(y_{1} y_{2}\right)^{i}$. The degree in $y_{1}$ of every monomial that appears here will be congruent to $i$ modulo $n$, and the same is true for $y_{2}$. Since $r\left(y_{1}^{n}, y_{2}^{n}, y_{1} y_{2}\right)=0$, and since the indices $i$ are distinct, $r_{i}\left(y_{1}^{n}, y_{2}^{n}\right)$ will be zero for every $i$. And if $r_{i}\left(y_{1}^{n}, y_{2}^{n}\right)$ is zero, then $r_{i}\left(u_{1}, u_{2}\right)=0$. So $r=0$, which means that $f$ divides $g$.

We go back to the operation of the cyclic group on $B$ and the algebra of invariants $A$. Let $Y$ denote the affine plane $\operatorname{Spec} B$, and let $X=\operatorname{Spec} A$. The group $G$ operates on $Y$, and except for the origin, which is a fixed point, the orbit of a point $\left(y_{1}, y_{2}\right)$ consists of the $n$ points $\left(\zeta^{i} y_{1}, \zeta^{-i} y_{2}\right), i=0, \ldots, n-1$. To show that $G$-orbits in $Y$ correspond bijectively to points of $X$, we fix complex numbers $u_{1}, u_{2}, w$ with $w^{n}=u_{1} u_{2}$, and we look for solutions of the equations 2.8.3. When $u_{1} \neq 0$, the equation $u_{1}=y_{1}^{n}$ has $n$ solutions for $y_{1}$, and then $y_{2}$ is determined by the equation $w=y_{1} y_{2}$. So the fibre has order $n$. Similarly, there are $n$ points in the fibre if $u_{2} \neq 0$. If $u_{1}=u_{2}=0$, then $y_{1}=y_{2}=w=0$. In all cases, the fibres are the $G$-orbits.
2.8.5. Theorem. Let $G$ be a finite group of automorphisms of a finite-type domain $B$, and let $A$ denote the algebra $B^{G}$ of invariant elements. Let $Y=\operatorname{Spec} B$ and $X=\operatorname{Spec} A$.
(i) $A$ is a finite-type domain and $B$ is a finite $A$-module.
(ii) $G$ operates by automorphisms on $Y$.
(iii) The morphism $Y \rightarrow X$ defined by the inclusion $A \subset B$ is surjective, and its fibres are the $G$-orbits of points of $Y$.

When a group $G$ operates on a set $Y$, one often denotes the set of $G$-orbits of $Y$ by $Y / G(\operatorname{read} Y \bmod G)$. With that notation, the theorem asserts that there is a bijective map

$$
Y / G \rightarrow X
$$

proof of 2.8 .5 (i): The invariant algebra $A=B^{G}$ is a finite-type algebra, and $B$ is a finite $A$-module.
This is an interesting indirect proof. To show that $A$ is a finite-type algebra, one constructs a finite-type subalgebra $R$ of $A$ such that $B$ is a finite $R$-module.

Let $\left\{z_{1}, \ldots, z_{k}\right\}$ be the $G$-orbit of an element $z_{1}$ of $B$. The orbit is the set of roots of the polynomial

$$
f(t)=\left(t-z_{1}\right) \cdots\left(t-z_{k}\right)=t^{k}-s_{1} t^{k-1}+\cdots \pm s_{k}
$$

spoly
whose coefficients are the elementary symmetric functions in $\left\{z_{1}, \ldots, z_{k}\right\}$. Let $R_{1}$ denote the algebra generated by those symmetric functions. Because the symmetric functions are invariant, $R_{1} \subset A$. Using the equation $f\left(z_{1}\right)=0$, we can write any power of $z_{1}$ as a polynomial in $z_{1}$ of degree less than $k$, with coefficients in $R_{1}$.

We choose a finite set of generators $\left\{y_{1}, \ldots, y_{r}\right\}$ for the algebra $B$. If the order of the orbit of $y_{j}$ is $k_{j}$, then $y_{j}$ will be the root of a monic polynomial $f_{j}$ of degree $k_{j}$ with coefficients in $A$. Let $R$ denote the finite-type algebra generated by all of the coefficients of all of the polynomials $f_{1}, \ldots, f_{r}$. We can write any power of $y_{j}$ as a polynomial in $y_{j}$ with coefficients in $R$, and of degree less than $k_{j}$. Using such expressions, we can write every monomial in $y_{1}, \ldots, y_{r}$ as a polynomial with coefficients in $R$, whose degree in each variable $y_{j}$ is less than $k_{j}$. Since $y_{1}, \ldots, y_{r}$ generate $B$, we can write every element of $B$ as such a polynomial. Then the finite set of monomials $y_{1}^{e_{1}} \cdots y_{r}^{e_{r}}$ with $e_{j}<k_{j}$ spans $B$ as an $R$-module. Therefore $B$ is a finite $R$-module.

Since $R$ is a finite-type algebra, it is noetherian. The algebra $A$ of invariants is a subalgebra of $B$ that contains $R$. So when regarded as an $R$-module, $A$ is a submodule of the finite $R$-module $B$. Since $R$ is noetherian, $A$ is also a finite $R$-module. When we put a finite set of algebra generators for $R$ together with a finite set of $R$-module generators for $A$, we obtain a finite set of algebra generators for $A$. So $A$ is a finite-type algebra. And, since $B$ is a finite $R$-module, it is also a finite module over the larger ring $A$.

## proof of 2.8 .5 (ii): The group $G$ operates on $Y$.

A group element $\sigma$ is a homomorphism $B \xrightarrow{\sigma} B$, which defines a morphism $Y \stackrel{u_{\sigma}}{\longleftrightarrow} Y$, as in Definition 2.7.5. Since $\sigma$ is an invertible homomorphism, i.e., an automorphism, $u_{\sigma}$ is also an automorphism. Thus $G$ operates on $Y$. However, there is a point that should be mentioned.

Let's write the operation of $G$ on $B$ on the left as usual, so that a group element $\sigma$ maps an element $\beta$ of $B$ to $\sigma b$. Then if $\sigma$ and $\tau$ are two group elements, the product $\sigma \tau$ acts as first do $\tau$, then $\sigma$ : $\quad(\sigma \tau) \beta=\sigma(\tau \beta)$.

$$
\begin{equation*}
B \xrightarrow{\tau} B \xrightarrow{\sigma} B \tag{2.8.6}
\end{equation*}
$$

We substitute $u=u_{\sigma}$ into Definition 2.7.5. If $q$ is a point of $Y$, the morphism $Y \stackrel{u_{\sigma}}{\leftarrow} Y$ sends $q$ to the point $p$ such that $\pi_{p}=\pi_{q} \sigma$. It seems permissible to drop the symbol $u$, and to write the morphism simply as $Y \stackrel{\sigma}{\longleftarrow} Y$. But since arrows are reversed when going from homomorphisms of algebras to morphisms of their spectra, the maps displayed in 2.8.6, give us morphisms

$$
\begin{equation*}
Y \stackrel{\tau}{\leftarrow} Y \stackrel{\sigma}{\longleftarrow} Y \tag{2.8.7}
\end{equation*}
$$

On $Y=\operatorname{Spec} B$, the product $\sigma \tau$ acts as first do $\sigma$, then $\tau$.
To get around this problem, we can put the symbol $\sigma$ on the right when it operates on $Y$, so that $\sigma$ sends a point $q$ to $q \sigma$. Then if $q$ is a point of $Y$, we will have $q(\sigma \tau)=(q \sigma) \tau$, as required of an operation.

- If $G$ operates on the left on $B$, then it operates on the right on $\operatorname{Spec} B$.

This is important only when one wants to compose morphisms. In Definition 2.7.5, we followed custom and wrote the morphism $u$ that corresponds to an algebra homomorphism $\varphi$ on the left. We will continue to write morphisms on the left when possible, but not here.

Let $\beta$ be an element of $B$ and let $q$ be a point of $Y$. The value of the function $\sigma \beta$ at a point $q$ is the same as the value of the function $\beta$ at the point $q \sigma:[\sigma \beta](q) \stackrel{\text { def } n}{=} \pi_{q}(\sigma \beta)=\pi_{q \sigma}(\beta) \stackrel{\text { defn }}{=} \beta(q \sigma)$ 2.7.6.:

$$
\begin{equation*}
[\sigma \beta](q)=\beta(q \sigma) \tag{2.8.8}
\end{equation*}
$$

proof of 2.8 .5 (iii): The fibres of the morphism $Y \rightarrow X$ are the $G$-orbits in $Y$.
We go back to the subalgebra $A=B^{G}$. For $\sigma$ in $G$, we have a diagram of algebra homomorphisms and the corresponding diagram of morphisms


The diagram of morphisms shows that the elements of $Y$ forming a $G$-orbit have the same image in $X$, and therefore that the set of $G$-orbits in $Y$, which we may denote by $Y / G$, maps to $X$. We show that the map $Y / G \rightarrow X$ is bijective.
2.8.10. Lemma. (i) Let $p_{1}, \ldots, p_{k}$ be distinct points of affine space $\mathbb{A}^{n}$, and let $c_{1}, \ldots, c_{k}$ be complex numbers. There is a polynomial $f\left(x_{1}, \ldots, x_{n}\right)$ such that $f\left(p_{i}\right)=c_{i}$ for $i=1, \ldots, n$.
(ii) Let $B$ be a finite-type algebra, let $q_{1}, \ldots, q_{k}$ be points of $\operatorname{Spec} B$, and let $c_{1}, \ldots, c_{k}$ be complex numbers. There is an element $\beta$ in $B$ such that $\beta\left(q_{i}\right)=c_{i}$ for $i=1, \ldots, k$.
injectivity of the map $Y / G \rightarrow X$ :
Let $O_{1}$ and $O_{2}$ be distinct $G$-orbits. Lemma 2.8 .10 tells us that there is an element $\beta$ in $B$ whose value is 0 at every point of $O_{1}$, and is 1 at every point of $O_{2}$. Since $G$ permutes the orbits, $\sigma \beta$ will also be 0 at points of $O_{1}$ and 1 at points of $O_{2}$. Then the product $\gamma=\prod_{\sigma} \sigma \beta$ will be 0 at points of $O_{1}$ and 1 at points of $O_{2}$, and the product $\gamma$ is invariant. If $p_{i}$ denotes the image in $X$ of the orbit $O_{i}$, the maximal ideal $\mathfrak{m}_{p_{i}}$ of $A$ is the intersection $A \cap \mathfrak{m}_{q}$, where $q$ is any point in $O_{i}$. Therefore $\gamma$ is in the maximal ideal $\mathfrak{m}_{p_{1}}$, but not in $\mathfrak{m}_{p_{2}}$. The images of the two orbits are distinct.
surjectivity of the map $Y / G \rightarrow X$ :
It suffices to show that the map $Y \rightarrow X$ is surjective.
2.8.11. Lemma. If I is an ideal of the invariant algebra $A$, and if the extended ideal $I B$ is the unit ideal of $B$, then $I$ is the unit ideal of $A$.

As before, the extended ideal $I B$ is the ideal of $B$ generated by $I$.
Let's assume the lemma for the moment, and use it to prove surjectivity of the map $Y \rightarrow X$. Let $p$ be a point of $X$. The lemma tells us that the extended ideal $\mathfrak{m}_{p} B$ isn't the unit ideal. So it is contained in a maximal ideal $\mathfrak{m}_{q}$ of $B$, where $q$ is a point of $Y$. Then $\mathfrak{m}_{p} \subset\left(\mathfrak{m}_{p} B\right) \cap A \subset \mathfrak{m}_{q} \cap A$.

The contraction $\mathfrak{m}_{q} \cap A$ is an ideal of $A$, and it isn't the unit ideal because 1 isn't in $\mathfrak{m}_{q}$. Since $\mathfrak{m}_{p}$ is a maximal ideal, $\mathfrak{m}_{p}=\mathfrak{m}_{q} \cap A$. This means that $q$ maps to $p$ in $X$.
proof of the lemma. If $I B=B$, there will be an equation $\sum_{i} z_{i} b_{i}=1$, with $z_{i}$ in $I$ and $b_{i}$ in $B$. The sums $\alpha_{i}=\sum_{\sigma} \sigma b_{i}$ are invariant, so they are elements of $A$, and the elements $z_{i}$ are invariant. Therefore $\sum_{\sigma} \sigma\left(z_{i} b_{i}\right)=z_{i} \sum_{\sigma} \sigma b_{i}=z_{i} \alpha_{i}$ is in $I$. Then

$$
\sum_{\sigma} 1=\sum_{\sigma} \sigma(1)=\sum_{\sigma, i} \sigma\left(z_{i} b_{i}\right)=\sum_{i} z_{i} \alpha_{i}
$$

The right side is in $I$, and the left side is the order of the group which, because $A$ contains the complex numbers, is an invertible element of $A$. So $I$ is the unit ideal.

# Chapter 3 PROJECTIVE ALGEBRAIC GEOMETRY 

3.1 Projective Varieties
3.2 Homogeneous Ideals
3.3 Product Varieties
3.4 Morphisms and Isomorphisms
3.5 Affine Varieties
3.6 Lines in Projective Three-Space

### 3.1 Projective Varieties

pvariety
xlambdax
pointclosed
definepoint
closedin-
line

The projective space $\mathbb{P}^{n}$ of dimension $n$ was described in Chapter 1 Its points are equivalence classes of nonzero vectors $\left(x_{0}, \ldots, x_{n}\right)$, the equivalence relation being that, for any nonzero complex number $\lambda$,

$$
\begin{equation*}
\left(x_{0}, \ldots, x_{n}\right) \sim\left(\lambda x_{0}, \ldots, \lambda x_{n}\right) . \tag{3.1.1}
\end{equation*}
$$

A subset of $\mathbb{P}^{n}$ is Zariski closed if it is the set of common zeros of a family of homogeneous polynomials $f_{1}, \ldots, f_{k}$ in the coordinate variables $x_{0}, \ldots, x_{n}$, or if it is the set of zeros of the ideal $\mathcal{I}$ generated by such a family. Homogeneity is required because the vectors $(x)$ and $(\lambda x)$ represent the same point of $\mathbb{P}^{n}$. As explained in 1.3.1, $f(\lambda x)=0$ for all $\lambda$ if and only if $f$ is homogeneous. The Zariski closed sets are the closed sets in the Zariski topology on $\mathbb{P}^{n}$. We usually refer to the Zariski closed sets simply as closed sets.

Because the polynomial ring $\mathbb{C}\left[x_{0}, \ldots, x_{n}\right]$ is noetherian, $\mathbb{P}^{n}$ is a noetherian space: Every strictly increasing family of ideals of $\mathbb{C}[x]$ is finite, and every strictly decreasing family of closed subsets of $\mathbb{P}^{n}$ is finite. Therefore every closed subset of $\mathbb{P}^{n}$ is a finite union of irreducible closed sets 2.2 .14 . The irreducible closed sets are the projective varieties, the closed subvarieties of $\mathbb{P}^{n}$. Thus every projective variety $X$ is an irreducible closed subset of some projective space.

We will want to know when two projective varieties are isomorphic. This will be explained in Section 3.4, where morphisms are defined.

The Zariski topology on a projective variety $X$ is induced from the topology on the projective space that contains it. Since a projective variety $X$ is closed in $\mathbb{P}^{n}$, a subset of $X$ is closed in $X$ if it is closed in $\mathbb{P}^{n}$.

### 3.1.2. Lemma. The one-point sets in projective space are closed.

proof. This simple proof illustrates a general method. Let $p$ be the point $\left(a_{0}, \ldots, a_{n}\right)$. The first guess might be that the one-point set $\{p\}$ is defined by the equations $x_{i}=a_{i}$, but the polynomials $x_{i}-a_{i}$ aren't homogeneous in $x$. This is reflected in the fact that, for any $\lambda \neq 0$, the vector $\left(\lambda a_{0}, \ldots, \lambda a_{n}\right)$ represents the same point, though it won't satisfy those equations. The equations that define the set $\{p\}$ are

$$
\begin{equation*}
a_{i} x_{j}=a_{j} x_{i}, \tag{3.1.3}
\end{equation*}
$$

for $i, j=0, \ldots, n$, which imply that the ratios $a_{i} / a_{j}$ and $x_{i} / x_{j}$ are equal.
3.1.4. Lemma. The proper closed subsets of the projective line are the nonempty finite subsets, and the proper closed subsets of the projective plane are finite unions of points and curves.

Though affine varieties are important, most of algebraic geometry concerns projective varieties. It won't be clear why this is so, but one property of projective space gives a hint of its importance: With its classical topology, projective space is compact.

A topological space is compact if it has these properties:
Hausdorff property: Distinct points $p, q$ of $X$ have disjoint open neighborhoods, and quasicompactness: If $X$ is covered by a family $\left\{U^{i}\right\}$ of open sets, then a finite subfamily covers $X$.

By the way, when we say that the sets $\left\{U^{i}\right\}$ cover a topological space $X$, we mean that $X$ is the union $\bigcup U^{i}$. We don't allow $U^{i}$ to contain elements that aren't in $X$, though that would be a customary English usage.

In the classical topology, affine space $\mathbb{A}^{n}$ isn't quasicompact, and therefore it isn't compact. The HeineBorel Theorem asserts that a subset of $\mathbb{A}^{n}$ is compact in the classical topology if and only if it is closed and bounded.

We'll show that $\mathbb{P}^{n}$ is compact, assuming that the Hausdorff property has been verified. The $2 n+1$ dimensional sphere $\mathbb{S}$ of unit length vectors in $\mathbb{A}^{n+1}$ is a bounded set, and because it is the zero locus of the equation $\bar{x}_{0} x_{0}+\cdots+\bar{x}_{n} x_{n}=1$, it is closed. The Heine-Borel Theorem tells us that $\mathbb{S}$ is compact. The map $\mathbb{S} \rightarrow \mathbb{P}^{n}$ that sends a vector $\left(x_{0}, \ldots, x_{n}\right)$ to the point of projective space with that coordinate vector is continuous and surjective. The next lemma of topology shows that $\mathbb{P}^{n}$ is compact.
3.1.5. Lemma. Let $Y \xrightarrow{f} X$ be a continuous map. Suppose that $Y$ is compact and that $X$ is a Hausdorff space. Then the image $Z=f(Y)$ is a closed and compact subset of $X$.

The rest of this section contains a few examples of projective varieties.

## (3.1.6) linear subspaces

If $W$ is a subspace of dimension $r+1$ of the vector space $V$, the points of $\mathbb{P}^{n}$ that are represented by the nonzero vectors in $W$ form a linear subspace $L$ of $\mathbb{P}^{n}$, of dimension $r$. If $\left(w_{0}, \ldots, w_{r}\right)$ is a basis of $W$, the linear subspace $L$ corresponds bijectively to a projective space of dimension $r$, by

$$
c_{0} w_{0}+\cdots+c_{r} w_{r} \longleftrightarrow\left(c_{0}, \ldots, c_{r}\right)
$$

For example, the set of points $\left(x_{0}, \ldots, x_{r}, 0, \ldots, 0\right)$ is a linear subspace of dimension $r$.

## (3.1.7) a quadric surface

A quadric in $\mathbb{P}^{3}$ is the locus of zeros of an irreducible homogeneous quadratic equation in four variables.
We describe a bijective map from the product $\mathbb{P}^{1} \times \mathbb{P}^{1}$ of projective lines to a quadric. Let coordinates in the two copies of $\mathbb{P}^{1}$ be $\left(x_{0}, x_{1}\right)$ and $\left(y_{0}, y_{1}\right)$, respectively, and let the four coordinates in $\mathbb{P}^{3}$ be $w_{i j}$, with $0 \leq i, j \leq 1$. The map is defined by $w_{i j}=x_{i} y_{j}$. Its image is the quadric $Q$ whose equation is

$$
\begin{equation*}
w_{00} w_{11}=w_{01} w_{10} \tag{3.1.8}
\end{equation*}
$$

Let's check that the map $\mathbb{P}^{1} \times \mathbb{P}^{1} \rightarrow Q$ is bijective. If $w$ is a point of $Q$, one of the coordinates, say $w_{00}$, will be nonzero. Then if $(x, y)$ is a point of $\mathbb{P}^{1} \times \mathbb{P}^{1}$ whose image is $w$, so that $w_{i j}=x_{i} y_{j}$, the coordinates $x_{0}$ and $y_{0}$ must be nonzero. When we normalize $w_{00}, x_{0}$, and $y_{0}$ to $1, w_{11}=w_{01} w_{10}$. There is a unique solution for $x$ and $y$ such that $w_{i j}=x_{i} y_{j}$, namely $x_{1}=w_{10}$ and $y_{1}=w_{01}$.

The quadric with the equation $\sqrt{3.1 .8}$ contains two families of lines (one dimensional linear subspaces), the images of the subsets $x \times \mathbb{P}^{1}$ and $\mathbb{P}^{1} \times y$ of $\mathbb{P} \times \mathbb{P}$.

Note. Equation (3.1.8) can be diagonalized by the substitution $w_{00}=s+t, w_{11}=s-t$, $w_{01}=u+v$, $w_{10}=u-v$. This substitution changes the equation (3.1.8) to $s^{2}-t^{2}=u^{2}-v^{2}$. When we look at the affine open set $\{u=1\}$, the equation becomes $s^{2}+v^{2}-t^{2}=1$. The real locus of this equation is a one-sheeted hyerboloid in $\mathbb{R}^{3}$, and the two families of complex lines in the quadric correspond to the familiar rulings of this hyboloid by real lines.
segreequations

## (3.1.9) hypersurfaces

A hypersurface in projective space $\mathbb{P}^{n}$ is the locus of zeros of an irreducible homogeneous polynomial $f\left(x_{0}, \ldots, x_{n}\right)$. The degree of $Y$ is the degree of the polynomial $f$.

Plane projective curves and quadric surfaces are hypersurfaces.

## (3.1.10) the Segre embedding of a product

The product $\mathbb{P}_{x}^{m} \times \mathbb{P}_{y}^{n}$ of projective spaces can be embedded by its Segre embedding into a projective space $\mathbb{P}_{w}^{N}$ that has coordinates $w_{i j}$, with $i=0, \ldots, m$ and $j=0, \ldots, n$. So $N=(m+1)(n+1)-1$. The Segre embedding is defined by

$$
\begin{equation*}
w_{i j}=x_{i} y_{j} \tag{3.1.11}
\end{equation*}
$$

We call the coordinates $w_{i j}$ the Segre variables.
The map from $\mathbb{P}^{1} \times \mathbb{P}^{1}$ to $\mathbb{P}^{3}$ that was described in 3 , is the simplest case of a Segre embedding.
3.1.12. Proposition. The Segre embedding maps the product $\mathbb{P}^{m} \times \mathbb{P}^{n}$ bijectively to the locus $S$ of the Segre equations

$$
\begin{equation*}
w_{i j} w_{k \ell}-w_{i \ell} w_{k j}=0 \tag{3.1.13}
\end{equation*}
$$

proof. The proof is the same as the one given above, in 3.1.7. When one substitutes 3.1.11) into the Segre equations, one obtains equations in $\left\{x_{i}, y_{j}\right\}$ that are true. So the image of the Segre embedding is contained in $S$.

Say that we have a point $p$ of the locus $S$, that is the image of a point $(x, y)$ of $\mathbb{P}^{m} \times \mathbb{P}^{n}$. Some coordinate of $p$, say $w_{00}$, will be nonzero, and then $x_{0}$ and $y_{0}$ are also nonzero. We normalize $w_{00}, x_{0}$, and $y_{0}$ to 1 . Then $w_{i j}=w_{i 0} w_{0 j}$ for all $i, j$. The unique solution of the Segre equations is $x_{i}=w_{i 0}$ and $y_{j}=w_{0 j}$.

The Segre embedding is important because it makes the product of projective spaces into a projective variety, the closed subvariety of $\mathbb{P}^{N}$ defined by the Segre equations. However, to show that the product is a variety, we need to show that the locus $S$ of the Segre equations is irreducible. This is less obvious than one might expect, so we defer the discussion to Section 3.3 (see Proposition 3.3.1).

## (3.1.14) the Veronese embedding of projective space

Let the coordinates in $\mathbb{P}^{n}$ be $x_{i}$, and let those in $\mathbb{P}^{N}$ be $v_{i j}$, with $0 \leq i \leq j \leq n$. Then $N=\binom{n+2}{2}-1$. The Veronese embedding is the map $\mathbb{P}^{n} \xrightarrow{f} \mathbb{P}^{N}$ defined by $v_{i j}=x_{i} x_{j}$. The Veronese embedding resembles the Segre embedding, but in the Segre embedding, there are distinct sets of coordinates $x$ and $y$, and there is no requirement that $i \leq j$.

The proof of the next proposition is similar to the proof of 3.1.12, once one has untangled the inequalities.
3.1.15. Proposition. The Veronese embedding $f$ maps $\mathbb{P}^{n}$ bijectively to the locus $X$ in $\mathbb{P}^{N}$ of the equations

$$
v_{i j} v_{k \ell}=v_{i \ell} v_{k j} \quad \text { for } \quad 0 \leq i \leq k \leq j \leq \ell \leq n
$$

For example, the Veronese embedding maps $\mathbb{P}^{1}$ bijectively to the conic $v_{00} v_{11}=v_{01}^{2}$ in $\mathbb{P}^{2}$.

## (3.1.16) the twisted cubic

There are higher order Veronese embeddings, defined in an analogous way by the monomials of some degree $d>2$. The first example is the embedding of $\mathbb{P}^{1}$ by the cubic monomials in two variables, which maps $\mathbb{P}_{x}^{1}$ to $\mathbb{P}_{v}^{3}$. Let the coordinates in $\mathbb{P}^{3}$ be $v_{0}, \ldots, v_{3}$. The cubic Veronese embedding is defined by

$$
v_{0}=x_{0}^{3}, \quad v_{1}=x_{0}^{2} x_{1}, \quad v_{2}=x_{0} x_{1}^{2}, \quad v_{3}=x_{1}^{3}
$$

Its image is a twisted cubic in $\mathbb{P}^{3}$, the locus $\left(v_{0}, v_{1}, v_{2}, v_{3}\right)=\left(x_{0}^{3}, x_{0}^{2} x_{1}, x_{0} x_{1}^{2}, x_{1}^{3}\right)$. This locus is the set of common zeros of the three polynomials

$$
\begin{equation*}
v_{0} v_{2}-v_{1}^{2}, \quad v_{1} v_{2}-v_{0} v_{3}, \quad v_{1} v_{3}-v_{2}^{2} \tag{3.1.17}
\end{equation*}
$$

which are the $2 \times 2$ minors of the $2 \times 3$ matrix

$$
\left(\begin{array}{lll}
v_{0} & v_{1} & v_{2}  \tag{3.1.18}\\
v_{1} & v_{2} & v_{3}
\end{array}\right)
$$

A $2 \times 3$ matrix has rank $\leq 1$ if and only if its $2 \times 2$ minors are zero. So a point $\left(v_{0}, v_{1}, v_{2}, v_{3}\right)$ lies on the twisted cubic if 3.1.18 has rank one. This means that if the vectors $\left(v_{0}, v_{1}, v_{2}\right)$ and $\left(v_{1}, v_{2}, v_{3}\right)$ are both nonzero, they represent the same point of $\mathbb{P}^{2}$. Setting $x_{0}=1$ and $x_{1}=t$, the twisted cubic becomes the locus of points $\left(1, t, t^{2}, t^{3}\right)$. There is also one point on the twisted cubic at which $x_{0}=0$, the point $(0,0,0,1)$.

### 3.2 Homogeneous Ideals

We denote the polynomial algebra $\mathbb{C}\left[x_{0}, \ldots, x_{n}\right]$ by $R$ here.
3.2.1. Lemma. Let $\mathcal{I}$ be an ideal of $R$. The following conditions are equivalent.
(i) $\mathcal{I}$ can be generated by homogeneous polynomials.
(ii) A polynomial is in $\mathcal{I}$ if and only if its homogeneous parts are in $\mathcal{I}$.

An ideal $\mathcal{I}$ of $R$ that satisfies these conditions is a homogeneous ideal.

### 3.2.2. Lemma. The radical 2.2.2 of a homogeneous ideal is homogeneous.

proof. Let $\mathcal{I}$ be a homogeneous ideal, and let $f$ be an element of its radical $\operatorname{rad} \mathcal{I}$. So $f^{r}$ is in $\mathcal{I}$ for some $r$. When $f$ is written as a sum $f_{0}+\cdots+f_{d}$ of its homogeneous parts, the highest degree part of $f^{r}$ is $\left(f_{d}\right)^{r}$. Since $\mathcal{I}$ is homogeneous, $\left(f_{d}\right)^{r}$ is in $\mathcal{I}$ and $f_{d}$ is in $\operatorname{rad} \mathcal{I}$. Then $f_{0}+\cdots+f_{d-1}$ is also in rad $\mathcal{I}$. By induction on $d$, all of the homogeneous parts $f_{0}, \ldots, f_{d}$ are in $\operatorname{rad} \mathcal{I}$.

If $f$ is a set of homogeneous polynomials, its set of zeros in $\mathbb{P}^{n}$ may be denoted by $V(f)$ or $V_{\mathbb{P}^{n}}(f)$, and the set of zeros of a homogeneous ideal $\mathcal{I}$ may be denoted by $V(\mathcal{I})$ or $V_{\mathbb{P}^{n}}(\mathcal{I})$. This is the same notation as is used for closed subsets of affine space.

The complement of the origin in the affine space $\mathbb{A}^{n+1}$ is mapped to the projective space $\mathbb{P}^{n}$ by sending a vector $\left(x_{0}, \ldots, x_{n}\right)$ to the point of $\mathbb{P}^{n}$ it defines. This map can be useful when one studies projective space.

A homogeneous ideal $\mathcal{I}$ has a zero locus in projective space $\mathbb{P}^{n}$ and also a zero locus in the affine space $\mathbb{A}^{n+1}$. We can't use the $V(\mathcal{I})$ notation for both of them here, so let's denote these two loci by $V$ and $W$, respectively. Unless $\mathcal{I}$ is the unit ideal, the origin $x=0$ will be a point of $W$, and the complement of the origin will map surjectively to $V$. If a point $x$ other than the origin is in $W$, then every point of the onedimensional subspace of $\mathbb{A}^{n+1}$ spanned by $x$ is in $W$, because a homogeneous polynomial $f$ vanishes at $x$ if and only if it vanishes at $\lambda x$. An affine variety that is the union of such lines through the origin is called an affine cone. If the locus $W$ contains a point $x$ other than the origin, it is an affine cone.

The familiar locus $x_{0}^{2}+x_{1}^{2}-x_{2}^{2}=0$ is a cone in $\mathbb{A}^{3}$. The zero locus of the polynomial $x_{0}^{3}+x_{1}^{3}-x_{2}^{3}$ is also called a cone.

Note. The real locus $x_{0}^{2}+x_{1}^{2}-x_{2}^{2}=0$ in $\mathbb{R}^{3}$ decomposes into two parts when the origin is removed. Because of this, it is sometimes called a "double cone". However, the complex locus doesn't decompose.

## (3.2.3) the irrelevant ideal

In the polynomial algebra $R=\mathbb{C}\left[x_{0}, \ldots, x_{n}\right]$, the maximal ideal $\mathcal{M}=\left(x_{0}, \ldots, x_{n}\right)$ that is generated by the variables is called the irrelevant ideal because its zero locus in projective space is empty.
nozeros

## homprime

3.2.4. Proposition. The zero locus in $\mathbb{P}^{n}$ of a homogeneous ideal $\mathcal{I}$ of $R$ is empty if and only if $\mathcal{I}$ contains $a$ power of the irrelevant ideal $\mathcal{M}$.

Another way to say this is that the zero locus $V(\mathcal{I})$ in projective space of a homogeneous ideal $\mathcal{I}$ is empty if and only if either $\mathcal{I}$ is the unit ideal $R$, or $\operatorname{rad} \mathcal{I}$ is the irrelevant ideal.
proof of Proposition 3.2.4 Let $Z$ be the zero locus of $\mathcal{I}$ in $\mathbb{P}^{n}$. If $\mathcal{I}$ contains a power of $\mathcal{M}$, it contains a power of each variable. Powers of the variables have no common zeros in projective space, so $Z$ is empty.

Suppose that $Z$ is empty, and let $W$ be the locus of zeros of $\mathcal{I}$ in the affine space $\mathbb{A}^{n+1}$ with the same coordinates $x_{0}, \ldots, x_{n}$. Since the complement of the origin in $W$ maps to the empty locus $Z$, it is empty. The origin is the only point that might be in $W$. If $W$ is the one point space consisting of the origin, then $\operatorname{rad} \mathcal{I}=\mathcal{M}$. If $W$ is empty, $\mathcal{I}$ is the unit ideal.
3.2.5. Lemma. Let $\mathcal{P}$ be a homogeneous ideal in the polynomial algebra $R$, not the unit ideal. The following conditions are equivalent:
(i) $\mathcal{P}$ is a prime ideal.
(ii) If $f$ and $g$ are homogeneous polynomials, and if $f g \in \mathcal{P}$, then $f \in \mathcal{P}$ or $g \in \mathcal{P}$.
(iii) If $\mathcal{A}$ and $\mathcal{B}$ are homogeneous ideals, and if $\mathcal{A B} \subset \mathcal{P}$, then $\mathcal{A} \subset \mathcal{P}$ or $\mathcal{B} \subset \mathcal{P}$.

In other words, a homogeneous ideal is a prime ideal if the usual conditions for a prime ideal are satisfied when the polynomials or ideals are homogeneous.
proof of the lemma. When the word homogeneous is omitted, (ii) and (iii) become the definition of a prime ideal. So (i) implies (ii) and (iii). The fact that (iii) $\Rightarrow$ (ii) is proved by considering the principal ideals generated by $f$ and $g$.
(ii) $\Rightarrow$ (i) Suppose that a homogeneous ideal $\mathcal{P}$ satisfies the condition (ii), and that the product $f g$ of two polynomials, not necessarily homogeneous, is in $\mathcal{P}$. If $f$ has degree $d$ and $g$ has degree $e$, the highest degree part of $f g$ is the product $f_{d} g_{e}$ of the homogeneous parts of $f$ and $g$ of maximal degree. Since $\mathcal{P}$ is a homogeneous ideal, it contains $f_{d} g_{e}$. Therefore one of the factors, say $f_{d}$, is in $\mathcal{P}$. Let $h=f-f_{d}$. Then $h g$ is in $\mathcal{P}$, and it has lower degree than $f g$. By induction on the degree of $f g, h$ or $g$ is in $\mathcal{P}$, and if $h$ is in $\mathcal{P}$, so is $f$.
3.2.6. Proposition. Let $Y$ be the zero locus in $\mathbb{P}^{n}$ of a homogeneous radical ideal $\mathcal{I}$ that isn't the irrelevant ideal. Then $Y$ is a projective variety (an irreducible closed subset of $\mathbb{P}^{n}$ ) if and only if $\mathcal{I}$ is a prime ideal. Thus a subset $Y$ of $\mathbb{P}^{n}$ is a projective variety if and only if it is the zero locus of a homogeneous prime ideal that isn't the irrelevant ideal.
proof. The locus $W$ of zeros of $\mathcal{I}$ in the affine space $\mathbb{A}^{n+1}$ is irreducible if and only if $Y$ is irreducible. This is easy to see. Proposition 2.2.19 tells us that $W$ is irreducible if and only if the radical ideal $\mathcal{I}$ is a prime ideal.

### 3.2.7. Strong Nullstellensatz, projective version.

(i) Let $g$ be a nonconstant homogeneous polynomial in $x_{0}, \ldots, x_{n}$, and let $\mathcal{I}$ be a homogeneous ideal of $\mathbb{C}[x]$. If $g$ vanishes at every point of the zero locus $V(\mathcal{I})$ in $\mathbb{P}^{n}$, then $\mathcal{I}$ contains a power of $g$.
(ii) Let $f$ and $g$ be homogeneous polynomials. If $f$ is irreducible and if $V(f) \subset V(g)$, then $f$ divides $g$.
(iii) Let $\mathcal{I}$ and $\mathcal{J}$ be homogeneous ideals, and suppose that $\operatorname{rad} \mathcal{I}$ isn't the irrelevant ideal or the unit ideal. Then $V(\mathcal{I})=V(\mathcal{J})$ if and only if $\operatorname{rad} \mathcal{I}=\operatorname{rad} \mathcal{J}$.
proof. (i) Let $W$ be the locus of zeros of $\mathcal{I}$ in the affine space $\mathbb{A}^{n+1}$ with coordinates $x_{0}, \ldots, x_{n}$. The homogeneous polynomial $g$ vanishes at every point of $W$ different from the origin, and since $g$ isn't a constant, it vanishes at the origin too. So the affine Strong Nullstellensatz applies.

## (3.2.8) quasiprojective varieties

We will also want to consider nonempty open subsets of a projective variety. We call such a subset a variety.

For example, the complement of a point in a projective variety is a variety. An affine variety $X=\operatorname{Spec} A$ may be embedded as a closed subvariety into the standard affine space $\mathbb{U}^{0}:\left\{x_{0} \neq 0\right\}$. It becomes an open subset of its closure in $\mathbb{P}^{n}$, which is a projective variety (Lemma 2.2 .16 (ii)). And of course, a projective variety is a variety. The topology on a (quasiprojective) variety is induced from the topology on projective space.

In more usual terminology such a variety is called a quasiprojective variety, but we drop the adjective 'quasiprojective'. There are varieties that aren't quasiprojective. They cannot be embedded into any projective space. But such varieties aren't very important. We will not study them. In fact, it is hard enough to find convincing examples of such varieties that we won't try to give one here. So the adjective 'quasiprojectve' is superfluous as well as ugly.
3.2.9. Lemma. The topology on the affine open subset $\mathbb{U}^{0}: x_{0} \neq 0$ of $\mathbb{P}^{n}$ that is induced from the Zariski topology on $\mathbb{P}^{n}$ is the Zariski topology obtained by viewing $\mathbb{U}^{0}$ as the affine space $\operatorname{Spec} \mathbb{C}\left[u_{1}, \ldots, u_{n}\right], u_{i}=$ $x_{i} / x_{0}$.

### 3.3 Product Varieties

The properties of products of varieties seem intuitive, but some of the proofs aren't obvious. The reason for this is that the (Zariski) topology on a product of varieties isn't the product topology.

The product topology on the product $X \times Y$ of topological spaces is the coarsest topology such that the projection maps $X \times Y \rightarrow X$ and $X \times Y \rightarrow Y$ are continuous. If $C$ and $D$ are closed subsets of $X$ and $Y$, then $C \times D$ is a closed subset of $X \times Y$ in the product topology, and every closed set in the product topology is a finite union of subsets of the form $C \times D$.

The first examples of closed subsets of $\mathbb{P}^{m} \times \mathbb{P}^{n}$ are products of the form $C \times D$, where $C$ is a closed subset of $\mathbb{P}^{m}$ and $D$ is a closed subset of $\mathbb{P}^{n}$. But the product topology on $\mathbb{P}^{m} \times \mathbb{P}^{n}$ is much coarser than the Zariski topology. For example, the proper (Zariski) closed subsets of $\mathbb{P}^{1}$ are the nonempty finite subsets. In the product topology, the proper closed subsets of $\mathbb{P}^{1} \times \mathbb{P}^{1}$ are finite unions of points and sets of the form $\mathbb{P}^{1} \times q$, $p \times \mathbb{P}^{1}$, and $p \times q$ ('horizontal lines', 'vertical' lines, and points). Most Zariski closed subsets of $\mathbb{P}^{1} \times \mathbb{P}^{1}$ aren't of this form. The diagonal $\Delta=\left\{(p, p) \mid p \in \mathbb{P}^{1}\right\}$ is a simple example.
3.3.1. Proposition. Let $X$ and $Y$ be irreducible topological spaces, and suppose that a topology is given on the product $\Pi=X \times Y$, with the following properties:

- The projections $\Pi \xrightarrow{\pi_{1}} X$ and $\Pi \xrightarrow{\pi_{2}} Y$ are continuous.
- For all $x$ in $X$ and all $y$ in $Y$, the fibres $x \times Y$ and $X \times y$ are homeomorphic to $Y$ and $X$, respectively.

Then $\Pi$ is an irreducible topological space.
The first condition means that the topology on $X \times Y$ is at least as fine as the product topology, and the second one assures us that the topology isn't too fine. (We don't want the discrete topology on П.)

The product of varieties has the two properties mentioned in the proposition.
3.3.2. Lemma. Let $X, Y$, and $\Pi$ be as in the proposition. If $W$ is an open subset of $\Pi$, its image $U$ via the projection $\Pi \rightarrow Y$ is an open subset of $Y$.
proof. The intersection ${ }_{x} W=W \cap(x \times Y)$ is an open subset of the fibre $x \times Y$, and its image ${ }_{x} U$ in the homeomorphic space $Y$ is open too. Since $W$ is the union of the sets ${ }_{x} W, U$ is the union of the open sets ${ }_{x} U$. So $U$ is open.
proof of Proposition 3.3.1 Let $C$ and $C^{\prime}$ be closed subsets of the product $\Pi$. Suppose that $C<\Pi$ and $C^{\prime}<\Pi$, and let $W=\Pi-C$ and $W^{\prime}=\Pi-C^{\prime}$ be the open complements of $C$ and $C^{\prime}$ in $\Pi$. To show that $\Pi$ is irreducible, we must show that $C \cup C^{\prime}<\Pi$. We do this by showing that $W \cap W^{\prime}$ is not the empty set.

Since $C<\Pi, W$ isn't empty. Similarly, $W^{\prime}$ isn't empty. The lemma tells us that the images $U$ and $U^{\prime}$ of $W$ and $W^{\prime}$ via projection to $Y$ are nonempty open subsets of $Y$. Since $Y$ is irreducible, $U \cap U^{\prime}$ is nonempty. Let $y$ be a point of $U \cap U^{\prime}$. The intersection $W_{y}=W \cap(X \times y)$ is an open subset of $X \times y$, and since its image $U$ contains $y, W_{y}$ contains a point of the form $p=(x, y)$. Thus $W_{y}$ is a nonempty open subst of $W_{y}$. Similarly, $W_{y}^{\prime}=W^{\prime} \cap(X \times y)$ is nonempty. Since $X \times y$ is homeomorphic to the irreducible space $X$, it is irreducible. So $W_{y} \cap W_{y}^{\prime}$ is nonempty, and therefore $W \cap W^{\prime}$ is nonempty, as was to be shown.
toponstandaff

We inspect the product $X \times Y$ of the affine varieties $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$. Say that $X$ is embedded as a closed subvariety of $\mathbb{A}^{m}$, so that $A=\mathbb{C}\left[x_{1}, \ldots, x_{m}\right] / P$ for some prime ideal $P$, and that $Y$ is embedded similarly into $\mathbb{A}^{n}$, and $B=\mathbb{C}\left[y_{1}, \ldots, y_{n}\right] / Q$. Then in affine $x, y$-space $\mathbb{A}^{m+n}, X \times Y$ is the locus of the equations $f(x)=0$ and $g(y)=0$ with $f \in P$ and $g \in Q$. Proposition 3.3.1 shows that $X \times Y$ is irreducible. Therefore it is a variety. Let $P^{\prime}$ be the ideal of $\mathbb{C}[x, y]$ generated by the elements of $P$. It consists of sums of products of elements of $P$ with polynomials in $x, y$. Let $Q^{\prime}$ be defined analgously using $Q$, and let $I=P^{\prime}+Q^{\prime}$.
fgprime 3.3.4. Proposition. The ideal $I=P^{\prime}+Q^{\prime}$ consists all elements of $\mathbb{C}[x, y]$ that vanish on the variety $X \times Y$. Therefore I is a prime ideal.
(The fact that $X \times Y$ is a variety tells us only that the radical of $I$ is a prime ideal.)
proof of Proposition 3.3.4 Let $A=\mathbb{C}[x] / P, B=\mathbb{C}[y] / Q$, and $R=\mathbb{C}[x, y] / I$. Any polynomial in $x, y$ can the written, in many ways, as a sum, each of whose terms is a product of a polynomial in $x$ with a polynomial in $y: \quad p(x, y)=\sum a_{i}(x) b_{i}(y)$. Therefore any element $p$ of $R$ can be written as a finite sum

$$
\begin{equation*}
p=\sum_{i=1}^{k} a_{i} b_{i} \tag{3.3.5}
\end{equation*}
$$

with $a_{i}$ in $A$ and $b_{i}$ in $B$. We show that if $p$ vanishes identically on $X \times Y$, then $p=0$. To do this, we show that the same element $p$ can also be written as a sum of $k-1$ products.

Suppose that $p=0$. If $a_{k}=0$, then $p=\sum_{i=1}^{k-1} a_{i} b_{i}$. If $a_{k} \neq 0$, the function defined by $a_{k}$ isn't identically zero on $X$. We choose a point $\bar{x}$ of $X$ such that $a_{k}(\bar{x}) \neq 0$. Let $\bar{a}_{i}=a_{i}(\bar{x})$ and $\bar{p}(y)=p(\bar{x}, y)$. So $\bar{p}(y)=\sum_{i=1}^{k} \bar{a}_{i} b_{i}$. Since $p$ vanishes on $X \times Y, \bar{p}$ vanishes on $Y$, and therefore $\bar{p}=0$. Since $\bar{a}_{k} \neq 0$, we can solve the equation $\left.\sum_{i=1}^{k}\right] o a_{i} b_{i}=0$ for $b_{k}: b_{k}=\sum_{i=1}^{k-1} c_{i} b_{i}$, where $c_{i}=-\bar{a}_{i} / \bar{a}_{k}$. Substituting into $p$ gives us an expression for $p$ as a sum of $k-1$ terms. Finally, when $k=1, \quad \bar{a}_{1} b_{1}=0$. Then $b_{1}=0$, and $p=0$.

## (3.3.6) the Zariski topology on $\mathbb{P}^{m} \times \mathbb{P}^{n}$

As mentioned above 3.1.10, the product of projective spaces $\mathbb{P}^{m} \times \mathbb{P}^{n}$ is made into a projective variety by identifying it with its Segre image, the locus of the Segre equations $w_{i j} w_{k \ell}=w_{i \ell} w_{k j}$. Since $\mathbb{P}^{m} \times \mathbb{P}^{n}$, with its Segre embedding, is a projective variety, we don't really need a separate definition of its Zariski topology. Its closed subsets are the zero sets of families of homogeneous polynomials in the Segre variables $w_{i j}$ that include the Segre equations.

One can also describe the closed subsets of $\mathbb{P}^{m} \times \mathbb{P}^{n}$ directly, in terms of bihomogeneous polynomials. A polynomial $f(x, y)$ is bihomogeneous if it is homogeneous in the variables $x$ and also in the variables $y$. For example, the polynomial $x_{0}^{2} y_{0}+x_{0} x_{1} y_{1}$ is bihomogeneous, of degree 2 in $x$ and degree 1 in $y$.
Because $(x, y)$ and $(\lambda x, \mu y)$ represent the same point of $\mathbb{P}^{m} \times \mathbb{P}^{n}$ for all nonzero $\lambda$ and $\mu$, we want to know that $f(x, y)=0$ if and only if $f(\lambda x, \mu y)=0$. This is true for all nonzero $\lambda$ and $\mu$ if and only if $f$ is bihomogeneous.
3.3.7. Lemma. The (Zariski) topology on $\mathbb{P}^{m} \times \mathbb{P}^{n}$ has the properties listed in Proposition 3.3.1.

- The projections $\mathbb{P}^{m} \times \mathbb{P}^{n} \rightarrow \mathbb{P}^{m}$ and $\mathbb{P}^{m} \times \mathbb{P}^{n} \rightarrow \mathbb{P}^{n}$ are continuous maps.
- For all $y$ in $\mathbb{P}^{n}$, the fibre $\mathbb{P}^{m} \times y$, with the topology induced from $\mathbb{P}^{m} \times \mathbb{P}^{n}$, is homeomorphic to $\mathbb{P}^{m}$, and the analogous statement is true for the fibre $x \times \mathbb{P}^{n}$.
proof. We look at the projection $\mathbb{P}^{m} \times \mathbb{P}^{n} \rightarrow \mathbb{P}^{m}$. If $X$ is the closed subset of $\mathbb{P}^{m}$ defined by a system of homogeneous polynomials $f_{i}(x)$, its inverse image in $\mathbb{P}^{m} \times \mathbb{P}^{n}$ is the zero set of the same system, considered as a family of bihomogeneous polynomials of degree zero in $y$. So the inverse image is closed.

For the second property, because the projection $\mathbb{P}^{m} \times \mathbb{P}^{n} \rightarrow \mathbb{P}^{m}$ is continuous, it suffices to show that the inclusion map $\mathbb{P}^{m} \rightarrow \mathbb{P}^{m} \times \mathbb{P}^{n}$ that sends $\mathbb{P}^{m}$ to $\mathbb{P}^{m} \times y$ is continuous. If $f(x, y)$ is a bihomogeneous polynomial and $\bar{y}$ is a point of $\mathbb{P}^{n}$, the zero set of $f$ in $\mathbb{P}^{m} \times \bar{y}$ is the zero set of $f(x, \bar{y})$. This polynomial also defines a closed subset of $\mathbb{P}^{m}$.

### 3.3.8. Proposition.

(i) A subset of $\mathbb{P}^{m} \times \mathbb{P}^{n}$ is closed if and only if it is the locus of zeros of a family of bihomogeneous polynomials.
(ii) If $X$ and $Y$ are closed subsets of $\mathbb{P}^{m}$ and $\mathbb{P}^{n}$, respectively, then $X \times Y$ is a closed subset of $\mathbb{P}^{m} \times \mathbb{P}^{n}$.
proof. (i) For this proof, we denote the Segre image of $\mathbb{P}^{m} \times \mathbb{P}^{n}$ by $\Pi$. Let $f(w)$ be a homogeneous polynomial in the Segre variables $w_{i j}$. When we substitute $w_{i j}=x_{i} y_{j}$ into into $f$, we obtain a polynomial $f\left(x_{i} y_{j}\right)$ that is bihomogeneous and that has the same degree as $f$ in $x$ and in $y$. Let's denote that bihomogeneous polynomial by $\widetilde{f}(x, y)$. The inverse image of the zero set of $f$ in $\Pi$ is the zero set of $\widetilde{f}$ in $\mathbb{P}^{m} \times \mathbb{P}^{n}$. Therefore the inverse image of a closed subset of $\Pi$ is the zero set of a family of bihomogeneous polynomials in $\mathbb{P}^{m} \times \mathbb{P}^{n}$.

Conversely, let $g(x, y)$ be a bihomogeneous polynomial, say of degrees $r$ in $x$ and degree $s$ in $y$. If $r=s$, we may collect variables that appear in $g$ in pairs $x_{i} y_{j}$ and replace each pair $x_{i} y_{j}$ by $w_{i j}$. We will obtain a homogeneous polynomial $G$ in $w$ such that $G(w)=g(x, y)$ when $w_{i j}=x_{i} y_{j}$. The zero set of $G$ in $\Pi$ is the image of the zero set of $g$ in $\mathbb{P}^{m} \times \mathbb{P}^{n}$.

Suppose that $r \geq s$, and let $k=r-s$. Because the variables $y$ cannot all be zero at any point of $\mathbb{P}^{n}$, the equation $g=0$ on $\mathbb{P}^{m} \times \mathbb{P}^{n}$ is equivalent with the system of equations $g y_{0}^{k}=g y_{1}^{k}=\cdots=g y_{n}^{k}=0$. The polynomials $g y_{i}^{k}$ are bihomogeneous, of same degree in $x$ and in $y$.
(ii) A polynomial $f(x)$ can be viewed as a bihomogeneous polynomial of degree zero in $y$, and a polynomial $g(y)$ as a bihomogeneous polynomial of degree zero in $x$. So $X \times Y$, which is the locus $f=g=0$ in $\mathbb{P}^{m} \times \mathbb{P}^{n}$, is closed in $\mathbb{P}^{m} \times \mathbb{P}^{n}$.
3.3.9. Corollary. Let $X$ and $Y$ be projective varieties, and let $\Pi$ denote the product $X \times Y$. This is a closed subset of $\mathbb{P}^{m} \times \mathbb{P}^{n}$.

- The projections $\Pi \rightarrow X$ and $\Pi \rightarrow Y$ are continuous.
- For all $x$ in $X$ and all $y$ in $Y$, the fibres $x \times Y$ and $X \times y$, with topologies induced from $\Pi$, are homeomorphic to $Y$ and $X$, respectively.

Therefore the product $X \times Y$ is a projective variety.
We will come back to products in Chapter ? ?

### 3.4 Morphisms and Isomorphisms

## (3.4.1) the function field

Let $X$ be a projective variety, and let $X^{i}$ be its intersection with the standard affine open subset $\mathbb{U}^{i}$ of the projective space with coordinates $x_{0}, \ldots, x_{n}$. If it is nonempty, $X^{i}$ will be an affine variety - an irreducible closed subset of $\mathbb{U}^{i}$. Let's omit the indices for which $X^{i}$ is empty. Then the intersection $X^{i j}=X^{i} \cap X^{j}$ will be a localization of $X^{i}$ and also a localization of $X^{j}$. If $X^{i}=\operatorname{Spec} A_{i}$ and $u_{i j}=x_{j} / x_{i}$, then $X^{i j}=\operatorname{Spec} A_{i j}$, where $A_{i j}=A_{i}\left[u_{i j}^{-1}\right]=A_{j}\left[u_{j i}^{-1}\right]$. So the fraction fields of the coordinate algebras $A_{i}$ are equal for all $i$ such that $X^{i}$ isn't empty.
3.4.2. Definition. The function field $K$ of a projective variety $X$ is the field of fractions of the coordinate algebra $A_{i}$ of any one of its nonempty affine open subsets $X^{i}=X \cap \mathbb{U}^{i}$. If $X^{\prime}$ is an open subvariety of a projective variety $X$, the function field of $X^{\prime}$ is the function field of $X$.

Thus all open subvarieties of a variety have the same function field. In particular, suppose that we regard an affine variety $X=\operatorname{Spec} A$ as a closed subvariety of $\mathbb{U}^{0}$. The function field of $X$ will be the field of fractions of $A$. This agrees with the definition given in Chapter 2 (See 2.7.1).)
3.4.3. Definition. A rational function on a variety $X^{\prime}$ is an element of the function field $K$ of $X^{\prime}$.

A rational function can be evaluated at some points of $X^{\prime}$, but probably not at all points. Suppose that $X^{\prime}$ is an open subvariety of a projective variety $X$, and that $p$ is a point of $X^{\prime}$ that lies in the affine open set $X^{i}=X \cap \mathbb{U}^{i}=\operatorname{Spec} A_{i}$, as above.

A rational function $\alpha$ on $X$ or on $X^{\prime}$ is regular at $p$ if it is a regular function at $p$ on one of the open sets $X^{i}$. This means that one can write $\alpha$ as a fraction $a / b$ of elements of $A_{i}$, with $b(p) \neq 0$. Then the value of $\alpha$ at $p$ is $\alpha(p)=a(p) / b(p)$.
projectconic
3.4.4. Lemma. The regularity of a rational function at $p$ doesn't depend on the choice of the open set $X^{i}$ that contains $p$.

Let $X=\operatorname{Spec} A$ be an affine variety. As has been noted, we may regard $X$ as a quasiprojective variety by embedding it as a closed subset of $\mathbb{U}^{0}$. The function field of $X$ will be the field of fractions of its coordinate algebra $A$, and a rational function $\alpha$ on $X$ will be regular at a point $p$ of $X$ if it can be written as a fraction $\alpha=a / s$, where $a$ and $s$ are in $A$ and $s$ isn't zero at $p$. Thus $\alpha$ is regular at $p$ if it is an element of the coordinate algebra of some localization $X_{s}$ that contains $p$. Proposition 2.7.2 shows that the regular functions on an affine variety $\operatorname{Spec} A$ are the elements of $A$.
3.4.5. Lemma. Let $X$ be a variety. A rational function that is regular on $X$ is determined by the function it defines on $X$.
proof. We show that if the function is identically zero, then $\alpha=0$. We may assume that $U$ is affine, say $X=\operatorname{Spec} A$, where $A$ is a finite-type domain. Then what is to be proved is that an element $\alpha$ of $A$ that defines the zero function is the zero element. Or, equivalently, that the only element of $A$ which is in every maximal ideal is zero. This is Corollary 2.5.16(ii).

## (3.4.6) points with values in a field

Let $K$ be a field that contains the complex numbers. A point of projective space $\mathbb{P}^{n}$ with values in $K$ is an equivalence class of nonzero vectors $\alpha=\left(\alpha_{0}, \ldots, \alpha_{n}\right)$ with $\alpha_{i}$ in $K$, the equivalence relation being analogous to the one for ordinary points: $\alpha \sim \alpha^{\prime}$ if $\alpha^{\prime}=\lambda \alpha$ for some $\lambda$ in $K$.

If $K$ is the function field of a variety $X$, the embedding of $X$ into projective space $\mathbb{P}^{n}$ defines a point $\left(\alpha_{0}, \ldots, \alpha_{n}\right)$ of $X$ with values in $K$. To get this point, one may choose a standard affine open set, say $\mathbb{U}^{0}$, of $\mathbb{P}^{m}$ such that $X^{0}=X \cap \mathbb{U}^{0}$ isn't empty. Then $X^{0}$ is affine, say $X^{0}=\operatorname{Spec} A$. The embedding of $X^{0}$ into the affine space $\mathbb{U}^{0}$ is defined by a homomorphism $\mathbb{C}\left[u_{1}, \ldots, u_{n}\right] \rightarrow A$, and the images $\alpha_{i}$ of the variables $u_{i}$ in $A$ give us a point $\left(\alpha_{0}, \alpha_{1}, \ldots, \alpha_{n}\right)$ with values in $K$, and with $\alpha_{0}=1$. This is the point.

## (3.4.7) morphisms to projective space

For the rest of this section, it will be helpful to have a separate notation for the point with values in the function field $K$ of a variety that is determined by a nonzero vector $\alpha=\left(\alpha_{0}, \ldots, \alpha_{n}\right)$, with $\alpha_{i} \in K$. We'll denote that point by $\underline{\alpha}$. So $\underline{\alpha}=\underline{\alpha}^{\prime}$ if $\alpha^{\prime}=\lambda \alpha$ for some nonzero $\lambda$ in $K$. We'll drop this notation later.

We define a morphism from a variety $X$ to projective space using a point of $\mathbb{P}^{n}$ with values in the function field $K$ of $X$. When doing this, we must keep in mind that the points of projective space are equivalence classes of vectors, not the vectors themselves. As we will see, this complication turns out to be useful.

We begin with a simple example.
3.4.8. Example. Let $C$ be the conic in the projective plane $\mathbb{P}^{2}$ defined by $f\left(x_{0}, x_{1}, x_{2}\right)=x_{0} x_{1}+x_{0} x_{2}+x_{1} x_{2}$. We project $C$ to the line $L_{0}:\left\{x_{0}=0\right\}$, defining $C \xrightarrow{\pi} \mathbb{P}^{1}$ by $\pi\left(x_{0}, x_{1}, x_{2}\right)=\left(x_{1}, x_{2}\right)$. The formula for $\pi$ is undefined at the point $q=(1,0,0)$, though the map extends to the whole conic $C$.

Let's write this projection using a point with values in the function field $K$ of $C$. The affine open set $\left\{x_{0} \neq 0\right\}$ of $\mathbb{P}^{2}$ is the polynomial algebra $\mathbb{C}\left[u_{1}, u_{2}\right]$, with $u_{1}=x_{1} / x_{0}$ and $u_{2}=x_{2} / x_{0}$. We also denote by $u_{i}$ the restriction of the function $u_{i}$ to $C^{0}=C \cap \mathbb{U}^{0}$. The restricted functions are related by the equation $u_{1}+u_{2}+u_{1} u_{2}=0$ that is obtained by dehomogenizing $f$. We solve for $u_{2}: u_{2}=-u_{1} /\left(1+u_{1}\right)$.

The projection is given by $\pi\left(x_{0}, x_{1}, x_{2}\right)=\pi\left(1, u_{1}, u_{2}\right)=\left(u_{1},-u_{1} /\left(1+u_{1}\right)\right)$. Multiplying by $\lambda=$ $\left(1+u_{1}\right) / u_{1}$, we see that $\pi\left(x_{0}, x_{1}, x_{2}\right)$ is the point $\left(1+u_{1},-1\right)$. This formula is defined at all points at which $x_{0} \neq 0$, including at $q$. Thus $\pi(q)=(1,-1)$. The image of $q$ is the point at which the tangent line $L_{q}$ to $C$ at $q$ intersects $L_{0}$.

To define $\pi$ at the remaining points, we look on another standard affine open set. Let $v_{i}=x_{i} / x_{1}$ and $w_{i}=$ $x_{i} / x_{2}$. Then $\left(x_{0}, x_{1}, x_{2}\right)=\left(v_{0}, 1, v_{2}\right)=\left(w_{0}, w_{1}, 1\right)$. The projection can also be written as $\pi(x)=\left(1, v_{2}\right)$, which is valid at points at which $x_{1} \neq 0$ or as $\pi(x)=\left(w_{1}, 1\right)$, which is valid at points at which $x_{2} \neq 0$.

Let $K$ be the function field of a variety $Y$, and let $\alpha=\left(\alpha_{0}, \ldots, \alpha_{n}\right)$ be a nonzero vector with entries in $K$. We try to define a morphism from $Y$ to projective space $\mathbb{P}^{n}$ using the point $\underline{\alpha}$ with values in $K$. To define the image $\underline{\alpha}(q)$ of a point $q$ of $Y$ (an ordinary point), we look for a vector $\alpha^{\prime}=\left(\alpha_{0}^{\prime}, \ldots, \alpha_{n}^{\prime}\right)$, with $\underline{\alpha}=\underline{\alpha}$, i.e., $\alpha^{\prime}=\lambda \alpha$, such that the rational functions $\alpha_{i}^{\prime}$ are all regular at $q$ and not all zero there. Such a vector may exist or not. If it exists, we define

$$
\begin{equation*}
\underline{\alpha}(q)=\left(\alpha_{0}^{\prime}(q), \ldots, \alpha_{n}^{\prime}(q)\right) \quad\left(=\underline{\alpha}^{\prime}(q)\right) \tag{3.4.9}
\end{equation*}
$$

If such a vector $\alpha^{\prime}$ exists for every point $q$ of $Y$, we call $\underline{\alpha}$ a good point.
3.4.10. Lemma. A point $\underline{\alpha}$ of $\mathbb{P}^{n}$ with values in the function field $K_{Y}$ of $Y$ is a good point if either one of the two following conditions holds for every point $q$ of $Y$ :

- There is an element $\lambda$ in $K_{Y}$ such that the rational functions $\alpha_{i}^{\prime}=\lambda \alpha_{i}, i=0, \ldots, n$, are regular and not all zero at $q$.
- There is an index $j, 0 \leq j \leq n$, such that the rational functions $\alpha_{i} / \alpha_{j}, j=0, \ldots, n$, are regular at $q$.
proof. The first condition simply restates the definition. We show that it is equivalent with the second one.
Suppose that $\alpha_{i} / \alpha_{j}$ are regular at $q$ for all $i$. Let $\lambda=\alpha_{j}^{-1}$, and let $\alpha_{i}^{\prime}=\lambda \alpha_{i}=\alpha_{i} / \alpha_{j}$. The rational functions $\alpha_{i}^{\prime}$ are regular at $q$, and they aren't all zero there because $\alpha_{j}^{\prime}=1$.

Conversely, suppose that $\alpha_{i}^{\prime}=\lambda \alpha_{i}$ are all regular at $q$ and that $\alpha_{j}^{\prime}$ isn't zero there. Then $\alpha_{j}^{\prime-1}$ is a regular function at $q$, so the rational functions $\alpha_{i}^{\prime} / \alpha_{j}^{\prime}$, which are equal to $\alpha_{i} / \alpha_{j}$, are regular at $q$ for all $i$.
3.4.11. Lemma. With notation as in (3.4.9, the point $\underline{\alpha}(q)$ is independent of the choice of the vector $\alpha^{\prime}$.
oesntdepend

You will be able to supply the proof of this lemma
3.4.12. Definition. Let $Y$ be a variety with function field $K$. A morphism from $Y$ to projective space $\mathbb{P}^{n}$ is a map that is defined, as above, by a good point $\underline{\alpha}$ with values in $K$. We denote that morphism by $\underline{\alpha}$ too. 3.4.9.
3.4.13. Example. The identity map $\mathbb{P}^{1} \rightarrow \mathbb{P}^{1}$.

Let $X=\mathbb{P}^{1}$, and let $\left(x_{0}, x_{1}\right)$ be coordinates in $X$. The function field of $X$ is the field $K=\mathbb{C}(t)$ of rational functions in the variable $t=x_{1} / x_{0}$. The identity map $X \rightarrow X$ is the map $\underline{\alpha}$ defined by the point $\alpha=(1, t)$ with values in $K$. For every point $p$ of $X$ except the point $(0,1), \underline{\alpha}(p)=\alpha(p)=(1, t(p))$. For the point $q=(0,1)$, we let $\alpha^{\prime}=t^{-1} \alpha=\left(t^{-1}, 1\right)$. Then $\underline{\alpha}(q)=\alpha^{\prime}(q)=\left(x_{0}(q) / x_{1}(q), 1\right)=(0,1)$.

## (3.4.14)

## morphisms to quasiprojective varieties

3.4.15. Definition. Let $Y$ be a variety, and let $X$ be a subvariety of a projective space $\mathbb{P}^{n}$. A morphism of varieties $Y \xrightarrow{\underline{\alpha}} X$ is the restriction of a morphism $Y \xrightarrow{\underline{\alpha}} \mathbb{P}^{n}$ whose image is contained in $X$.

When a projective variety $X$ is the locus of zeros of a family $f$ of homogeneous polynomials, a morphism $Y \xrightarrow{\alpha} \mathbb{P}^{n}$ defines a morphism $Y \rightarrow X$ if $f(\alpha)=0$.

A word of caution: A morphism $Y \xrightarrow{\underline{\alpha}} X$ won't define a map of function fields $K_{X} \rightarrow K_{Y}$ unless the image of $Y$ is dense in $X$.
3.4.16. Proposition. A morphism of varieties $Y \xrightarrow{\alpha} X$ is a continuous map in the Zariski topology, and a continuous map in the classical topology.
proof. This proposition is rather trivial, once one has unraveled the notation. Let $\mathbb{U}^{i}$ be the standard affine open subset of $\mathbb{P}^{m}$, and let $Y^{i}$ be an affine open subset of the inverse image of $\mathbb{U}^{i}$. If $X=\mathbb{P}^{m}$, the restriction $Y^{i} \rightarrow \mathbb{U}^{i}$ of $\alpha$ is continuous in either topology because it is a morphism of affine varieties, as was defined in Section 2.7. Since $Y$ can be covered by affine open sets such as $Y^{i}, \underline{\alpha}$ is continuous. Continuity for a morphism to a subvariety $X$ of $\mathbb{P}^{m}$ follows, because the topology on $X$ is the induced topology.
3.4.17. Proposition. Let $X, Y$, and $Z$ be varieties and let $Z \xrightarrow{\underline{\beta}} Y$ and $Y \xrightarrow{\underline{\alpha}} X$ be morphisms. The composed map $Z \xrightarrow{\alpha \beta} X$ is a morphism.
proof. The proof is easy. Say that $X$ is a subvariety of $\mathbb{P}^{m}$. The morphism $\alpha$ is the restriction of a morphism $Y \rightarrow \mathbb{P}^{m}$ whose image is in $X$, and that is defined by a good point $\underline{\alpha}, \alpha=\left(\alpha_{0}, \ldots, \alpha_{m}\right)$ of $\mathbb{P}^{m}$ with values in the function field $K_{Y}$ of $Y$. Similarly, if $Y$ is a subvariety of $\mathbb{P}^{n}$, the morphism $\beta$ is the restriction of a morphism $Z \rightarrow \mathbb{P}^{n}$ whose image is contained in $Y$, and that is defined by a good point $\underline{\beta}, \beta=\left(\beta_{0}, \ldots, \beta_{n}\right)$ of $\mathbb{P}^{n}$ with values in the function field $K_{Z}$ of $Z$.

Let $z$ be a point (an ordinary point) of $Z$. Since $\beta$ is a good point, we may adjust $\beta$ by a factor in $K_{Z}$ so that the rational functions $\beta_{i}$ are regular and not all zero at $z$. Then $\underline{\beta}(z)$ is the point $\left(\beta_{0}(z), \ldots, \beta_{n}(z)\right)$. Let's denote that point by $q=\left(q_{0}, \ldots, q_{n}\right)$. So $q_{i}=\beta_{i}(z)$. The elements $\overline{\alpha_{j}}$ are rational functions on $Y$. We may adjust $\alpha$ by a factor in $K_{Y}$, so that they are regular and not all zero at $q$. Then $[\underline{\alpha \beta}](z)=\underline{\alpha}(q)=\left(\alpha_{0}(q), \ldots, \alpha_{m}(q)\right)$, and $\alpha_{j}(q)=\alpha_{j}\left(\beta_{0}(z), \ldots, \beta_{n}(z)\right)=\alpha_{j}(\beta(z))$ are not all zero. When these adjustments have been made, the point of $\mathbb{P}^{m}$ with values in $K_{Z}$ that defines $\underline{\alpha \beta}$ is $\left(\alpha_{0}(\beta(z)), \ldots, \alpha_{m}(\beta(z))\right)$.

This next is a lemma of topology.
3.4.18. Lemma. Let $\left\{X^{i}\right\}$ be a covering of a topological space $X$ by open sets. $A$ subset $Y$ of $X$ is open (or closed) if and only if $Y \cap X^{i}$ is open (or closed) in $X^{i}$ for every $i$. In particular, if $\left\{\mathbb{U}^{i}\right\}$ is the standard affine cover of $\mathbb{P}^{n}$, a subset $Y$ of $\mathbb{P}^{n}$ is open (or closed) if and only if $Y \cap \mathbb{U}^{i}$ is open (or closed) in $\mathbb{U}^{i}$ for every $i$.

### 3.4.19. Lemma.

(i) The inclusion of an open or a closed subvariety $Y$ into a variety $X$ is a morphism.
(ii) Let $Y \xrightarrow{f} X$ be a map whose image lies in an open or a closed subvariety $Z$ of $X$. Then $f$ is a morphism if and only if its restriction $Y \rightarrow Z$ is a morphism.
(iii) Let $\left\{Y^{i}\right\}$ be open an open covering of a variety $Y$, and let $Y^{i} \xrightarrow{f^{i}} X$ be morphisms. If the restrictions of $f^{i}$ and $f^{j}$ to the intersections $Y^{i} \cap Y^{j}$ are equal for all $i, j$, there is a unique morphism $f$ whose restriction to $Y^{i}$ is $f^{i}$.

We omit the proof, noting only that (iii) is trivial because the points with values in $K$ are all the same.

## (3.4.20) isomorphisms

A bijective morphism $Y \xrightarrow{u} X$ of quasiprojective varieties whose inverse function is also a morphism is an isomorphism. Isomorphisms are important because they allow us to identify different incarnations of the "same" variety, i.e., to describe an isomorphism class of varieties. For example, the projective line $\mathbb{P}^{1}$, a conic in $\mathbb{P}^{2}$, and the twisted cubic in $\mathbb{P}^{3}$ are isomorphic.

### 3.4.21. Example.

Let $Y$ denote the projective line, with coordinates $y_{0}, y_{1}$. As before, the function field of $Y$ is the field $K=\mathbb{C}(t)$ of rational functions in $t=y_{1} / y_{0}$. The degree 3 Veronese map $Y \longrightarrow \mathbb{P}^{3}$ 3.1.16 defines an isomorphism of $Y$ to its image, a twisted cubic $X$. The Veronese map is defined by the point $\alpha=\left(1, t, t^{2}, t^{3}\right)$ of $\mathbb{P}^{3}$ with values in $K$. On the open set $\left\{y_{0} \neq 0\right\}$ of $Y$, the rational functions $1, t, t^{2}, t^{3}$ are regular and not all zero. Let $\lambda=t^{-3}$ and $\alpha^{\prime}=\lambda \alpha=\left(t^{-3}, t^{-2}, t^{-1}, 1\right)$. The functions $t^{-k}$ are regular on the open set $\left\{y_{1} \neq 0\right\}$. So $\alpha$ is a good point that defines a morphism $Y \xrightarrow{\underline{\alpha}} X$.

The twisted cubic $X$ is the locus of zeros of the equations 3.1.17:

$$
v_{0} v_{2}=v_{1}^{2}, v_{2} v_{1}=v_{0} v_{3}, v_{1} v_{3}=v_{2}^{2}
$$

To identify the function field $F$ of $X$, we put $v_{0}=1$, obtaining relations $v_{2}=v_{1}^{2}, v_{3}=v_{1}^{3}$. Then $F$ is the field $\mathbb{C}\left(v_{1}\right)$. The point of $Y=\mathbb{P}^{1}$ with values in $F$ that defines the inverse of the morphism $\underline{\alpha}$ is $\beta=\left(1, v_{1}\right)$.
3.4.22. Lemma. Let $Y \xrightarrow{f} X$ be a morphism of varieties, let $\left\{X^{i}\right\}$ be an open covering of $X$, and let $Y^{i}=f^{-1} X^{i}$. If the restrictions $Y^{i} \xrightarrow{f^{i}} X^{i}$ of $f$ are isomorphisms, then $f$ is an isomorphism.
proof. Let $g^{i}$ denote the inverse of the morphism $f^{i}$. Then $g^{i}=g^{j}$ on $X^{i} \cap X^{j}$ because $f^{i}=f^{j}$ on $Y^{i} \cap Y^{j}$. By 3.4.19) (iii), there is a unique morphism $X \xrightarrow{g} Y$ whose restriction to $Y^{i}$ is $g^{i}$. That morphism is the inverse of $f$.

## (3.4.23) the diagonal

Let $X$ be a variety. The diagonal $X_{\Delta}$ is the set of points $(p, p)$ in $X \times X$. It is an example of a subset of $X \times X$ that is closed in the Zariski topology, but not closed in the product topology.
3.4.24. Proposition. Let $X$ be a variety. The diagonal $X_{\Delta}$ is a closed subvariety of the product variety $X \times X$.
proof. Let $\mathbb{P}$ denote the projective space $\mathbb{P}^{n}$ that contains $X$, and let $x_{0}, \ldots, x_{n}$ and $y_{0}, \ldots, y_{n}$ be coordinates in the two factors of $\mathbb{P} \times \mathbb{P}$. The diagonal $\mathbb{P}_{\Delta}$ in $\mathbb{P} \times \mathbb{P}$ is the closed subvariety defined by the bilinear equations $x_{i} y_{j}=x_{j} y_{i}$, or in the Segre variables, by the equations $w_{i j}=w_{j i}$, which show that the ratios $x_{i} / x_{j}$ and $y_{i} / y_{j}$ are equal.

Next, suppose that $X$ is the closed subvariety of $\mathbb{P}$ defined by a system of homogeneous equations $f(x)=$ 0 . The diagonal $X_{\Delta}$ can be identified as the intersection of the product $X \times X$ with the diagonal $\mathbb{P}_{\Delta}$ in $\mathbb{P} \times \mathbb{P}$, so it is a closed subvariety of $X \times X$. As a closed subvariety of $\mathbb{P} \times \mathbb{P}$, the diagonal $X_{\Delta}$ is defined by the equations

$$
\begin{equation*}
x_{i} y_{j}=x_{j} y_{i} \quad \text { and } \quad f(x)=0 \tag{3.4.25}
\end{equation*}
$$

The equations $f(y)=0$ are redundant. Finally, $X_{\Delta}$ is irreducible because it is homeomorphic to $X$.
It is interesting to compare Proposition 3.4 .24 with the Hausdorff condition for a topological space. The proof of the next lemma is often assigned as an exercise in topology.
3.4.26. Lemma. A topological space $X$ is a Hausdorff space if and only if, when $X \times X$ is given the product topology, the diagonal $X_{\Delta}$ is a closed subset of $X \times X$.

Though a variety $X$ with its Zariski topology isn't a Hausdorff space unless it is a point, Lemma 3.4.26 doesn't contradict Proposition 3.4.24 because the Zariski topology on $X \times X$ is finer than the product topology.

## (3.4.27) the graph of a morphism

Let $Y \xrightarrow{f} X$ be a morphism of varieties. The graph $\Gamma$ of $f$ is the subset of $Y \times X$ of pairs $(q, p)$ such that $p=f(q)$.
3.4.28. Proposition. The graph $\Gamma_{f}$ of a morphism $Y \xrightarrow{f} X$ is a closed subvariety of $Y \times X$, that is isomorphic to $Y$.
proof. We form a diagram of morphisms

graphdiagram
where $v$ sends a point $(q, p)$ of $\Gamma_{f}$ with $f(q)=p$ to $(p, p)$. The graph $\Gamma_{f}$ is the inverse image in $Y \times X$ of the diagonal $X_{\Delta}$. Since the diagonal is closed in $X \times X, \Gamma_{f}$ is closed in $Y \times X$.

Let $\pi_{1}$ denote the projection from $X \times Y$ to $Y$. The composition of the morphisms $Y \xrightarrow{(i d, f)} Y \times X \xrightarrow{\pi_{1}} Y$ is the identity map on $Y$, and the image of the map $(i d, f)$ is the graph $\Gamma_{f}$. Therefore $Y$ maps bijectively to $\Gamma_{f}$. The two maps $Y \rightarrow \Gamma_{f}$ and $\Gamma_{f} \rightarrow Y$ are inverses, so $\Gamma_{f}$ is isomorphic to $Y$.

$$
\mathbb{P}^{n} \xrightarrow{\pi} \mathbb{P}^{n-1}
$$

that drops the last coordinate of a point: $\pi\left(x_{0}, \ldots, x_{n}\right)=\left(x_{0}, \ldots, x_{n-1}\right)$ is called a projection. It is defined at all points of $\mathbb{P}^{n}$ except at the point $q=(0, \ldots, 0,1)$, which is called the center of projection. So $\pi$ is a morphism from the complement $U=\mathbb{P}^{n}-\{q\}$ to $\mathbb{P}^{n-1}$.

Let the coordinates in $\mathbb{P}^{n}$ and $\mathbb{P}^{n-1}$ be $x=x_{0}, \ldots, x_{n}$ and $y=y_{0}, \ldots, y_{n-1}$, respectively. The fibre $\pi^{-1}(y)$ over a point $\left(y_{0}, \ldots, y_{n-1}\right)$ is the set of points $\left(x_{0}, \ldots, x_{n}\right)$ such that $\left(x_{0}, \ldots, x_{n-1}\right)=\left(\lambda y_{0}, \ldots, \lambda y_{n-1}\right)$, while $x_{n}$ is arbitrary. It is the line in $\mathbb{P}^{n}$ through the points $\left(y_{1}, \ldots, y_{n-1}, 0\right)$ and $q=(0, \ldots, 0,1)$, with the center of projection $q$ omitted.

In Segre coordinates, the graph $\Gamma$ of $\pi$ in $U \times \mathbb{P}_{y}^{n-1}$ is the locus of solutions of the equations $w_{i j}=w_{j i}$ for $0 \leq i, j \leq n-1$, which imply that the vectors $\left(x_{0}, \ldots, x_{n-1}\right)$ and $\left(y_{0}, \ldots, y_{n-1}\right)$ are proportional.
3.4.32. Proposition. In $\mathbb{P}_{x}^{n} \times \mathbb{P}_{y}^{n-1}$, the locus $\bar{\Gamma}$ of the equations $x_{i} y_{j}=x_{j} y_{i}$, or $w_{i j}=w_{j i}$, with $0 \leq i, j \leq$ $n-1$, is the closure of the graph $\Gamma$ of $\pi$.
proof. The equations are true at points $(x, y)$ of $\Gamma$ at which $x \neq q$, and also at all points $(q, y)$. So the locus $\bar{\Gamma}$, a closed set, is the union of the graph $\Gamma$ and the set $q \times \mathbb{P}^{n-1}$. We must show that a homogeneous polynomial $g(w)$ that vanishes on $\Gamma$ vanishes at all points of $q \times \mathbb{P}^{n-1}$. Given $y$, let $x=\left(t y_{0}, \ldots, t y_{n-1}, 1\right)$. For all $t \neq 0$, the point $(x, y)$ is in $\Gamma$ and therefore $g(x, y)=0$. Since $g$ is a continuous function, $g(x, y)$ approaches $g(q, y)$ as $t \rightarrow 0$. So $g(q, y)=0$.

The projection $\bar{\Gamma} \rightarrow \mathbb{P}_{x}^{n}$ that sends a point $(x, y)$ to $x$ is bijective except when $x=q$. The fibre over $q$, which is $q \times \mathbb{P}^{n-1}$, is a projective space of dimension $n-1$. Because the point $q$ of $\mathbb{P}^{n}$ is replaced by a projective space in $\bar{\Gamma}$, the map $\bar{\Gamma} \rightarrow \mathbb{P}_{x}^{n}$ is called a blowup of the point $q$.

## figure: projection with closure of graph??

3.4.33. Proposition. Let $Y \xrightarrow{\underline{\alpha}} X$ and $Z \xrightarrow{\underline{\beta}} W$ be morphisms of varieties. The product map $Y \times Z \xrightarrow{\underline{\alpha} \times \beta}$ $X \times W$ that sends $(y, z)$ to $(\underline{\alpha}(y), \underline{\beta}(z))$ is a morphism
proof. Let $P$ and $q$ be points of $X$ and $Y$, respectively. We may assume that $\alpha_{i}$ are regular and not all zero at $p$ and that $\beta_{j}$ are regular and not all zero at $q$. Then, in the Segre coordinates $w_{i j}, \quad[\alpha \times \beta](p, q)$ is the point $w_{i j}=\alpha_{i}(p) \beta_{j}(q)$. We must show that $\alpha_{i} \beta_{j}$ are all regular at $(p, q)$ and are not all zero there. This follows from the analogous properties of $\alpha_{i}$ and $\beta_{j}$.

When defining morphisms varieties, one must keep in mind that points of projective space are equivalence classes of vectors, not the vectors themselves. This complication turns out to be very useful.

Some morphisms are sufficiently obvious that they don't require discussion. They include the projection from a product variety $X \times Y$ to $X$, the inclusion of $X$ into the product $X \times Y$ as the set $X \times y$ for some point $y$ of $Y$, the morphism of products $X \times Y \rightarrow X^{\prime} \times Y$ when a morphism $X \rightarrow X^{\prime}$ is given, and of course, the analogous maps when $Y$ replaces $X$.

If $X$ and $Y$ are subvarieties of projective spaces $\mathbb{P}^{m}$ and $\mathbb{P}^{n}$, respectively, a morphism $Y \rightarrow X$ will be determined by a morphism from $Y$ to $\mathbb{P}^{m}$ whose image is contained in $X$. However, it is important to note that a morphism $Y \xrightarrow{f} X$ needn't be the restriction of a morphism from $\mathbb{P}^{n}$ to $\mathbb{P}^{m}$. There will often be no way to extend the morphism from $Y$ to $\mathbb{P}^{n}$. It may not be possible to define $f$ using polynomials in the coordinate variables of $\mathbb{P}^{n}$.

For example, the Veronese map from the projective line $\mathbb{P}^{1}$ to $\mathbb{P}^{2}$, defined by $\left(x_{0}, x_{1}\right) \rightsquigarrow\left(x_{0}^{2}, x_{0} x_{1}, x_{1}^{2}\right)$, is an obvious morphism. Its image is the conic $C: v_{00} v_{11}-v_{01}^{2}=0$ in the projective plane $\mathbb{P}^{2}$. The Veronese defines a bijective morphism $\mathbb{P}^{1} \xrightarrow{f} C$ whose inverse function sends a point $\left(v_{00}, v_{01}, v_{11}\right)$ of $C$ with $v_{00} \neq 0$ to the point $\left(x_{0}, x_{1}\right)=\left(v_{01}, v_{11}\right)$. There is no way to extend the inverse function $f^{-1}$ to $\mathbb{P}^{2}$, though it is a morphism. In fact, there is no nonconstant morphism from $\mathbb{P}^{2}$ to $\mathbb{P}^{1}$.

In order to have a definition that includes all cases, we will define morphisms using points with values in a field.

## (3.4.34) the function field of a product

To define the function field of a product $X \times Y$ of projective varieties, one can use the Segre embedding $\mathbb{P}_{x}^{m} \times \mathbb{P}_{y}^{n} \rightarrow \mathbb{P}^{N}$. We use notation as in 3.1.10, and let's denote the product $X \times Y$ by $\Pi$. So $x_{i}, y_{j}$, and $w_{i j}$ are coordinates in the three projective spaces. The Segre map is defined by $w_{i j}=x_{i} y_{j}$. Let $\mathbb{U}^{i}, \mathbb{V}^{j}$, and $\mathbb{W}^{i j}$ be the standard affine open sets $x_{i} \neq 0, y_{j} \neq 0$ and $w_{i j} \neq 0$, respectively. The function field will be the field of fractions of any of the nonempty intersection $\Pi \cap \mathbb{W}^{i j}=\Pi^{i j}$, and $\Pi^{i j} \approx X^{i} \times Y^{j}$, where $X^{i}=X \cap \mathbb{U}^{i}$ and $Y^{j}=Y \cap \mathbb{V}^{j}$.

Since $\Pi^{i j}=X^{i} \times Y^{j}$, all that remains to do is to describe the field of fractions of a product of affine varieties $\Pi=X \times Y$, when $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$. If $A=\mathbb{C}[x] / P$ and $B=\mathbb{C}[y] / Q$, then in the notation of Proposition 3.3.4, the coordinate algebra of $\Pi$ is the algebra $\mathbb{C}[x, y] /\left(P^{\prime}+Q^{\prime}\right)$. This is the tensor product algebra $A \otimes B$. We don't need to know much about the tensor product algebra here, but let's use the tensor notation.

The function field $K_{X}$ of $X$ is the field of fractions of the coordinate algebra $A$. Similarly, $K_{Y}$ is the field of fractions of $B$ and $K_{X \times Y}$ is the field of fractions of $A \otimes B$. The one important fact to note is that $K_{X \times Y}$ isn't generated by $K_{X}$ and $K_{Y}$. For example, if $A=\mathbb{C}[x]$ and $B=\mathbb{C}[y]$ (one $x$ and one $y$ ), then $K_{X \times Y}$ is the field of rational functions in two variables $\mathbb{C}(x, y)$. The algebra generated by the fraction fields $\mathbb{C}(x)$ and $\mathbb{C}(y)$ consists of the rational functions $p(x, y) / q(x, y)$ in which $q(x, y)$ is a product $f(x) g(y)$ of a polynomial in $x$ and a polynomial in $y$. Most rational functions, $1 /(x+y)$ for example, aren't of this type.

But, $K_{X \times Y}$ is the fraction field of $A \otimes B$.

## (3.4.35) interlude: rational functions on projective space

Let $R$ denote the polynomial ring $\mathbb{C}\left[x_{0}, \ldots, x_{n}\right]$. A homogeneous polynomial $f$ of positive degree $d$ doesn't define a function on $\mathbb{P}^{n}$, because $f(\lambda x)=\lambda^{d} f(x)$. It does make sense to say that $f$ vanishes at a point of $\mathbb{P}^{n}$.

On the other hand, a fraction $g / h$ of homogeneous polynomials of the same degree $d$ does define a function wherever $h$ isn't zero, because

$$
g(\lambda x) / h(\lambda x)=\lambda^{d} g(x) / \lambda^{d} h(x)=g(x) / h(x)
$$

A homogeneous fraction $f$ is a fraction of homogeneous polynomials. The degree of a homogeneous fraction $f=g / h$ is the difference of degrees: $\operatorname{deg} f=\operatorname{deg} g-\operatorname{deg} h$.

A homogeneous fraction $f$ is regular at a point $p$ of $\mathbb{P}^{n}$ if, when it is written as a fraction $g / h$ of relatively prime homogeneous polynomials, the denominator $h$ isn't zero at $p$, and $f$ is regular on a subset $U$ if it is regular at every point of $U$. This definition agrees with the one given above, in Definition 3.4.2
3.4.36. Lemma. (i) Let $h$ be a homogeneous polynomial of positive degree $d$, and let $V$ be the open subset of $\mathbb{P}^{n}$ of points at which $h$ isn't zero. The nonzero rational functions that are regular on $V$ are those of the form $g / h^{k}$, where $k \geq 0$ and $g$ is a homogeneous polynomial of degree $d k$.
ratfnspspace
ho-
(ii) The only rational functions that are regular at every point of $\mathbb{P}^{n}$ are the constant functions.

For example, the homogeneous polynomials that don't vanish at any point of the standard affine open set $\mathbb{U}^{0}$ are the scalar multiples of powers of $x_{0}$. So the rational functions that are regular on $\mathbb{U}^{0}$ are those of the form $g / x_{0}^{k}, g$ homogeneous of degree $k$. This agrees with the fact that the coordinate algebra of $\mathbb{U}^{0}$ is the polynomial ring $\mathbb{C}\left[u_{1}, \ldots, u_{n}\right]$, with $u_{i}=x_{i} / x_{0}$ :, because $g\left(x_{0}, \ldots, x_{m}\right) / x_{0}^{k}=g\left(u_{0}, \ldots, u_{n}\right)$.
proof of Lemma 3.4 .36 (i) Let $\alpha$ be a regular function on the open set $U$, say $g_{1} / h_{1}$, where $g_{1}$ and $h_{1}$ are relatively prime homogeneous polynomials. Then $h_{1}$ doesn't vanish on $U$, so its zero locus in $\mathbb{P}^{n}$ is contained in the zero locus of $h$. According to the Strong Nullstellensatz 3.2.7, $h_{1}$ divides a power of $h$, say $h^{k}=f h_{1}$. Then $g_{1} / h_{1}=f g_{1} / f h_{1}=f g_{1} / h^{k}$.
(ii) If a rational function $f$ is regular at every point of $\mathbb{P}^{n}$, then it is regular on $\mathbb{U}^{0}$. so it will have the form $g / x_{0}^{k}$, where $g$ has degree $k$ and isn't divisible by $x_{0}$. And since $f$ is also regular on $\mathbb{U}^{1}$, it will have the form $h / x_{1}^{\ell}$, where $x_{1}$ doesn't divide $h$. Then $g x_{1}^{\ell}=h x_{0}^{k}$. Since $x_{0}$ doesn't divide $g, k=0, g$ is a constant, and $f=g$.

It is also true that the only rational functions on a projective variety $X$ that are regular at every point are the constant functions. The proof of this will be given later (see Corollary 8.3.10. When studying projective varieties, the constant functions are useless, so one has to look at at regular functions on open subsets. One way that affine varieties appear in projective algebraic geometry is as open subsets on which there are enough regular functions.

### 3.5 Affine Varieties

We have used the term 'affine variety' in several contexts:
A closed subset of affine space $\mathbb{A}_{x}^{n}$ is an affine variety, the set of zeros of a prime ideal $P$ of $\mathbb{C}[x]$. Its coordinate algebra is $A=\mathbb{C}[x] / P$.

The spectrum $\operatorname{Spec} A$ of a finite type domain $A$ is an affine variety that becomes a closed subvariety of affine space when one chooses a presentation $A=\mathbb{C}[x] / P$.

An affine variety becomes a variety in projective space when the ambient affine space $\mathbb{A}^{n}$ is identified with the standard open subset $\mathbb{U}^{0}$.

We combine these definitions now: An affine variety $X$ is a variety that is isomorphic to a variety of the form $\operatorname{Spec} A$.

If $X=\operatorname{Spec} A$ is an affine variety with function field $K$, its coordinate algebra $A$ will be the subalgebra of $K$ consisting of the regular functions on $X$. So $A$ and $\operatorname{Spec} A$, are determined uniquely by $X$, and the isomorphism Spec $A \rightarrow X$ is determined uniquely too. When $X$ is affine, it seems permissible to identify $X$ with $\operatorname{Spec} A$.

## (3.5.1) regular functions on affine varieties

Let $X=\operatorname{Spec} A$ be an affine variety. Its function field $K$ is the field of fractions of the coordinate algebra $A$. A rational function $\alpha$ is regular at a point $p$ of $X$ if it can be written as a fraction $a / s$ where $a, s$ are in $A$ and $s(p) \neq 0$, and $\alpha$ is regular on $X$ if it is regular at every point of $X$. As Proposition 2.7 .2 shows, the regular functions on an affine variety $\operatorname{Spec} A$ are the elements of $A$.

## Xmaps

3.5.2. Lemma.
(i) Let $R$ be the algebra of regular functions on a variety $Y$, and let $A$ be a finite-type domain. A homomorphism $A \rightarrow R$ defines a morphism $Y \xrightarrow{f} \operatorname{Spec} A$.
(ii) When $X$ and $Y$ are affine varieties, say $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$, morphisms $Y \rightarrow X$, as defined in (3.4.15), correspond bijectively to algebra homomorphisms $A \rightarrow B$, as in Definition 2.7.5

Note. Since $Y$ isn't affine, all that we know about the algebra $R$ is that it is a the subring of the function field of $Y$ of elements that are regular everywhere.
proof of Lemma 3.5.2 (i) Let $\left\{Y^{i}\right\}$ be an affine open covering of $Y$, and let $R_{i}$ be the coordinate algebra of $Y^{i}$. The inclusions $A \subset R \subset R_{i}$ define morphisms $Y^{i}=\operatorname{Spec} R_{i} \xrightarrow{f^{i}} \operatorname{Spec} A$. It is true that $f^{i}=f^{j}$ on $Y^{i} \cap Y^{j}$, so Lemma 3.4.19(iii) applies.
(ii) We choose a presentation of $A$, to embed $X$ as a closed subvariety of affine space, and we identify that affine space with the standard affine open set $\mathbb{U}^{0}$ of $\mathbb{P}^{n}$. Let $K$ be the function field of $Y$ - the field of fractions of $B$. A morphism $Y \xrightarrow{u} X$ is determined by a good point $\alpha$ with values in $K$, for which $\alpha_{0} \neq 0$. We may suppose that this point has the form $\alpha=\left(1, \alpha_{1} \ldots, \alpha_{n}\right)$. Then the rational functions $\alpha_{i}$ will be regular at every point of $Y$. They are elements of $B$. The coordinate algebra $A$ of $X$ is generated by the residues of the coordinate variables $x_{1}, \ldots, x_{n}$, with $x_{0}=1$. Sending $x_{i} \rightarrow \alpha_{i}$ defines a homomorphism $A \xrightarrow{\varphi} B$. Conversely, if $\varphi$ is such a homomorphism, the good point that defines the morphism $Y \xrightarrow{u} X$ is $\left(1, \varphi\left(x_{1}\right), \ldots, \varphi\left(x_{n}\right)\right)$.

An affine open subset of a variety $X$ is an open subset that is an affine variety. If $V$ is a nonempty open subset of $X$ and $R$ is the algebra of rational functions that are regular on $V$, then $V$ is an affine open subset if and only if $R$ is a finite-type domain and $V$ is isomorphic to $\operatorname{Spec} R$.

### 3.5.4. Proposition. The complement of a hypersurface is an affine open subvariety of $\mathbb{P}^{n}$.

\#\#\#rewrite this proof \#\#\#
proof. Let $V$ be the complement of the hypersurface $\{f=0\}$, where $f$ is an irreducible homogeneous polynomial of degree $d$, let $R$ be the algebra of regular functions on $V$, and let $K$ be its fraction field, the field of rational functions on $V$.

The regular functions on $V$ are the homogeneous fractions of degree zero of the form $g / f^{k}$ 3.4.35, and the fractions $m / f$, where $m$ is a monomial of degree $d$, generate $R$. Since there are finitely many monomials of degree $d, R$ is a finite-type domain. Let $w$ be an arbitrary monomial of degree $d-1$, and let $s_{i}=x_{i} w / f$. The point $\left(x_{0}, \ldots, x_{n}\right)$ of $V$ can also be written as $\left(s_{0}, \ldots, s_{n}\right)$, and the fractions $s_{i}$ are among the generators for $R$. Let $W=\operatorname{Spec} R$. Then $\left(s_{0}, \ldots, s_{n}\right)$ is a point of $W$ with values in $K$ that defines a morphism $W \xrightarrow{z} V$. We show that $z$ is an isomorphism.
3.5.5. Lemma. Let $\mathbb{U}^{i}$ be the standard affine open subset of $\mathbb{P}^{n}$. With $s_{i}$ as above, the intersection $V^{i}=V \cap \mathbb{U}^{i}$ is isomorphic to the localization $W_{s_{i}}$ of $W$.

Assuming the lemma, the morphism $W \xrightarrow{z} V$ restricts to an isomorphism $V^{i} \rightarrow$ Spec $R_{s_{i}}$. Since the sets $V^{i}=V \cap \mathbb{U}^{i}$ cover $V z$ is an isomorphism 3.4.22.
proof of Lemma ??. We work with the index $i=0$, as usual. Let $s=x_{0}^{d} / f$ and $t=s^{-1}=f / x_{0}^{d}$. Also, let $P$ be the coordinate algebra of $\mathbb{U}^{0}$. Then $V^{0}=V \cap \mathbb{U}^{0}$ is the set of points of $\mathbb{U}^{0}$ at which $t$ isn't zero. Its coordinate algebra is the localization $P_{t}$, and $V^{0}$ is the affine variety Spec $P_{t}$.

According to Lemma 3.4.22, the lemma will follow when we show that $P_{t}$ is the localization $R_{s}$ of $R$. With coordinates $u_{j}=x_{j} / x_{0}$ for $\mathbb{U}^{0}$, a fraction $m / f$, where $m$ is a monomial $x_{j_{1}} \cdots x_{j_{d}}$, can also be written as $u_{j_{1}} \cdots u_{j_{d}} / t$. These fractions generate $R$, so $R \subset P_{t}$, and since $s^{-1}=t$ is in $P_{t}, R_{s} \subset P_{t}$. For the other inclusion, we write $u_{j}=\left(x_{j} x_{0}^{d-1} / f\right) s^{-1}$. Because $x_{j} x_{0}^{d-1} / f$ is in $R, u_{j}$ is in $R_{s}$. Therefore $P \subset R_{s}$ and $P_{t} \subset R_{s}$. So $P_{t}=R_{s}$, as claimed.
3.5.6. Lemma. The affine open subsets of a variety $X$ form a basis for the topology on $X$.
proof. See Proposition ??.
3.5.7. Theorem Let $U$ and $V$ be affine open subvarieties of a variety $X$, say $U \approx \operatorname{Spec} A$ and $V \approx \operatorname{Spec} B$. The intersection $U \cap V$ is an affine open subvariety whose coordinate algebra is generated by the two algebras $A$ and $B$.
proof. We denote the algebra generated by two subalgebras $A$ and $B$ of the function field $K$ of $X$ by $[A, B]$. The elements of $[A, B]$ are finite sums of products of elements of $A$ and $B$. If $A=\mathbb{C}\left[a_{1}, \ldots, a_{r}\right]$, and $B=$ $\mathbb{C}\left[b_{1}, \ldots, b_{s}\right]$, then $[A, B]$ is generated by the set $\left\{a_{i}\right\} \cup\left\{b_{j}\right\}$. Let $W=\operatorname{Spec} R$, where $R=[A, B]$. We are to show that $W$ is isomorphic to $U \cap V$. The varieties $U, V$, and $W$ have the same function field $K$ as $X$, and the inclusions of coordinate algebras $A \rightarrow R$ and $B \rightarrow R$ give us morphisms $W \rightarrow U$ and $W \rightarrow V$. We also have inclusions $U \subset X$ and $V \subset X$, and $X$ is a subvariety of a projective space $\mathbb{P}^{n}$. Let $\alpha$ be the $\mathbb{P}^{n}$ with values in $K$ that defines the projective embedding $X \xrightarrow{\varphi} \mathbb{P}^{n}$. This point also defines morphisms $U \xrightarrow{u} \mathbb{P}^{n}$ and $V \xrightarrow{v} \mathbb{P}^{n}$ and $W \xrightarrow{\psi} \mathbb{P}^{n}$. The morphisms $u$ and $v$ are the restrictions of $\varphi$ to the open subsets $U$ and $V$ of $X$, respectively.

The morphism $W \xrightarrow{\psi} \mathbb{P}^{n}$ can be obtained as the composition of the morphisms $W \rightarrow U \subset X \rightarrow \mathbb{P}^{n}$, and also as the analogous composition, in which $V$ replaces $U$. Therefore the image of $W$ in $X$ is contained in $U \cap V$. (I suggest this slightly confusing point as an exercise.) Thus $\psi$ defines a morphism $W \xrightarrow{\epsilon} U \cap V$. We show that $\epsilon$ is an isomorphism.

Let $p$ be a point of $U \cap V$. We choose an affine open subset $Z$ of $U \cap V$ that is a localization of $U$ and a localization of $V$, and that contains $p$ 2.6.2)(ii). Let $S$ be the coordinate ring of $Z$. So $S=A_{s}$ for some nonzero $s$ in $A$ and also $S=B_{t}$ for some nonzero $t$ in $B$. Then

$$
R_{s}=[A, B]_{s}=\left[A_{s}, B\right]=[S, B]=S
$$

So $\epsilon$ maps the localization $W_{s}=\operatorname{Spec} R_{s}$ of $W$ isomorphically to the open subset $Z$ of $U \cap V$. Since we can cover $U \cap V$ by open sets such as $Z$, Lemma 3.4.19(ii) shows that $\epsilon$ is an isomorphism.

### 3.6 Lines in Projective Three-Space

The Grassmanian $G(m, n)$ is a variety whose points correspond to subspaces of dimension $m$ of the vector space $\mathbb{C}^{n}$, and to linear subspaces of dimension $m-1$ of $\mathbb{P}^{n-1}$. One says that $G(m, n)$ parametrizes those subspaces. For example, the Grassmanian $G(1, n+1)$ is the projective space $\mathbb{P}^{n}$. Points of $\mathbb{P}^{n}$ parametrize one-dimensional subspaces of $\mathbb{C}^{n+1}$.

The Grassmanian $G(2,4)$ parametrizes two-dimensional subspaces of $\mathbb{C}^{4}$, or lines in $\mathbb{P}^{3}$. In this section we describe that Grassmanian, denoting it by $\mathbb{G}$. The point of $\mathbb{G}$ that corresponds to a line $\ell$ in $\mathbb{P}^{3}$ will be denoted by $[\ell]$.

One can get some insight into the structure of $\mathbb{G}$ using row reduction. Let $V=\mathbb{C}^{4}$, let $u_{1}$, $u_{2}$ be a basis of a two-dimensional subspace $U$ of $V$ and let $M$ be the $2 \times 4$ matrix whose rows are $u_{1}, u_{2}$. The rows of the matrix $M^{\prime}$ obtained from $M$ by row reduction span the same space $U$, and the row-reduced matrix $M^{\prime}$ is uniquely determined by $U$. Provided that the left hand $2 \times 2$ submatrix of $M$ is invertible, $M^{\prime}$ will have the form

$$
M^{\prime}=\left(\begin{array}{llll}
1 & 0 & * & *  \tag{3.6.1}\\
0 & 1 & * & *
\end{array}\right)
$$

So the Grassmanian $\mathbb{G}$ contains, as an open subset, a four-dimensional affine space whose coordinates are the variable entries of $M^{\prime}$.

In any $2 \times 4$ matrix $M$ with independent rows, some pair of columns will be independent. Those columns can be used in place of the first two in a row reduction. So $\mathbb{G}$ is covered by six four-dimensional affine spaces that we denote by $\mathbb{W}^{i j}, 1 \leq i<j \leq 4$, $\mathbb{W}^{i j}$ being the space of $2 \times 4$ matrices such that column $n_{i}=(1,0)^{t}$ and column $n_{j}=(0,1)^{t}$. Since $\mathbb{P}^{4}$ and the Grassmanian are both covered by affine spaces of dimension four, they may seem similar, but they aren't the same.

## (3.6.2) the exterior algebra

Let $V$ be a complex vector space. The exterior algebra $\bigwedge V$ (read 'wedge $V$ ') is a noncommutative ring that contains the complex numbers and is generated by the elements of $V$, with the relations

$$
\begin{equation*}
v w=-w v \quad \text { for all } v, w \text { in } V . \tag{3.6.3}
\end{equation*}
$$

3.6.4. Lemma. The condition 3.6 .3 is equivalent with: $v v=0$ for all $v$ in $V$.
proof. To get $v v=0$ from 3.6.3), one sets $w=v$. Suppose that $v v=0$ for all $v$ in $V$. Then $(v+w)(v+w)=$ $v v=w w=0$. Since $(v+w)(v+w)=v v+v w+w v+w w, \quad v w+w v=0$.

To familiarize yourself with computation in $\Lambda V$, verify that $v_{2} v_{3} v_{1} v_{4}=v_{1} v_{2} v_{3} v_{4}$ and that $v_{2} v_{3} v_{4} v_{1}=$ $-v_{1} v_{2} v_{3} v_{4}$.

Let $\bigwedge^{r} V$ denote the subspace of $\bigwedge V$ spanned by products of length $r$ of elements of $V$. The exterior algebra $\bigwedge V$ is the direct sum of those subspaces $\bigwedge^{r} V$. An algebra $A$ that is a direct sum of subspaces $A^{i}$, and such that multiplication maps $A^{i} \times A^{j}$ to $A^{i+j}$ is called a graded algebra. Since its multiplication law isn't commutative, the exterior algebra is a noncommutative graded algebra.
3.6.5. Proposition. If $\left(v_{1}, \ldots, v_{n}\right)$ is a basis for $V$, the products $v_{i_{1}} \cdots v_{i_{r}}$ of length $r$ with increasing indices $i_{1}<i_{2}<\cdots<i_{r}$ form a basis for $\bigwedge^{r} V$.

The proof is at the end of the section.
3.6.6. Corollary. Let $v_{1}, \ldots, v_{r}$ be elements of $V$. The product $v_{1} \cdots v_{r}$ is zero in $\bigwedge^{r} V$ if and only if the set $\left(v_{1}, \ldots, v_{r}\right)$ is dependent.

For the rest of the section, we let $V$ be a vector space of dimension four with basis $\left(v_{1}, \ldots, v_{4}\right)$. Proposition 3.6 .5 tells us that

$$
\begin{aligned}
& \text { (3.6.7) } \\
& \bigwedge^{0} V=\mathbb{C} \text { is a space of dimension } 1, \text { with basis }\{1\} \\
& \bigwedge^{1} V=V \text { is a space of dimension } 4, \text { with basis }\left\{v_{1}, v_{2}, v_{3}, v_{4}\right\} \\
& \bigwedge^{2} V \text { is a space of dimension } 6, \text { with basis }\left\{v_{i} v_{j} \mid i<j\right\}=\left\{v_{1} v_{2}, v_{1} v_{3}, v_{1} v_{4}, v_{2} v_{3}, v_{2} v_{4}, v_{3} v_{4}\right\} \\
& \bigwedge^{3} V \text { is a space of dimension } 4, \text { with basis }\left\{v_{i} v_{j} v_{k} \mid i<j<k\right\}=\left\{v_{1} v_{2} v_{3}, v_{1} v_{2} v_{4}, v_{1} v_{3} v_{4}, v_{2} v_{3} v_{4}\right\} \\
& \bigwedge^{4} V \text { is a space of dimension } 1, \text { with basis }\left\{v_{1} v_{2} v_{3} v_{4}\right\} \\
& \bigwedge^{q} V=0 \text { when } q>4 .
\end{aligned}
$$

The elements of $\bigwedge^{2} V$ are combinations

$$
\begin{equation*}
w=\sum_{i<j} a_{i j} v_{i} v_{j} \tag{3.6.8}
\end{equation*}
$$

wedgetwo

We regard $\bigwedge^{2} V$ as an affine space of dimension 6 , identifying the combination $w$ with the vector whose coordinates are the six coefficients $a_{i j}(i<j)$. We use the same symbol $w$ to denote the point of the projective space $\mathbb{P}^{5}$ with those coordinates: $w=\left(a_{12}, a_{13}, a_{14}, a_{23}, a_{24}, a_{34}\right)$.
3.6.9. Definition. An element $w$ of $\bigwedge^{2} V$ is decomposable if it is a product of two elements of $V$.
3.6.10. Proposition. The decomposable elements of $\bigwedge^{2} V$ are those such that $w w=0$, and the relation $w w=0$ is given by the following equation in the coefficients $a_{i j}$ of $w=\sum_{i<j} a_{i j} v_{i} v_{j}$ :

$$
\begin{equation*}
a_{12} a_{34}-a_{13} a_{24}+a_{14} a_{23}=0 \tag{3.6.11}
\end{equation*}
$$

proof. If $w$ is decomposable, say $w=u_{1} u_{2}$, then $w^{2}=u_{1} u_{2} u_{1} u_{2}=-u_{1}^{2} u_{2}^{2}$ is zero because $u_{1}^{2}=0$. For the converse, we compute $w^{2}$ when $w=\sum_{i<j} a_{i j} v_{i} v_{j}$. The answer is

$$
w w=2\left(a_{12} a_{34}-a_{13} a_{24}+a_{14} a_{23}\right) v_{1} v_{2} v_{3} v_{4}
$$

To show that $w$ is decomposable if $w^{2}=0$, it seems simplest to factor $w$ explictly. Since the assertion is trivial when $w=0$, we may suppose that some coefficient of $w$, say $a_{12}$, is nonzero. Then if $w^{2}=0, w$ is the product

$$
\begin{equation*}
w=\frac{1}{a_{12}}\left(a_{12} v_{2}+a_{13} v_{3}+a_{14} v_{4}\right)\left(-a_{12} v_{1}+a_{23} v_{3}+a_{24} v_{4}\right) \tag{3.6.12}
\end{equation*}
$$

3.6.13. Corollary. (i) Let $w$ be a nonzero decomposable element of $\bigwedge^{2} V$, say $w=u_{1} u_{2}$, with $u_{i}$ in $V$. Then $\left(u_{1}, u_{2}\right)$ is a basis for a two-dimensional subspace of $V$.
(ii) If $\left(u_{1}, u_{2}\right)$ and $\left(u_{1}^{\prime}, u_{2}^{\prime}\right)$ are bases for the same subspace $U$ of $V$, then $w=u_{1} u_{2}$ and $w^{\prime}=u_{1}^{\prime} u_{2}^{\prime}$ differ by a scalar factor. Their coefficients represent the same point of $\mathbb{P}^{5}$.
(iii) Let $u_{1}, u_{2}$ be a basis for a two-dimensional subspace $U$ of $V$, and let $w=u_{1} u_{2}$. The rule $\epsilon(U)=w$ defines a bijection $\epsilon$ from $\mathbb{G}$ to the quadric $Q$ in $\mathbb{P}^{5}$ whose equation is (3.6.11).

Thus the Grassmanian $\mathbb{G}$ can be represented as the quadric 3.6.11.
proof. (i) If an element $w$ of $\bigwedge^{2} V$ is decomposable, say $w=u_{1} u_{2}$, and if $w$ is nonzero, then $u_{1}$ and $u_{2}$ must be independent 3.6.6. They span a two-dimensional subspace.
(ii) When we write the second basis in terms of the first one, say $\left(u_{1}^{\prime}, u_{2}^{\prime}\right)=\left(a u_{1}+b u_{2}, c u_{2}+d u_{2}\right)$, the product $u_{1}^{\prime} u_{2}^{\prime}$ becomes $(a d-b c) u_{1} u_{2}$, and $a d-b c \neq 0$.
linesinasurface
scontainslzero
(iii) In view of (i) and (ii), all that remains to show is that, if $\left(u_{1}, u_{2}\right)$ and $\left(u_{1}^{\prime}, u_{2}^{\prime}\right)$ are bases for distinct twodimensional subspaces $U$ and $U^{\prime}$, then in $\bigwedge^{2} V, u_{1} u_{2} \neq u_{1}^{\prime} u_{2}^{\prime}$. When $U \neq U^{\prime}$, the intersection $W=U \cap U^{\prime}$ will have dimension at most 1 . Then at least three of the vectors $u_{1}, u_{2}, u_{1}^{\prime}, u_{2}^{\prime}$ will be independent. Therefore $u_{1} u_{2} \neq u_{1}^{\prime} u_{2}^{\prime}$.

For the rest of this section, we use the algebraic dimension of a variety, a concept that will be studied in the next chapter. We refer to the algebraic dimension simply as the dimension. The dimension of a variety $X$ can be defined as the length $d$ of the longest chain $C_{0}>C_{1}>\cdots>C_{d}$ of closed subvarieties of $X$.

As was mentioned in Chapter 1, the topological dimension of $X$, its dimension in the classical topology, is always twice the algebraic dimension. Because the Grassmanian $\mathbb{G}$ is covered by affine spaces of dimension 4 , its algebraic dimension is 4 and its topological dimension is 8 .
3.6.14. Proposition. Let $\mathbb{P}^{3}$ be the projective space associated to a four dimensional vector space $V$. In the product $\mathbb{P}^{3} \times \mathbb{G}$, the locus $\Gamma$ of pairs $p,[\ell]$ such that the point $p$ of $\mathbb{P}^{3}$ lies on the line $\ell$ is a closed subset of dimension 5 .
proof. Let $\ell$ be the line in $\mathbb{P}^{3}$ that corresponds to the subspace $U$ with basis ( $u_{1}, u_{2}$ ), and say that $p$ represented by the vector $x$ in $V$. Let $w=u_{1} u_{2}$. Then $p \in \ell$ means $x \in U$, which is true if and only if $\left(x, u_{1}, u_{2}\right)$ is a dependent set, and this happens if and only if $x w=0 \sqrt{3.6 .5}$. So $\Gamma$ is the closed subset of points $(x, w)$ of $\mathbb{P}^{3} \times \mathbb{P}^{5}$ defined by the bihomogeneous equations $w^{2}=0$ and $x w=0$.

When we project $\Gamma$ to $\mathbb{G}$, The fibre over a point $[\ell]$ of $\mathbb{G}$ is the set of points $p,[\ell]$ such that $p$ is a point of $\ell$. The fibre maps bijectively to the line $\ell$. Thus $\Gamma$ can be viewed as a family of lines, parametrized by the four-dimensional variety $\mathbb{G}$ of $X$. Its dimension is $\operatorname{dim} \ell+\operatorname{dim} \mathbb{G}=1+4=5$.

## (3.6.15) lines on a surface

One may ask whether or not a given surface in $\mathbb{P}^{3}$ contains a line. One surface that contains lines is the quadric $Q$ in $\mathbb{P}^{3}$ with equation $w_{01} w_{10}=w_{00} w_{11}$, the image of the Segre embedding $\mathbb{P}^{1} \times \mathbb{P}^{1} \rightarrow \mathbb{P}_{w}^{3} \sqrt{\text { 3.1.7 }}$. It contains two families of lines, corresponding to the two "rulings" $p \times \mathbb{P}^{1}$ and $\mathbb{P}^{1} \times q$ of $\mathbb{P}^{1} \times \mathbb{P}^{1}$. There are surfaces of arbitrary degree that contain lines, but, that a generic surface of degree four or more doesn't contain any line.

We use coordinates $x_{i}$ with $i=1,2,3,4$ for $\mathbb{P}^{3}$ here. There are $N=\binom{d+3}{3}$ monomials of degree $d$ in four variables, so homogeneous polynomials of degree $d$ are parametrized by an affine space of dimension $N$, and surfaces of degree $d$ in $\mathbb{P}^{3}$ by a projective space of dimension $N-1$. Let $\mathbb{S}$ denote that projective space, and let $[S]$ denote the point of $\mathbb{S}$ that corresponds to a surface $S$. The coordinates of $[S]$ are the coefficients of the monomials in the defining polynomial $f$ of $S$. Speaking infomally, we say that a point of $\mathbb{S}$ "is" a surface of degree $d$ in $\mathbb{P}^{3}$. (When $f$ is reducible, its zero locus isn't a variety. Let's not worry about this.)

Consider the line $\ell_{0}$ defined by $x_{3}=x_{4}=0$. Its points are those of the form ( $x_{1}, x_{2}, 0,0$ ), so a surface $S:\{f=0\}$ will contain $\ell_{0}$ if and only if $f\left(x_{1}, x_{2}, 0,0\right)=0$ for all $x_{1}, x_{2}$. Substituting $x_{3}=x_{4}=0$ into $f$ leaves us with a polynomial in two variables:

$$
\begin{equation*}
f\left(x_{1}, x_{2}, 0,0\right)=c_{0} x_{1}^{d}+c_{1} x_{1}^{d-1} x_{2}+\cdots+c_{d} x_{2}^{d} \tag{3.6.16}
\end{equation*}
$$

where $c_{i}$ are some of the coefficients of the polynomial $f$. If $f\left(x_{1}, x_{2}, 0,0\right)$ is identically zero, all of those coefficients must be zero. So the surfaces that contain $\ell_{0}$ correspond to the points of the linear subspace $\mathbb{L}_{0}$ of $\mathbb{S}$ defined by the equations $c_{0}=\cdots=c_{d}=0$. Its dimension is $(N-1)-(d+1)=N-d-2$. This is a satisfactory answer to the question of which surfaces contain $\ell_{0}$, and we can use it to make a guess about lines in a generic surface of degree $d$.

### 3.6.17. Lemma. In the product variety $\mathbb{G} \times \mathbb{S}$, the set $\Gamma$ of pairs $[\ell],[S]$ such that $\ell \subset S$ is closed.

proof. Let $\mathbb{W}^{i j}, 1 \leq i<j \leq 4$ denote the six affine spaces that cover the Grassmanian, as at the beginning of this section. It suffices to show that the intersection $\Gamma^{i j}=\Gamma \cap\left(\mathbb{W}^{i j} \times \mathbb{S}\right)$ is closed in $\mathbb{W}^{i j} \times \mathbb{S}$ 3.4.18). We inspect the case $i, j=1,2$.

A line $\ell$ such that $[\ell]$ is in $\mathbb{W}^{12}$ corresponds to a subspace of $\mathbb{C}^{2}$ with basis of the form $u_{1}=\left(1,0, a_{2}, a_{3}\right)$, $u_{2}=\left(0,1, b_{2}, b_{3}\right)$ and $\ell$ is the line whose points are combinations $r u_{1}+s u_{2}$ of $u_{1}, u_{2}$. Let $f\left(x_{1}, x_{2}, x_{3}, x_{4}\right)$ be the polynomial that defines a surface $S$. The line $\ell$ is contained in $S$ if and only if $f\left(r, s, r a_{2}+s b_{2}, r a_{3}+s b_{3}\right)$
is zero for all $r$ and $s$. This is a homogeneous polynomial of degree $d$ in $r$, s. Let's call it $\widetilde{f}(r, s)$. If we write $\widetilde{f}(r, s)=z_{0} r^{d}+z_{1} r^{d-1} s+\cdots+z_{d} s^{d}$, the coefficients $z_{\nu}$ will be polynomials in $a_{i}, b_{i}$ and in the coefficients of $f$. The locus $z_{0}=\cdots=z_{d}=0$ is the closed set $\Gamma^{12}$ of $\mathbb{W}^{12} \times \mathbb{S}$.

The set of surfaces that contain our special line $\ell_{0}$ corresponds to the linear space $\mathbb{L}_{0}$ of $\mathbb{S}$ of dimension $N-d-2$, and $\ell_{0}$ can be carried to any other line $\ell$ by a linear map $\mathbb{P}^{3} \rightarrow \mathbb{P}^{3}$. So the sufaces that contain another line $\ell$ also form a linear subspace of $\mathbb{S}$ of dimension $N-d-2$. Those subspaces are the fibres of $\Gamma$ over $\mathbb{G}$. The dimension of the Grassmanian $\mathbb{G}$ is 4 . Therefore the dimension of $\Gamma$ is $\operatorname{dim} \Gamma=\operatorname{dim} \mathbb{L}_{0}+\operatorname{dim} \mathbb{G}=$ $(N-d-2)+4$. Since $\mathbb{S}$ has dimension $N-1$,

$$
\begin{equation*}
\operatorname{dim} \Gamma=\operatorname{dim} \mathbb{S}-d+3 \tag{3.6.18}
\end{equation*}
$$

We project the product $\mathbb{G} \times \mathbb{S}$ and its subvariety $\Gamma$ to $\mathbb{S}$. The fibre of $\Gamma$ over a point $[S]$ is the set of pairs [ $\ell],[S]$ such that $\ell$ is contained in $S$ - the set of lines in $S$.

When the degree $d$ of the surfaces we are studying is $1, \quad \operatorname{dim} \Gamma=\operatorname{dim} \mathbb{S}+2$. Every fibre of $\Gamma$ over $\mathbb{S}$ will have dimension at least 2 . In fact, every fibre has dimension equal to 2 . Surfaces of degree 1 are planes, and the lines in a plane form a two-dimensional family.

When $d=2, \operatorname{dim} \Gamma=\operatorname{dim} \mathbb{S}+1$. We can expect that most fibres of $\Gamma$ over $\mathbb{S}$ will have dimension 1 . This is true: A smooth quadric contains two one-dimensional families of lines. (All smooth quadrics are equivalent with the quadric (3.1.8 .) But if a quadratic polynomial $f\left(x_{1}, x_{2}, x_{3}, x_{4}\right)$ is the product of linear polynomials, its locus of zeros will be a union of planes. It will contain two-dimensional families of lines. Some fibres have dimension 2.

When $d \geq 4, \operatorname{dim} \Gamma<\operatorname{dim} \mathbb{S}$. The projection $\Gamma \rightarrow \mathbb{S}$ cannot be surjective. Most surfaces of degree 4 or more contain no lines.

The most interesting case is that $d=3$. In this case, $\operatorname{dim} \Gamma=\operatorname{dim} \mathbb{S}$. Most fibres will have dimension zero. They will be finite sets. In fact, a generic cubic surface contains 27 lines. We have to wait to see why the number is precisely 27 (see Theorem 4.8.16).

Our conclusions are intuitively plausible, but to be sure about them, we need to study dimension carefully. We do this in the next chapters.
proof of Proposition 3.6.5 Let $v=\left(v_{1}, \ldots, v_{n}\right)$ be a basis of the vector space $V$. The proposition asserts that the products $v_{i_{1}} \cdots v_{i_{r}}$ of length $r$ with increasing indices $i_{1}<i_{2}<\cdots<i_{r}$ form a basis for $\bigwedge^{r} V$.

To prove this, we need to be more precise about the definition of the exterior algebra $\wedge V$. We start with the algebra $T(V)$ of noncommutative polynomials in the basis $v$, which is also called the tensor algebra on $V$. The part $T^{r}(V)$ of $T(V)$ of degree $r$ has as basis the $n^{r}$ noncommutative monomials of degree $r$, products $v_{i_{1}} \cdots v_{i_{r}}$ of length $r$ of elements of the basis $v$. Its dimension is $n^{r}$. When $n=r=2, T^{2}(V)$ has the basis $\left(x_{1}^{2}, x_{1} x_{2}, x_{2} x_{1}, x_{2}^{2}\right)$.

The exterior algebra $\bigwedge V$ is the quotient of $T(V)$ obtained by forcing the relations $v w+w v=0$ 3.6.3. Using the distributive law, one sees that the relations $v_{i} v_{j}+v_{j} v_{i}=0,1 \leq i, j \leq n$, are sufficient to define this quotient. The relations $v_{i} v_{i}=0$ are included when $i=j$.

To obtain $\bigwedge^{r} V$, we multiply the relations $v_{i} v_{j}+v_{j} v_{i}$ on left and right by arbitrary noncommutative monomials $p(v)$ and $q(v)$ in $v_{1}, \ldots, v_{n}$ whose degrees add to $r-2$. The noncommutative polynomials

$$
\begin{equation*}
p\left(v_{i} v_{j}+v_{j} v_{i}\right) q \tag{3.6.19}
\end{equation*}
$$

span the kernel of the linear map $T^{r}(V) \rightarrow \bigwedge^{r} V$. So in $\bigwedge^{r} V, p\left(v_{i} v_{j}\right) q=-p\left(v_{j} v_{i}\right) q$. Using these relations, any product $v_{i_{1}} \cdots v_{i_{r}}$ in $\bigwedge^{r} V$ is, up to sign, equal to a product in which the elements $v_{i_{\nu}}$ are listed in increasing order. Thus the products with indices in increasing order span $\bigwedge^{r} V$, and because $v_{i} v_{i}=0$, such a product will be zero unless the indices are strictly increasing.

We go to the proof now. Let $v=\left(v_{1}, \ldots, v_{n}\right)$ be a basis for $V$. We show first that the product $w=v_{1} \cdots v_{n}$ in increasing order of the basis elements of $V$ is a basis of $\bigwedge^{n} V$. We have shown that this product spans $\bigwedge^{n} V$, and it remains to show that $w \neq 0$, or that $\bigwedge^{n} V \neq 0$.

Let's use multi-index notation: $(i)=\left(i_{1}, \ldots, i_{r}\right)$, and $v_{(i)}=v_{i_{1}} \cdots v_{i_{r}}$. We define a surjective linear map $T^{n}(V) \xrightarrow{\varphi} \mathbb{C}$ on the basis of $T^{n}(V)$ of products $v_{(i)}=\left(v_{i_{1}} \cdots v_{i_{n}}\right)$ of length $n$. If there is no repetition among the indices $i_{1}, \ldots, i_{n}$, then $(i)$ will be a permutation of the indices $1, \ldots, n$. In that case, we set $\varphi\left(v_{(i)}\right)=$ $\varphi\left(v_{i_{1}} \cdots v_{i_{n}}\right)=\operatorname{sign}(i)$. If there is a repetition, we set $\varphi\left(v_{(i)}\right)=0$.

Let $p$ and $q$ be noncommutative monomials whose degrees add to $n-2$. If the product $p\left(v_{i} v_{j}\right) q$ has no repeated index, the indices in $p\left(v_{i} v_{j}\right) q$ and $p\left(v_{j} v_{i}\right) q$ will be permutations of $1, \ldots, n$, and those permutations will have opposite signs. Then $p\left(v_{i} v_{j}+v_{j} v_{i}\right) q$ will be in the kernel of $\varphi$. Since these elements span the space of relations, $\varphi$ defines a surjective linear map $\bigwedge^{n} V \rightarrow \mathbb{C}$. Therefore $\bigwedge^{n} V \neq 0$.

To prove (3.6.5), we must show that for $r \leq n$, the products $v_{i_{1}} \cdots v_{i_{r}}$ with $i_{1}<i_{2}<\cdots<i_{r}$ form a basis for $\bigwedge^{r} V$, and we know that those products span $\bigwedge^{r} V$. We must show that they are independent. Suppose that a combination $z=\sum c_{(i)} v_{(i)}$ is zero, the sum being over sets of strictly increasing indices. We choose a set $\left(j_{1}, \ldots, j_{r}\right)$ of $n$ strictly increasing indices, and we let $(k)=\left(k_{1}, \ldots, k_{n-r}\right)$ be the set of $n-r$ indices that don't occur in $(j)$, listed in arbitrary order. Then all terms in the sum $z v_{(k)}=\sum c_{(i)} v_{(i)} v_{(k)}$ will be zero except the term with $(i)=(j)$. On the other hand, since $z=0, \quad z v_{(k)}=0$. Therefore $c_{(j)} v_{(j)} v_{(k)}=0$, and since $v_{(j)} v_{(k)}$ differs by sign from $v_{1} \cdots v_{n}$, it isn't zero. It follows that $c_{(j)}=0$. This is true for all $(j)$, so $z=0$.

# Chapter 4 INTEGRAL MORPHISMS OF AFFINE VARIETIES 

## july 18

4.1 The Nakayama Lemma
4.2 Integral Extensions
4.3 Finiteness of the Integral Closure
4.4 Geometry of Integral Morphisms
4.5 Dimension
4.6 Krull's Theorem
4.7 Chevalley's Finiteness Theorem

Double Planes

The concept of an algebraic integer was one of the important ideas that contributed to the development of algebraic number theory in the 19th century. Then, largely through the work of Noether and Zariski, an analog was seen to be essential in algebraic geometry. We study the analog in this chapter.

## Section 4.1 The Nakayama Lemma

nakaya-
malem
(4.1.1) eigenvectors

It won't surprise you that eigenvectors are important, but the way that they are used to study modules may be unfamiliar.

Let $P$ be an $n \times n$ matrix with entries in a ring $A$. The concept of an eigenvector for $P$ makes sense when the entries of a vector are in an $A$-module. A column vector $v=\left(v_{1}, \ldots, v_{n}\right)^{t}$ with entries in an $A$-module $M$ is an eigenvector of $P$ with eigenvalue $\lambda$ if $P v=\lambda v$.

When the entries of a vector are in a module, it becomes hard to adapt the usual requirement that an eigenvector must be nonzero, so we drop it, though the zero eigenvector tells us nothing.
4.1.2. Lemma. Let $p(t)$ be the characteristic polynomial $\operatorname{det}(t I-P)$ of a square matrix $P$. If $v$ is an eigenvector of $P$ with eigenvalue $\lambda$, then $p(\lambda) v=0$.

The usual proof, in which one multiplies the equation $(\lambda I-P) v=0$ by the cofactor matrix of $(\lambda I-P)$, carries over.

The next lemma is a cornerstone of the theory of modules.
4.1.3. Nakayama Lemma. Let $M$ be a finite module over a ring $A$, and let $J$ be an ideal of $A$ such that $M=J M$. There is an element $z$ in $J$ such that $m=z m$ for all $m$ in $M$, or such that $(1-z) M=0$.

By definition, $J M$ denotes the set of (finite) sums $\sum a_{i} m_{i}$ with $a_{i}$ in $J$ and $m_{i}$ in $M$.
Because it is always true that $M \supset J M$, the hypothesis $M=J M$ can be replaced by $M \subset J M$.
proof of the Nakayama Lemma. Let $v_{1}, \ldots, v_{n}$ be generators for the finite $A$-module $M$, and let $v$ be the vector $\left(v_{1}, \ldots, v_{n}\right)^{t}$. The equation $M=J M$ tells us that there are elements $p_{i j}$ in $J$ such that $v_{i}=\sum p_{i j} v_{j}$.

In matrix notation, $v=P v$. So $v$ is an eigenvector of $P$ with eigenvalue 1 , and if $p(t)$ is the characteristic polynomial of $P$, then $p(1) v=0$. Since the entries of $P$ are in $J$, inspection of the determinant of $I-P$ shows that $p(1)$ has the form $1-z$, with $z$ in $J$. Then $(1-z) v_{i}=0$ for all $i$. Since $v_{1}, \ldots$, $v_{n}$ generate $M$, $(1-z) M=0$.
4.1.4. Corollary. With notation as in the Nakayama Lemma, let $s=1-z$, so that $s M=0$. The localized module $M_{s}$ is the zero module.
4.1.5. Corollary. (i) Let $I$ and $J$ be ideals of a noetherian domain $A$. If $I=J I$, then either $I$ is the zero ideal or $J$ is the unit ideal.
(ii) Let $A \subset B$ be rings, and suppose that $B$ is a finite $A$-module. If $J$ is an ideal of $A$, and if the extended ideal $J B$ is the unit ideal of $B$, then $J$ is the unit ideal of $A$.
proof. (i) Since $A$ is noetherian, $I$ is a finite $A$-module. If $I=J I$, the Nakayama Lemma tells us that there is an element $z$ of $J$ such that $z x=x$ for all $x$ in $I$. Suppose that $I$ isn't the zero ideal. We choose a nonzero element $x$ of $I$. Because $A$ is a domain, we can cancel $x$ from the equation $z x=x$, obtaining $z=1$. Then 1 is in $J$, and $J$ is the unit ideal.
(ii) The elements of the extended ideal $J B$ are sums $\sum u_{i} b_{i}$ with $u_{i}$ in $J$ and $b_{i}$ in $B$ 2.6.3. Suppose that $B=J B$. The Nakayama Lemma tells us that there is an element $z$ in $J$ such that $z b=b$ for all $b$ in $B$. Setting $b=1$ shows that $z=1$. So $J$ is the unit ideal.
4.1.6. Corollary. Let $x$ be an element of a noetherian domain $A$, not a unit, and let $J$ be the principal ideal $x A$.
(i) The intersection $\bigcap J^{n}$ is the zero ideal.
(ii) If $y$ is a nonzero element of $A$, the integers $k$ such that $x^{k}$ divides $y$ in $A$ are bounded.
(iii) For every $k>0, J^{k}>J^{k+1}$.
proof. Let $I=\bigcap J^{n}$. The elements of $I$ are the ones that are divisible by $x^{n}$ for every $n$. Let $y$ be such an element. So for every $n$, there is an element $a_{n}$ in $A$ such that $y=a_{n} x^{n}$. Then $y / x=a_{n} x^{n-1}$, which is an element of $J^{n-1}$. Since this is true for every $n, y / x$ is in $I$, and $y$ is in $J I$. Since $y$ can be any nonzero element of $I, I=J I$. But since $x$ isn't a unit, $J$ isn't the unit ideal. So (i) tells us that $I=0$. This proves (i) and (ii). For (iii), we note that if $J^{k}=J^{k+1}$, then, multiplying by $J^{n-k}$ shows that $J^{n}=J^{n+1}$ for every $n \geq k$, and therefore that $\bigcap J^{n}=J^{k}$. But since $A$ is a domain and $x \neq 0, J^{k}=x^{k} A \neq 0$.
4.1.7. Corollary. Let I be a nonzero ideal of a noetherian domain $A$, and let $B$ be a domain that contains $A$ as subring. If $\beta$ is an element of $B$ and if $\beta I \subset I$, then $\beta$ is integral over $A$.
proof. Because $A$ is noetherian, $I$ is finitely generated. Let $v=\left(v_{1}, \ldots, v_{n}\right)^{t}$ be a vector whose entries generate $I$. The hypothesis $\beta I \subset I$ allows us to write $\beta v_{i}=\sum p_{i j} v_{j}$ with $p_{i j}$ in $A$, or in matrix notation, $P v=\beta v$. So $v$ is an eigenvector of $P$ with eigenvalue $\beta$, and if $p(t)$ is the characteristic polynomial of $P$, then $p(\beta) v=0$. Since at least one $v_{i}$ is nonzero and since $A$ is a domain, $p(\beta)=0$. The characteristic polynomial $p(t)$ is a monic polynomial with coefficients in $A$, so $\beta$ is integral over $A$.

## Section 4.2 Integral Extensions

Let $A$ be a domain. An extension of $A$ is a ring that contains $A$ as a subring. An element $\beta$ of an extension $B$ is integral over $A$ if it is a root of a monic polynomial

$$
\begin{equation*}
f(x)=x^{n}+a_{n-1} x^{n-1}+\cdots+a_{0} \tag{4.2.1}
\end{equation*}
$$

with coefficients $a_{i}$ in $A$, and an extension $B$ is an integral extension of $A$ if all of its elements are integral over $A$.
4.2.2. Lemma. Let $A \subset B$ be an extension of domains.
(i) An element b of $B$ is integral over $A$ if and only if the subring $A[b]$ of $B$ generated by $b$ is a finite $A$-module.
(ii) The set of elements of $B$ that are integral over $A$ is a subring of $B$.
(iii) If $B$ is generated as $A$-algebra by finitely many integral elements, it is a finite $A$-module.
(iv) Let $R \subset A \subset B$ be rings, and suppose that $A$ is an integral extension of $R$. An element of $B$ is integral over $A$ if and only if it is integral over $R$.

$$
\begin{equation*}
f(x, y)=a_{0}(x) y^{n}+a_{1}(x) y^{n-1}+\cdots+a_{n}(x) \tag{4.2.7}
\end{equation*}
$$

Let $x_{0}$ be a point of $X$. The fibre of $Y$ over $x_{0}$ consists of the points $\left(x_{0}, y_{0}\right)$ such that $y_{0}$ is a root of the one-variable polynomial $f\left(x_{0}, y\right)=\widetilde{f}(y)$. Because $f$ is irreducible, its discriminant $\delta(x)$ with respect to the variable $y$ isn't identically zero (1.7.19). For all but finitely many values of $x, a_{0}$ and $\delta$ will be nonzero. Then $\widetilde{f}(y)$ will have $n$ distinct roots.

When $f(x, y)$ is a monic polynomial in $y, u$ will be an integral morphism. If so, the leading term $y^{n}$ of $f$ will be the dominant term, when $y$ is large. Near to any point $x_{0}$ of $X$, there will be a positive real number $B$ such that

$$
\left|y^{n}\right|>\left|a_{1}(x) y^{n-1}+\cdots+a_{n}(x)\right|
$$

when $|y|>B$, and there fore $\widetilde{f}(y) \neq 0$. So the roots $y$ of $f(x, y)$ are bounded for all $x$ near to $x_{0}$.
On the other hand, when the leading coefficient $a_{0}(x)$ isn't a constant, $B$ won't be integral over $A$. If $x_{0}$ is a root of $a_{0}(x), f\left(x_{0}, y\right)$ will have degree less than $n$. What happens then is that, as a point $x_{1}$ of $X$ approaches $x_{0}$, at least one root of $f\left(x_{1}, y\right)$ tends to infinity. In calculus, one says that the locus $f(x, y)=0$ has a vertical asymptote at $x_{0}$.

To see this, we divide $f$ by its leading coefficient. Let $g(x, y)=f(x, y) / a_{0}=y^{n}+c_{1} y^{n-1}+\cdots+c_{n}$ with $c_{i}(x)=a_{i}(x) / a_{0}(x)$. For any $x$ at which $a_{0}(x)$ isn't zero, the roots of $g$ are the same as those of $f$. However, let $x_{0}$ be a root of $a_{0}$. Because $f$ is irreducible. At least one coefficient $a_{j}(x)$ doesn't have $x_{0}$ as a root. Then $c_{j}(x)$ is unbounded near $x_{0}$, and because the coefficient $c_{j}$ is an elementary symmetric function in the roots, the roots can't all be bounded.

This is the general picture: The roots of a polynomial remain bounded where the leading coefficient isn't zero. If the leading coefficient vanishes at a point, some roots are unbounded near that point.

## figure : nonmonic polynomial, but compare with figure for Hensel's Lemma

4.2.8. Noether Normalization Theorem. Let $A$ be a finite-type algebra over an infinite field $k$. There exist elements $y_{1}, \ldots, y_{n}$ in $A$ that are algebraically independent over $k$, such that $A$ is a finite module over its polynomial subalgebra $k\left[y_{1}, \ldots, y_{n}\right]$, i.e., $A$ is an integral extension of $k\left[y_{1}, \ldots, y_{n}\right]$.

When $K=\mathbb{C}$, the theorem can be stated by saying that every affine variety $X$ admits an integral morphism to an affine space.

The Noether Normalization Theorem is also true when $k$ is a finite field, though the proof given below needs to be modified.
4.2.9. Lemma. Let $k$ be an infinite field, and let $f(x)$ be a nonzero polynomial of degree $d$ in $x_{1}, \ldots, x_{n}$, with coefficients in $k$. After a suitable linear change of variable, the coefficient of $x_{n}^{d}$ in $f$ will be nonzero.
proof. Let $f_{d}$ be the homogeneous part of $f$ of maximal degree $d$. We regard $f_{d}$ as a polynomial function. Since $k$ is infinite, this function isn't identically zero. We choose coordinates $x_{1}, \ldots, x_{n}$ so that the point $q=(0, \ldots, 0,1)$ isn't a zero of $f_{d}$. Then $f_{d}\left(0, \ldots, 0, x_{n}\right)=c x_{n}^{d}$, and the coefficient $c$, which is $f_{d}(0, \ldots, 0,1)$, will be nonzero. By scaling $x_{n}$, we can make $c=1$.
proof of the Noether Normalization Theorem. Say that the finite-type algebra $A$ is generated by elements $x_{1}, \ldots, x_{n}$. If those elements are algebraically independent over $k, A$ will be isomorphic to the polynomial algebra $\mathbb{C}[x]$, and we will be done. If not, they will satisfy a polynomial relation $f(x)=0$ of some degree $d$, with coefficients in $k$. The lemma tells us that, after a suitable change of variable, the coefficient of $x_{n}^{d}$ in $f$ will be 1 . Then $f$ will be a monic polynomial in $x_{n}$ with coefficients in the subalgebra $R$ generated by $x_{1}, \ldots, x_{n-1}$. So $x_{n}$ will be integral over $R$, and $A$ will be a finite $R$-module. By induction on $n$, we may assume that $R$ is a finite module over a polynomial subalgebra $P$. Then $A$ will be a finite module over $P$ too.

The next proposition is an example of the general principle $\mathbf{5 . 1 . 1 5}$, that a construction involving finitely many operations can be done in a simple localization.
4.2.10. Proposition. Let $A \subset B$ be finite-type domains. There is a nonzero element $s$ in $A$ such that $B_{s}$ is a finite module over a polynomial subring $A_{s}\left[y_{1}, \ldots, y_{r}\right]$.
proof. Let $S$ be the set of nonzero elements of $A$, so that $K=A S^{-1}$ is the fraction field of $A$, and let $B_{K}=B S^{-1}$ be the ring obtained from $B$ by inverting all elements of $S$. Also, let $\beta=\left(\beta_{1}, \ldots, \beta_{k}\right)$ be a set of elements of the finite-type algebra $B$ that generates $B$ as algebra. Then $B_{K}$ is generated as $K$-algebra by $\beta$. It is a finite-type $K$-algebra. (A $K$-algebra is a ring that contains $K$ as subring.) The Noether Normalization Theorem tells us that $B_{K}$ is a finite module over a polynomial subring $P=K\left[y_{1}, \ldots, y_{r}\right]$. So $B_{K}$ is an integral extension of $P$. Any element of $B$ will be in $B_{K}$, and therefore it will be the root of a monic polynomial, say

$$
f(x)=x^{n}+c_{n-1} x^{n-1}+\cdots+c_{0}=0
$$

where the coefficients $c_{j}(y)$ are elements of $P$. Each coefficient $c_{j}$ is a combination of finitely many monomials in $y$, with coefficients in $K$. If $d \in A$ is a common denominator for those coefficients, $c_{j}(x)$ will have coefficients in $A_{d}[y]$. Since the generators $\beta$ of $B$ are integral over $P$, we may choose a denominator $s$ so that the generators $\beta_{1}, \ldots, \beta_{k}$ are all integral over $A_{s}[y]$. The algebra $B_{s}$ is generated over $A_{s}$ by $\beta$, so $B_{s}$ will be an integral extension of $A_{s}[y]$.

## Section 4.3 Normalization

Let $A$ be a domain with fraction field $K$. The normalization $A^{\#}$ of $A$ is the set of al elements of $K$ that are integral over $A$. It follows from Lemma 4.2 .2 (ii) that the normalization is a domain, and it contains $A$.

A domain $A$ is normal if it is equal to its normalization, and a normal variety $X$ is a variety that has an affine covering $\left\{X^{i}=\operatorname{Spec} A_{i}\right\}$ in which the algebras $A_{i}$ are normal domains.

To justify the definition of normal variety, we need to show that if an affine variety $X=\operatorname{Spec} A$ has an affine covering $X^{i}=\operatorname{Spec} A_{i}$, in which $A_{i}$ are normal domains, then $A$ is normal. This follows from Lemma 4.3 .3 (iii) below.

Our goal here is the next theorem, whose proof is at the end of the section.
4.3.1. Theorem. Let $A$ be a finite-type domain with fraction field $K$ of characteristic zero. The normalization $A^{\#}$ of $A$ is a finite $A$-module and a finite-type domain.

Thus there will be an integral morphism $\operatorname{Spec} A^{\#} \rightarrow \operatorname{Spec} A$.
The proof given here makes use of the characteristic zero hypothesis, though the theorem is true for a finitetype algebra over any field $k$.
4.3.2. Example. (normalization of a nodal cubic curve) The algebra $A=\mathbb{C}[u, v] /\left(v^{2}-u^{3}-u^{2}\right)$ can be embedded into the one-variable polynomial algebra $B=\mathbb{C}[x]$, by $u=x^{2}-1$ and $v=x^{3}-x$. The fraction

$$
\begin{equation*}
\beta^{n}+a_{1} \beta^{n-1}+\cdots+a_{n-1} \beta+a_{n}=0 \tag{4.3.4}
\end{equation*}
$$

with $a_{i}$ in $A$. We write $\beta=r / s$, where $r$ and $s$ are relatively prime elements of $A$. Multiplying by $s^{n}$ gives us the equation

$$
r^{n}=-s\left(a_{1} r^{n-1}+\cdots+a_{n} s^{n-1}\right)
$$

This equation shows that if a prime element of $A$ divides $s$, it also divides $r$. Since $r$ and $s$ are relatively prime, there is no such element. So $s$ is a unit, and $\beta$ is in $A$.
(ii) Let $\beta$ be an element of the fraction field of $A$ that is integral over $A_{s}$. There will be a polynomial relation of the form 4.3.4, except that the coefficients $a_{i}$ will be elements of $A_{s}$. The element $\gamma=s^{k} \beta$ satisfies the polynomial equation

$$
\gamma^{n}+\left(s^{k} a_{1}\right) \gamma^{n-1}+\cdots+\left(s^{(n-1) k} a_{n-1}\right) \gamma+\left(s^{n k} a_{n}\right)=0
$$

Since $a_{i}$ are in $A_{s}$, all coefficients of this polynomial will be in $A$ when $k$ is sufficiently large, and then $\gamma$ will be integral over $A$. Since $A$ is normal, $\gamma$ will be in $A$, and $\beta=s^{-k} \gamma$ will be in $A_{s}$.
(iii) This proof follows a common pattern. Suppose that $A_{s_{i}}$ is normal for every $i$. If an element $\beta$ of $K$ is integral over $A$, it will be in $A_{s_{i}}$ for all $i$, and $s_{i}^{n} \beta$ will be an element of $A$ if $n$ is large. We can use the same exponent $n$ for all $i$. Since $s_{1}, \ldots, s_{k}$ generate the unit ideal, so do their powers $s_{i}^{n}, \ldots, s_{k}^{n}$. Say that $\sum r_{i} s_{i}^{n}=1$, with $r_{i}$ in $A$. Then $\beta=\sum r_{i} s_{i}^{n} \beta$ is in $A$.

For the proof, it is convenient to state Theorem 4.3.1 more generally. The more general statement is essentially the same.

Let $A$ be a finite type domain with fraction field $K$, and let $L$ be a finite field extension of $K$. The integal closure of $A$ in $L$ is the set of all elements of $L$ that are integral over $A$. As Lemma 4.2.2 (ii) shows, the integral closure is a domain that contains $A$.
intclo 4.3.5. Theorem. Let $A$ be a finite type domain with fraction field $K$, and let $L$ be a finite field extension of $K$. The integal closure $B$ of $A$ in $L$ is a finite $A$-module.
abouttracetwo
4.3.3. Lemma. (i) A unique factorization domain is normal. In particular, a polynomial algebra over a field is normal.
(ii) If $s$ is a nonzero element of a normal domain $A$. The localization $A_{s}$ is normal.
(iii) Let $s_{1}, \ldots, s_{k}$ be nonzero elements of a domain $A$ that generate the unit ideal. If the localizations $A_{s_{i}}$ are normal for all $i$, then $A$ is normal.
proof. (i) Let $A$ be a unique factorization domain, and let $\beta$ be an element of its fraction field that is integral over $A$. Say that
4.3.6. Lemma. Let $A$ be a normal noetherian domain with fraction field $K$ of characteristic zero, and let $L$ be an algebraic field extension of $K$. An element $\beta$ of $L$ is integral over $A$ if and only if the coefficients of the
fields of $A$ and $B$ are equal because $x=v / u$, and the equation $x^{2}-(u+1)=0$ shows that $x$ is integral over $A$. The algebra $B$ is normal, so it is the normalization of $A$ (see Lemma 4.3.3(i)).

In this example, $\operatorname{Spec} B$ is the affine line $\mathbb{A}_{x}^{1}$, and the plane curve $C=\operatorname{Spec} A$ has a node at the origin $p=(0,0)$. The inclusion $A \subset B$ defines an integral morphism $\mathbb{A}_{x}^{1} \rightarrow C$ whose fibre over $p$ is the point pair $x= \pm 1$. The morphism is bijective at all other points. I think of $C$ as the variety obtained by gluing the points $x= \pm 1$ of the affine line together.

## figure: curve, not quite glued

In this example, the effect of normalization can be visualized geometrically. This isn't always so. Normalization is an algebraic process. Its effect on geometry may be subtle. monic irreducible polynomial $f$ for $\beta$ over $K$ are in $A$.
proof. If the monic polynomial $f$ has coefficients in $A$, then $\beta$ is integral over $A$. Suppose that $\beta$ is integral over $A$. Since we may replace $L$ by any field extension that contains $\beta$, we may assume that $L$ is a finite extension of $K$. A finite extension embeds into a Galois extension, so we may assume that $L / K$ is a Galois
extension. Let $G$ be its Galois group, and let $\left\{\beta_{1}, \ldots, \beta_{r}\right\}$ be the $G$-orbit of $\beta$, with $\beta=\beta_{1}$. The irreducible polynomial for $\beta$ over $K$ is

$$
\begin{equation*}
f(x)=\left(x-\beta_{1}\right) \cdots\left(x-\beta_{r}\right) \tag{4.3.7}
\end{equation*}
$$

Its coefficients are symmetric functions of the roots. If $\beta$ is integral over $A$, then all elements of the orbit are integral over $A$. Therefore the symmetric functions are integral over $A(4.2 .2)$ (iii), and since $A$ is normal, they are in $A$. So $f$ has coefficients in $A$.
4.3.8. Example. A polynomial in $A=\mathbb{C}[x, y]$ is square-free if it has no nonconstant square factors and isn't a constant. Let $f(x, y)$ be a square-free polynomial, and let $B$ denote the integral extension $\mathbb{C}[x, y, w] /\left(w^{2}-f\right)$ of $A$. Let $K$ and $L$ be the fraction fields of $A$ and $B$, respectively. Then $L=K[w] /\left(w^{2}-f\right)$ is a Galois extension of $K$. Its Galois group is generated by the automorphism $\sigma$ of order 2 defined by $\sigma(w)=-w$. The elements of $L$ have the form $\beta=a+b w$ with $a, b \in K$, and $\sigma(\beta)=\beta^{\prime}=a-b w$.

We show that $B$ is the integral closure of $A$ in $L$. Suppose that $\beta=a+b w$ is integral over $A$. If $b=0$, then $\beta=a$. This is an element of $A$ and therefore it is in $B$. If $b \neq 0$, the irreducible polynomial for $\beta=a+b w$ will be

$$
(x-\beta)\left(x-\beta^{\prime}\right)=x^{2}-2 a x+\left(a^{2}-b^{2} f\right)
$$

Because $\beta$ is integral over $A, 2 a$ and $a^{2}-b^{2} f$ are in $A$. Because the characteristic isn't 2 , this is true if and only if $a$ and $b^{2} f$ are in $A$. We write $b=u / v$, with $u, v$ relatively prime elements of $A$, so $b^{2} f=u^{2} f / v^{2}$. If $v$ weren't a constant, then since $f$ is square-free, it couldn't cancel $v^{2}$. So from $b^{2} f$ in $A$ we can conclude that $v$ is a constant and that $b$ is in $A$. Summing up, $\beta$ is integral if and only if $a$ and $b$ are in $A$, which means that $\beta$ is in $B$.

## (4.3.9) trace

Let $L$ be a finite field extension of a field $K$ and let $\beta$ be an element of $K$. When $L$ is viewed as a $K$-vector space, multiplication by $\beta$ becomes a linear operator $L \stackrel{\beta}{\longrightarrow} L$. The trace of this operator will be denoted by $\operatorname{tr}(\beta)$. The trace is a $K$-linear map $L \rightarrow K$.
4.3.10. Lemma. Let $L / K$ be a finite field extension, let $\beta$ be an element of $L$ of degree $r$ over $K$, let $f(x)=x^{r}+a_{1} x^{r-1}+\cdots+a_{r}$ be its irreducible polynomial over $K$, and let $K(\beta)$ be the extension of $K$ generated by $\beta$. Say that $[L: K(\beta)]=d$ and $[L: K]=n(=r d)$. Then $\operatorname{tr}(\beta)=-d a_{1}$. If $\beta$ is an element of $K$, then $\operatorname{tr}(\beta)=n \beta$.
proof. The set $\left(1, \beta, \ldots, \beta^{r-1}\right)$ is a $K$-basis for $K(\beta)$. On this basis, the matrix $M$ of multiplication by $\beta$ has the form illustrated below for the case $r=3$. Its trace is $-a_{1}$.

$$
M=\left(\begin{array}{lll}
0 & 0 & -a_{3} \\
1 & 0 & -a_{2} \\
0 & 1 & -a_{1}
\end{array}\right) .
$$

Next, let $\left(u_{1}, \ldots, u_{d}\right)$ be a basis for $L$ over $K(\beta)$. Then $\left\{\beta^{i} u_{j}\right\}$, with $i=0, \ldots, r-1$ and $j=1, \ldots, d$, will be a basis for $L$ over $K$. When this basis is listed in the order

$$
\left(u_{1}, u_{1} \beta, \ldots, u_{1} \beta^{n-1} ; u_{2}, u_{2} \beta, \ldots u_{2} \beta^{n-1} ; \ldots ; u_{d}, u_{d} \beta, \ldots, u_{d} \beta^{n-1}\right)
$$

the matrix of multiplication by $\beta$ will be made up of $d$ blocks of the matrix $M$.
4.3.11. Corollary. Let $A$ be a normal domain with fraction field $K$ and let $L$ be a finite field extension of $K$. If an element $\beta$ is integral over $A$, its trace is in $A$.

This follows from Lemmas 4.3.6 and 4.3.10
4.3.12. Lemma. Let $A$ be a normal noetherian domain with fraction field $K$ of characteristic zero, and let $L$ be a finite field extension of $K$. The form $L \times L \rightarrow K$ defined by $\langle\alpha, \beta\rangle=\operatorname{tr}(\alpha \beta)$ is $K$-bilinear, symmetric, and nondegenerate. If $\alpha$ and $\beta$ are integral over $A$, then $\langle\alpha, \beta\rangle$ is an element of $A$.

## clearde-

$$
\begin{equation*}
T: L \rightarrow K^{n} \tag{4.3.14}
\end{equation*}
$$

be the map defined by $T(\beta)=\left(\left\langle v_{1}, \beta\right\rangle, \ldots,\left\langle v_{n}, \beta\right\rangle\right)$, where $\langle$,$\rangle is the form defined in Lemma 4.3.12. This$ map is $K$-linear. If $\left\langle v_{i}, \beta\right\rangle=0$ for all $i$, then because $\left(v_{1}, \ldots, v_{n}\right)$ is a basis for $L,\langle\gamma, \beta\rangle=0$ for all $\gamma$ in $L$, and since the form is nondegenerate, $\beta=0$. Therefore $T$ is injective.

Let $B$ be the integral closure of $A$ in $L$. The basis elements $v_{i}$ are in $B$, and if $\beta$ is in $B, v_{i} \beta$ will be in $B$ too. Then $\left\langle v_{i}, \beta\right\rangle=\operatorname{tr}\left(v_{i} b\right)$ will be in $A$, and $T(\beta)$ will be in $A^{n} 4.3 .12$. When we restrict $T$ to $B$, we obtain an injective map $B \rightarrow A^{n}$ that we denote by $T_{0}$. Since $T$ is $K$-linear, $T_{0}$ is a $A$-linear. It is an injective homomorphism of $A$-modules that maps $B$ isomorphically to its image, a submodule of $A^{n}$. Since $A$ is noetherian, every submodule of the finite $A$-module $A^{n}$ is finitely generated. Therefore the image of $T_{0}$ is a finite $A$-module, and so is the isomorphic $A$-module $B$.

## Section 4.4 Geometry of Integral Morphisms

prmint
4.3.13. Lemma. Let $A$ be a domain with fraction field $K$, let $L$ be a field extension of $K$, and let $\beta$ be an element of $L$ that is algebraic over $K$. If $\beta$ is a root of a polynomial $f=a_{n} x^{n}+a_{n-1} x^{n-1}+\cdots+a_{0}$ with $a_{i}$ in $A$, then $\gamma=a_{n} \beta$ is integral over $A$.
proof. One finds a monic polynomial with root $\gamma$ by substituting $x=y / a_{n}$ into $f$ and multiplying by $a_{n}^{n-1}$.
proof of Theorem 4.3.1. Let $A$ be a finite-type domain with fraction field $K$ of characteristic zero, and let $L$ be a finite field extension of $K$. We are to show that the integral closure of $A$ in $L$ is a finite $A$-module.

Step 1. We may assume that $A$ is normal.
We use the Noether Normalization Theorem to write $A$ as a finite module over a polynomial subalgebra $R=\mathbb{C}\left[y_{1}, \ldots, y_{d}\right]$. Let $F$ be the fraction field of $R$. Then $K$ and $L$ are finite extensions of $F$. An element of $L$ will be integral over $A$ if and only if it is integral over $R(\sqrt[4.2 .2]{ }$ (iv)). So the integral closure of $A$ in $L$ is the same as the integral closure of $R$ in $L$. We replace $A$ by the normal algebra $R$ and $K$ by $F$.

## Step 2. Bounding the integral extension.

We assume that $A$ is normal. Let $\left(v_{1}, \ldots, v_{n}\right)$ be a $K$-basis for $L$ whose elements are integral over $A$. Such a basis exists because we can multiply any element of $L$ by a nonzero element of $K$ to make it integral (Lemma 4.3.13). Let

The main geometric properties of an integral morphism of affine varieties are summarized in the theorems in
proof. The form is obviously symmetric, and it is $K$-bilinear because multiplication is $K$-bilinear and trace is $K$-linear. A form is nondegenerate if its nullspace is zero, which means that when $\alpha$ is a nonzero element, there is an element $\beta$ such that $\langle\alpha, \beta\rangle \neq 0$. We let $\beta=\alpha^{-1}$. Then $\langle\alpha, \beta\rangle=\operatorname{tr}(1)$, which, according to 4.3.10, , is the degree $[L: K]$ of the field extension. It is here that the hypothesis on the characteristic of $K$ enters: The degree is a nonzero element of $K$.

If $\alpha$ and $\beta$ are integral over $A$, so is their product $\alpha \beta$ 4.2.2 (ii). Corollary 4.3.11 shows that $\langle\alpha, \beta\rangle$ is an element of $A$. this section, which show that the geometry is as nice as could be expected.

Let $Y \xrightarrow{u} X$ be an integral morphism. We say that a closed subvariety $D$ of $Y$ lies over a closed subvariety $C$ of $X$ if $C$ is the image of $D$.

We translate this definition to algebra: Let $A \subset B$ be an extension of finite-type domains. As before, if $I$ is an ideal of $A$, the extended ideal $I B$ is the ideal of $B$ generated by $I$. Its elements are finite sums $\sum u_{i} b_{i}$ with $u_{i}$ in $I$ and $b_{i}$ in $B$. The contraction of an ideal $J$ of $B$ is the ideal $J \cap A$ of $A$. The contraction of a prime ideal is a prime ideal.

Closed subvarieties of the affine variety $X=\operatorname{Spec} A$ correspond bijectively to prime ideals of $A$. In analogy with the terminology for closed subvarieties, we say that a prime ideal $Q$ of $B$ lies over a prime ideal $P$ of $A$ if its contraction is $P$. For example, if a point $y$ of $Y=\operatorname{Spec} B$ has image $x$ in $X$, the maximal ideal $\mathfrak{m}_{y}$ lies over the maximal ideal $\mathfrak{m}_{x}$.
4.4.1. Lemma. Let $A \rightarrow B$ be an integral extension, and let $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$. Also, let $P$ and $Q$ be prime ideals of $A$ and $B$, respectively, and let $C=V_{X}(P)=\operatorname{Spec} A / P$ and $D=V_{Y}(Q)=\operatorname{Spec} B / Q$. Then $Q$ lies over $P$ if and only if $D$ lies over $C$.
proof. If $Q$ lies over $P$, i.e., $P=Q \cap A$, then the canonical map $A / P \rightarrow B / Q$ is an integral extension. By Theorem 4.4.3. Spec $B / Q=D$ maps surjectively to $\operatorname{Spec} A / P=C$, which means that $D$ lies over $C$. Conversely, if $D=\operatorname{Spec} B / Q$ lies over $C=\operatorname{Spec} A / P$, the morphism $D \xrightarrow{\varphi} C$ gives us a map $A / P \rightarrow B / Q$, and since $D \rightarrow C$ is surjective, that map must be injective. This impies that $P=Q \cap A$.
4.4.2. Proposition. Let $Y \xrightarrow{u} X$ be an integral morphism of affine varieties, and let $D$ and $D^{\prime}$ be closed subvarieties of $Y$ that lie over closed subvarieties $C$ and $C^{\prime}$ of $X$, respectively. Then $D^{\prime}<D$ if and only if $C^{\prime}<C$.
proof. Since $C$ and $C^{\prime}$ are the images of $D$ and $D^{\prime}$, respectively, it is clear that if $C^{\prime}<C$, then $D^{\prime}<D$. For the other implication, we may replace $X$ and $Y$ by $C$ and $D$, respectively. Then what has to be shown is that if $C^{\prime}<X$, then $D^{\prime}<Y$. We go over to ideals. Say that $Q^{\prime}$ and $P^{\prime}$ are the prime ideals of $B$ and $A$ corresponding to $D^{\prime}$ and $C^{\prime}$, respectively. So $P^{\prime}$ is the contraction $Q^{\prime} \cap A$ of $Q^{\prime}$. What has to be shown is that if $Q^{\prime}$ is nonzero, then $P^{\prime}$ is nonzero.

Let $\beta$ be a nonzero element of $Q^{\prime}$. Then $\beta$ is integral over $A$, say $\beta^{n}+a_{n-1} \beta^{n-1}+\cdots+a_{0}=0$, with $a_{i} \in A$. If $a_{0}=0$, then because $B$ is a domain, we can cancel $\beta$ from the equation. So we may assume $a_{0} \neq 0$. The equation shows that $a_{0}$ is in $Q^{\prime}$, and since it is also in $A$, it is in $P^{\prime}$.
4.4.3. Theorem. Let $Y \xrightarrow{u} X$ be an integral morphism of affine varieties.
(i) The morphism $u$ is surjective, and its fibres have bounded cardinality.
(ii) The image of a closed subvariety of $Y$ is a closed subvariety of $X$, and the image of a closed subset of $Y$ is a closed subset of $X$.
(iii) The set of closed subvarieties of $Y$ that lie over a closed subvariety of $X$ is finite and nonempty.
proof. Let $Y \xrightarrow{u} X$ be the integral morphism, $Y=\operatorname{Spec} B$ and $X=\operatorname{Spec} A$, that corresponds to the inclusion $A \subset B$.
(i) (bounding the fibres) Let $\mathfrak{m}_{x}$ be the maximal ideal at point $x$ of $X$. Corollary 4.1 .5 (ii) shows that the extended ideal $\mathfrak{m}_{x} B$ isn't the unit ideal of $B$, so it is contained in a maximal ideal of $B$, say $\mathfrak{m}_{y}$, where $y$ is a point of $Y$, and then $x$ is the image of $y$. Therefore $u$ is surjective.

Let $y_{1}, \ldots, y_{r}$ be the points of $Y$ in the fibre over a point $x$ of $X$. Then for each $i$, the maximal ideal $\mathfrak{m}_{x}$ of $A$ is the contraction of the maximal ideal $\mathfrak{m}_{y_{i}}$ of $B$. To bound the number $r$, we use the Chinese Remainder Theorem to show that $B$ cannot be spanned as $A$-module by fewer than $r$ elements.

Let $k_{i}$ and $k$ denote the residue fields $B / \mathfrak{m}_{y_{i}}$, and $A / \mathfrak{m}_{x}$, respectively, all of these fields being abstractly isomorphic to $\mathbb{C}$. Let $\bar{B}=k_{1} \times \cdots \times k_{r}$. We form a diagram of algebra homomorphisms

which we interpret as a diagram of $A$-modules. The minimal number of generators of the $A$-module $\bar{B}$ is equal to its dimension as $k$-module, which is $r$. The Chinese Remainder Theorem asserts that $\varphi$ is surjective, so $B$ cannot be spanned by fewer than $r$ elements.
(ii) (the image of a closed set is closed) It is clear that the image of an irreducible set via a continuous map is irreducible, so it suffices to show that the image of a closed subvariety is closed. Let $D$ be the closed subvariety of $Y$ that corresponds to a prime ideal $Q$ of $B$, and let $P=Q \cap A$ be its contraction, which is a prime ideal of $A$. Let $C$ be the variety of zeros of $P$ in $X$. The coordinate algebras of the affine varieties $D$ and $C$ are $\bar{B}=B / Q$ and $\bar{A}=A / P$, respectively, and because $B$ is an integral extension of $A, \bar{B}$ is an integral extension of $\bar{A}$. By (i), the map $D \rightarrow C$ is surjective. Therefore $C$ is the image of $D$.
(iii) (subvarieties that lie over a closed subvariety) We will refer to Proposition 4.4 .2 here. The inverse image $Z=u^{-1} C$ of a closed subvariety $C$ is closed in $Y$. It is the union of finitely many irreducible closed subsets,
say $Z=D_{1}^{\prime} \cup \cdots \cup D_{k}^{\prime}$. Part (i) tells us that the image $C_{i}^{\prime}$ of $D_{i}^{\prime}$ is a closed subvariety of $X$. Since $u$ is surjective, $C=\bigcup C_{i}^{\prime}$, and since $C$ is irreducible, it is equal to at least one $C_{i}^{\prime}$. The components $D_{i}^{\prime}$ such that $C_{i}^{\prime}=C$ are the subvarieties that lie over $C$. Moreover, any subvariety $D$ that lies over $C$ will be contained in the inverse image $Z=\bigcup D_{i}^{\prime}$ of $C$. According to Proposition 4.4.2, there are no inclusions among subvarieties that lie over $C$. Therefore $D$ must be one of the $D_{i}^{\prime}$. So it is an element of a finite set.

## (4.4.4) finite group actions

Let $G$ be a finite group of automorphisms of a normal, finite-type domain $B$, let $A$ be the algebra of invariant elements of $B$. According to Theorem 2.8.5, $A$ is a finite-type domain, and $B$ is a finite $A$-module. Let $Y=\operatorname{Spec} B$ and $X=\operatorname{Spec} A$. Points of $X$ correspond to $G$-orbits of points of $Y$. Theorem4.4.6 below extends this correspondence to closed subvarieties.
4.4.5. Lemma. Let $G$ be a finite group of automorphisms of a normal, finite-type domain $B$ and let $A$ be the subalgebra of invariant elements. Let $L$ and $K$ be the fraction fields of $B$ and $A$, respectively.
(i) The algebra $A=B^{G}$ is normal.
(ii) Every element of $L$ can be written as a fraction $b / s$, with $b$ in $B$, and $s$ in $A$.
(iii) $L$ is a Galois extension of $K$, with Galois group $G$. The ring of $G$-invariants $L^{G}$ is $K$.

Since $B$ is a finite $A$-module, the results of Section 4.4 apply to the integral morphism $Y \xrightarrow{u} X$. Recall that, when $G$ operates on the left on $B$, it operates on the right on $Y$ 2.8.9.
4.4.6. Theorem. Let $G$ be a finite group of automorphisms of a normal, finite-type domain $B$, let $A$ be the algebra of invariant elements of $B$, and let $Y \xrightarrow{u} X$ be the integral morphism of varieties corresponding to the inclusion $A \subset B$.
(i) Let $\left\{D_{1}^{\prime}, \ldots, D_{r}^{\prime}\right\}$ and $\left\{D_{1}, \ldots, D_{s}\right\}$ be orbits of closed subvarieties $Y$ that lie over $C^{\prime}$ and $C$, respectively, in $X$. If $D_{i}^{\prime} \subset D_{j}$ for some $i$ and $j$, then $C^{\prime} \subset C^{\prime}$. If $C^{\prime} \subset C$, then every $D_{i}^{\prime}$ is contained in some $D_{j}$.
(ii) There is a bijective correspondence between closed subvarieties of $X$ and $G$-orbits of closed subvarieties of $Y$ :

$$
\{\text { closed subvarieties of } Y\} / G \quad \longleftrightarrow \quad \text { \{closed subvarieties of } X\}
$$

The orbit that corresponds to a closed subvariety $C$ of $X$ is the set of closed subvarieties of $Y$ that lie over $C$.
proof. (i) The first assertion is clear: If $D_{i}^{\prime} \subset D_{j}$, then $C^{\prime} \subset C$. Suppose that $D_{i}^{\prime} \not \subset D_{j}$ for all $i$ and $j$. We show that $C^{\prime} \not \subset C$. The lemma below, whose proof is left as an exercise, shows that there is an element $\beta$ that is identically zero on every $D_{i}$ and is not identically zero on any $D_{j}^{\prime}$. Then for all $\sigma$ in $G, \sigma \beta$ has the same property. So $\alpha=\prod \sigma \beta$ is an element of $A$ that is identically zero on every $D_{i}$ but not on any $D_{j}^{\prime}$. Therefore $\alpha$ is identically zero on $C$ but not on $C^{\prime}$. So $C^{\prime}$ isn't contained in $C$.
(ii) Let $D$ be a closed subvariety of $Y$ with image $C$ in $X$. Because points of $X$ correspond to $\sigma$-orbits in $Y$, the subvarieties $D \sigma$ in the orbit of $D$ have the same image $C$. The fact that distinct orbits of closed subvarieties of $Y$ lie over distinct closed subvarieties of $X$ will follow from (ii).
4.4.7. Lemma. Let $Y=\operatorname{Spec} B$ be an affine variety, let $D_{1}, \ldots, D_{n}$ be distinct closed subvarieties of $Y$ and let $V$ be a closed subset of $Y$. Assume that $V$ doesn't contain any of the subvarieties $D_{j}$. There is an element $\beta$ of $B$ that vanishes on $V$, but isn't identically zero on any $D_{j}$.

The next theorem concerns inclusions among closed subvarieties. It refers to the diagram below:

$$
\begin{array}{ll}
D^{\prime} \subset D & Y \\
C^{\prime} \subset C & X
\end{array}
$$

4.4.8. Theorem. Let $Y \xrightarrow{u} X$ be an integral morphism of affine varieties, and let $C^{\prime} \subset C$ be closed subvarieties of $X$.
(i) Every closed subvariety $D$ of $Y$ that lies over $C$ contains a closed subvariety $D^{\prime}$ that lies over $C^{\prime}$.
(ii) Suppose that $X$ normal. Every closed subvariety $D^{\prime}$ of $Y$ that lies over $C^{\prime}$ is contained in a closed subvariety $D$ that lies over $C$.

Property (ii) is more subtle than (i), as is indicated by the fact that $X$ is assumed normal. There is an example at the end of the section showing that the hypothesis of normality cannot be dropped.

In commutative algebra books, Theorem 4.4 .8 is stated in terms of prime ideals. Suppose given prime ideals $P \subset P^{\prime}$ of $A$ : Do there exist prime ideals $Q \subset Q^{\prime}$ of $B$ that lie over $P$ and $P^{\prime}$, respectively?

$$
\begin{array}{ll}
Q \subset Q^{\prime} & B \\
P \subset P^{\prime} & A
\end{array}
$$

The translation of Theorem 4.4.8 to prime ideals reads as follows:
4.4.9. Theorem. Let $A \subset B$ be an integral extension of finite-type domains, and let $P \subset P^{\prime}$ be prime ideals of $A$.
(i) Every prime ideal $Q$ that lies over $P$ is contained in a prime ideal $Q^{\prime}$ that lies over $P^{\prime}$.
(ii) Suppose that $A$ is normal. Then every prime ideal $Q^{\prime}$ that lies over $P^{\prime}$ contains a prime ideal $Q$ that lies over $P$.

The statements (i) and (ii) of this theorem are often called "going up", and "going down", respectively. Since inclusions are reversed when one passes to closed subvarieties, those terms aren't appropriate in Theorem 4.4.8.
proof of Theorem4.4.8 (i). We are given $C^{\prime} \subset C$ in $X$ and $D$ in $Y$ that lies over $C$, and we must find $D^{\prime}$. We replace $Y$ and $X$ by $D$ and $C$, respectively. Then what is to be proved is that there is a closed subvariety $D^{\prime}$ of $Y$ that lies over a given closed subvariety $C^{\prime}$. This is part (ii) of Theorem 4.4.3
(ii). Here, we are given an integral morphism $\operatorname{Spec} B=Y \xrightarrow{u} X=\operatorname{Spec} A$ with $X$ normal, and we are given closed subvarieties $C \supset C^{\prime}$ of $X$ and a closed subvariety $D^{\prime}$ of $Y$ that lies over $C^{\prime}$. We are to find a closed subvariety $D$ that lies over $C$, and that contains $D^{\prime}$. Let $K$ and $L$ denote the fraction fields of $A$ and $B$, respectively. Since $B$ is a finite $A$-module, $L$ is a finite extension of $K$.

Case 1: $L$ is a Galois extension of $K$ and $B$ is normal. Then $A$ and $B^{G}$ have the same fields of fractions, $B^{G}$ is an integral extension of $A$, and $A$ is normal. So $B^{G}=A$. This case follows from Theorem 4.4.6.

Case 2: the general case. We put $L$ into a Galois extension $F$, and we let $R$ be the integral closure of $B$ in $F$. Then $R$ is a finite $B$-module and a finite $A$-module. Let $Z=\operatorname{Spec} R$ and let $E$ be a closed subvariety of $Z$ that lies over $D$. Then $E$ also lies over $C$. By Case 1 , there is a closed subvariety $E^{\prime}$ of $Z$ that lies over $C^{\prime}$ and is contained in $E$. The image $D^{\prime}$ of $E^{\prime}$ in $Y$ is the required closed subvariety of $Y$.

$$
\begin{array}{rr}
E^{\prime} \subset E & Z \\
D & Y \\
C^{\prime} \subset C & X
\end{array}
$$

\#\#out of place ?? \#\#
4.4.10. Corollary. Let $Y \xrightarrow{u} X$ be an integral morphism of affine varieties with $X$ normal, and let $C$ be a closed subvariety of $X$. The subvarieties of $Y$ that lie over $C$ are the irreducible components of the inverse image of $C$.
proof. Let $Z$ be the inverse image of $C$, let $D^{\prime}$ be a component of $Z$, and let $C^{\prime}$ be its image in $X$. Then $C^{\prime} \subset C$, so by part (ii) of Theorem 4.4.8, $D^{\prime}$ is contained in a subvariety $D$ of $Y$ that lies over $C$. Because $D$ is contained in $Z$ and $D^{\prime}$ is a component of $Z, \quad D^{\prime}=D$.
4.4.11. Example. In this example, $B$ is the normalization of a finite-type domain $A, A$ isn't normal, and the conclusion of Theorem 4.4.8(ii) fails.

In the affine plane $Y=\operatorname{Spec} \mathbb{C}[u, v]$, let $L_{1}$ and $L_{2}$ be the lines $v=1$ and $v=-1$, respectively, let $D$ be the diagonal line $u=v$, and let $D^{\prime}$ be the point $(-1,1)$. We form an affine variety $X$ by gluing $L_{1}$ to $L_{2}$,
integralprimeprop
identifying the points $(u, 1)$ and $(u,-1)$. You can work out its coordinate algebra. Let $Y \xrightarrow{\pi} X$ be the gluing map. The real loci of $Y \rightarrow X$ are depicted below. In the figure, $C$ is the image of $D$, and $C^{\prime}$ is the image of $D^{\prime}$. Then $C^{\prime}$ is ontained in $C$, but because $D$ is the only subvariety of $Y$ that lies over $C$, there is no variety that lies over $C$ and that contains $D^{\prime}$.

## \#\#figure\#\#

### 4.5 Dimension

dimtheorem

Every variety has a dimension, and though dimension is a very coarse measure, but as is true for the dimension of a vector space, it is important.

A chain of closed subvarieties of a variety $X$ is a strictly decreasing sequence of closed subvarieties (of irreducible closed subsets)

$$
\begin{equation*}
C_{0}>C_{1}>C_{2}>\cdots>C_{k} \tag{4.5.1}
\end{equation*}
$$

The length of this chain is defined to be $k$. The chain is maximal if it cannot be lengthened by inserting another closed subvariety, which means that $C_{0}=X$, that there is no closed subvariety $\widetilde{C}$ with $C_{i}>\widetilde{C}>C_{i+1}$ for $i<k$, and that $C_{k}$ is a point. Theorem 4.5 .3 below shows that all maximal chains have the same length. The dimension of $X$, often denoted by $\operatorname{dim} X$, is the length of a maximal chain.

When $X$ is the affine variety $\operatorname{Spec} A$, the decreasing chain 4.5.1) corresponds to an increasing chain

$$
\begin{equation*}
P_{0}<P_{1}<P_{2}<\cdots<P_{k} \tag{4.5.2}
\end{equation*}
$$

of prime ideals of $A$ of length $k$, a prime chain. This prime chain is maximal if it cannot be lengthened by inserting another prime ideal, which means that $P_{0}$ is the zero ideal, that there is no prime ideal $\widetilde{P}$ with $P_{i}<\widetilde{P}<P_{i+1}$ for $i<k$, and that $P_{k}$ is a maximal ideal. The dimension $\operatorname{dim} A$ of a finite-type domain $A$ is the length $k$ of a maximal chain $(4.5 .2)$ of prime ideals. Thus if $X=\operatorname{Spec} A$, then $\operatorname{dim} X=\operatorname{dim} A$.
4.5.3. Theorem. (i) Let A be a finite-type domain whose fraction field $K$ has transcendence degree $n$. All prime chains in $A$ have length at most $n$, and all maximal prime chains have length equal to $n$. Therefore the dimension of $A$ is the transcendence degree of $A$.
(ii) Let $X$ be a variety whose function field $K$ has transcencence degree $n$. All chains of closed subvarieties of $X$ have length at most $n$, and all maximal chains have length $n$. Therefore the dimension of $X$ is equal to the transcendence degree of $K$.

The proof is below.
For example, the transcendence degree of the polynomial algebra $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ in $n$ variables is $n$, so the polynomial algebra has dimension $n$. The chain of prime ideals

$$
\begin{equation*}
0<\left(x_{1}\right)<\left(x_{1}, x_{2}\right)<\cdots<\left(x_{1}, \ldots, x_{n}\right) \tag{4.5.4}
\end{equation*}
$$

is a maximal prime chain. The corresponding chain

$$
\mathbb{P}^{n}>\mathbb{P}^{n-1}>\cdots>\mathbb{P}^{0}
$$

is a maximal chain of closed subvarieties of projective space $\mathbb{P}^{n}$.
The maximal chains of closed subvarieties of $\mathbb{P}^{2}$ have the form $\mathbb{P}^{2}>C>p$, where $C$ is a plane curve and $p$ is a point.

If (4.5.1) is a maximal chain in $X$, then $C_{0}=X$, and

$$
\begin{equation*}
C_{1}>C_{2}>\cdots>C_{k} \tag{4.5.5}
\end{equation*}
$$

will be a maximal chain in the variety $C_{1}$. So when $X$ has dimension $k$, the dimension of $C_{1}$ is $k-1$. Similarly, let 4.5.2 be a maximal chain of prime ideals in a finite-type domain $A$, let $\bar{A}=A / P_{1}$ and let $\bar{P}_{j}$ denote the image $P_{j} / P_{1}$ of $P_{j}$ in $\bar{A}$, for $j \geq 1$. Then

$$
\overline{0}=\bar{P}_{1}<\bar{P}_{2}<\cdots<\bar{P}_{k}
$$

will be a maximal prime chain in $\bar{A}$, and therefore the dimension of the domain $\bar{A}$ is $k-1$. There is a bijective correspondence between maximal prime chains in $\bar{A}$ and maximal prime chains in $A$ whose first term is $P_{0}$.

One more term: A closed subvariety $C$ of a variety $X$ has codimension codimension 1 if $C<X$ and if there is no closed set $\widetilde{C}$ with $C<\widetilde{C}<X$. A prime ideal $P$ of a noetherian domain has codimension 1 if it is not the zero ideal, and if there is no prime ideal $\widetilde{P}$ with $(0)<\widetilde{P}<P$. In the polynomial algebra $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$, the prime ideals of codimension 1 are the principal ideals generated by irreducible polynomials.
4.5.6. Proposition. Let $Y \xrightarrow{u} X$ be an integral morphism of varieties. Every chain of closed subvarieties of
proof. It follows from Proposition 4.4 .2 that the image of a prime chain in $Y$ is a prime chain in $X$. Theorem 4.4.8(i) and Proposition 4.4.2 show that \#\#\#
4.5.7. Lemma. Let $A=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ be a polynomial algebra, let $f$ be an irreducible element of $A$, and let $\bar{A}=A /(f)$. The transcendence degree of $\bar{A}$ is $n-1$.
proof. We may choose coordinates so that $f$ becomes a monic polynomial in $x_{n}$ with coefficients in $\mathbb{C}\left[x_{1}, \ldots, x_{n-1}\right]$, say $f=x_{n}^{k}+c_{k-1} x_{n}^{k-1}+\cdots+c_{0}$ (Lemma 4.2.9). Then $\bar{A}$ will be integral over $\mathbb{C}\left[x_{1}, \ldots, x_{n-1}\right]$, so it will have the same transcendence degree.
proof of theorem 4.5 .3 . (i) Induction allows us to assume the theorem true for a finite-type domain whose transcendence degree is less than $n$. Let $A$ be a finite-type domain of transcendence degree $n$.

Case 1: The case that $A$ is a polynomial algebra $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$.
Let $P_{0}<P_{1}<\cdots<P_{k}$ be a prime chain in $A$. We are to show that $k \leq n$, and that $k=n$ if the chain is maximal. We may assume that $P_{0}=0$ and that $P_{1}$ is a codimension 1 prime, generated by an irreducible polynomial $f$. If not, we insert a prime ideal into the chain. Let $\bar{A}=A / P_{1}$, and for $i \geq 1$, let $\bar{P}_{i}=P_{i} / P_{1}$. Then $\bar{P}_{1}<\bar{P}_{2}<\cdots<\bar{P}_{k}$ is a prime chain in $\bar{A}$ of length $k-1$, and if the chain $\left\{P_{i}\right\}$ is maximal, the chain $\left\{\bar{P}_{i}\right\}$ will be a maximal chain too. Lemma 4.5 .7 shows that $\bar{A}$ has transcendence degree $n-1$. So by induction, the length of the chain $\left\{\bar{P}_{i}\right\}$ is at most $n-1$ and is equal to $n-1$ if the chain $\left\{P_{i}\right\}$ is maximal. Therefore the chain $\left\{P_{i}\right\}$ has length at most $n$ and has length $n$ if it is maximal.

Case 2: The general case.
Let $B$ be a finite-type domain of transcendence degree $n$, and let $Q_{0}<Q_{1}<\cdots<Q_{k}$ be a prime chain in $B$. Again, we are to show that $k \leq n$ and that if the chain is maximal, then $k=n$. We apply the Noether Normalization Theorem: $B$ is a finite module over a polynomial subring $A$. The transcendence degree of $A$ is $n$, and the contractions $P_{i}=Q_{i} \cap A$ form a prime chain in $A$. This follows when one translates Proposition 4.5.6 to prime ideals. Therefore $k \leq n$.

Next, suppose that the chain $\left\{Q_{i}\right\}$ is maximal. Then $Q_{0}=0$ and therefore $P_{0}=0$. If $P_{1}$ were not a codimension 1 prime, we could choose a nonzero prime ideal $\widetilde{P}$ contained in $P_{1}$. Since $A$ is normal, we could apply Theorem 4.4.8 (ii): there would be a nonzero prime ideal $\widetilde{Q}$ of $B$ that lies over $\widetilde{P}$ and is contained in $Q_{1}$. This would imply that the chain $\left\{Q_{i}\right\}$ wasn't maximal, contrary to hypothesis. So $P_{1}$ is a codimension 1 prime, and by Lemma 4.5.7, $\bar{A}=A / P_{1}$ has transcendence degree $n-1$. Since $\bar{B}=B / Q_{1}$ is a finite $\bar{A}$-module, it also has transcendence degree $n-1$. Let $\bar{Q}_{i}=Q_{i} / Q_{1}$ for $i \geq 1$. By induction, the length of the maximal chain $\bar{Q}_{1}<\cdots<\bar{Q}_{k}$ in $\bar{B}$ is $n-1$, and therefore $k=n$.

Part (ii) of Theorem 4.5 .3 follows from the next lemma.
4.5.8. Lemma. Let $X^{\prime}$ be an open subvariety of a variety $X$. There is a bijective correspondence between chains $C_{0}>\cdots>C_{k}$ of closed subvarieties of a variety $X$ such that $C_{k} \cap X^{\prime} \neq \emptyset$ and chains $C_{0}^{\prime}>\cdots>C_{k}^{\prime}$ of closed subvarieties of $X^{\prime}$, defined by $C_{i}^{\prime}=C_{i} \cap X^{\prime}$. Given a chain $C_{i}^{\prime}$ in $X^{\prime}$, the corresponding chain in $X$ consists of the closures $C_{i}$ in $X$ of the varieties $C_{i}^{\prime}$.
proof. Suppose given a chain $C_{i}$ and that $C_{k} \cap X^{\prime} \neq \emptyset$. Then the intersections $C_{i}^{\prime}=C_{i} \cap X^{\prime}$ are dense open subsets of the irreducible closed sets $C_{i}$ 2.2.13. So the closure of $C_{i}^{\prime}$ is $C_{i}$, and since $C_{i}>C_{i+1}$, it is also true that $C_{i}^{\prime}>C_{i+1}^{\prime}$. Therefore $C_{0}^{\prime}>\cdots>C_{k}^{\prime}$ is a chain of closed subsets of $X^{\prime}$. Conversely, if $C_{0}^{\prime}>\cdots>C_{k}^{\prime}$ is a chain in $X^{\prime}$, the closures in $X$ will form a chain in $X$.
inte-graldimequal
krull
krullthm
chainavoids
4.5.9. Corollary. (i) If $X_{s}$ is a localization of an affine variety $X$, then $\operatorname{dim} X_{s}=\operatorname{dim} X$.
(ii) If $Y$ is a proper closed subvariety of a variety $X$, then $\operatorname{dim} Y<\operatorname{dim} X$.
(iii) If $Y \rightarrow X$ is an integral morphism of varieties, then $\operatorname{dim} Y=\operatorname{dim} X$.

### 4.6 Krull's Theorem

Krull's Principal Ideal Theorem completes our discussion of dimension. It asserts that the zero set of a principal ideal can't have a low dimension. Though the statement is natural, the proof isn't very easy.
4.6.1. Krull's Theorem. Let $X$ be an affine variety of dimension $d$, and let $\alpha$ be a nonzero element of its coordinate ring $A$. Every irreducible component of the zero locus $V_{X}(\alpha)$ of $\alpha$ in $X$ has dimension $d-1$.
proof. Let $C$ be an irreducible component of $V_{X}(\alpha)$. Since $\alpha$ isn't zero, $C$ is a proper subset of $X$, and its dimension is less than $d$. We must show that the dimension is $d-1$, and we prove this by contradiction. So we assume that $\operatorname{dim} C<d-1$.

Step 1: Reduction to the case that $A$ is normal.
Let $B$ be the normalization of $A$ and let $Y=\operatorname{Spec} B$. The dimension of $Y$ is $d$. The integral morphism $Y \rightarrow X$ is surjective, and it sends closed sets to closed sets 4.4.3). So the zero locus of $\alpha$ in $Y$ maps surjectively to the zero locus in $X$, and at least one irreducible component of $V_{Y}(\alpha)$, call it $D$, will map surjectively to $C$. The map $D \rightarrow C$ is also an integral morphism, so the dimension of $D$ is the same as that of $C$. We may therefore replace $X$ by $Y$ and $C$ by $D$. Hence we may assume that $A$ is normal.

Step 2: Reduction to the case that the zero locus of $\alpha$ is irreducible.
We do this by localizing. Say that the zero locus is $C \cup V$, where $C$ is a closed subvariety of codimension at least two, and $V$ is the union of the other irreducible components. We choose an element $s$ of $A$ that is identically zero on $V$ but not identically zero on $C$. Inverting $s$ eliminates the points of $V$, but $X_{s} \cap C=C_{s}$ will be nonempty. If $X$ is normal, so is $X_{s}$. Since localization doesn't change dimensions, we may replace $X$ and $C$ by $X_{s}$ and $C_{s}$.

Step 3: Completion of the proof.
This is the main step: We assume that $X$ is a normal affine variety, $X=\operatorname{Spec} A$, and that the zero locus of $\alpha$ in $X$ is a closed subvariety $C$ of codimension at least two. Then $C$ is the zero locus of a prime ideal $P$, and also the zero locus of $\alpha$. So $P$ is the radical of the principal ideal $\alpha A$, and $P^{n} \subset A \alpha$ if $n$ is large (see 2.5.11).

By what we know about dimension, $C$ will be contained in a closed subvariety $Z$ of codimension one. Let $Q$ be the prime ideal whose locus is $Z$. Then $P \supset Q$ because $C \subset Z$. On the other hand, $\alpha \notin Q$ because $\alpha$ vanishes only on $C$, and it follows that the principal ideal $A \alpha$ isn't contained in $Q$.
4.6.2. Lemma. With notation as above, There is an element $\gamma$ in $A$ such that $\gamma \notin A \alpha$ but $P \gamma \subset A \alpha$.
proof. Let $\beta$ be an element of $Q$. Corollary 4.1.6tells us that the powers of $\alpha$ that divide $\beta$ are bounded. Let $\alpha^{k}$ is the largest such power, and let $\beta^{\prime}=\beta / \alpha^{k}$. Then $\beta^{\prime}$ vanishes on the dense open subset $Z^{\prime}=Z-C$ of $Z$, so it vanishes on $Z$. So $\beta^{\prime}$ is in $Q$ but not in $A \alpha$.

Next, since $P^{n} \subset A \alpha$ for large $n$, it is also true that $P^{n} Q \subset A \alpha$ for large $n$. Let $r$ be the largest integer $\geq 0$ such that $P^{r} Q \not \subset A \alpha$. We choose an element $\gamma$ of $P^{r} Q$ that isn't in $A \alpha$ : $\gamma \notin A \alpha$, but $P \gamma \subset A \alpha$.

We finish the proof of Krull's Theorem now. Let $\delta=\gamma / \alpha$, where $\gamma$ is as in the lemma. Then $\delta \notin A$, but $P \delta \subset A$. Since $\gamma$ vanishes on $Z$ while $\alpha$ vanishes only on $C$, every element of $P \delta$ vanishes on the dense complement $Z^{\prime}$ of $C$ in $Z$, and therefore on $Z$. So $P \delta \subset Q \subset P$. Corollary 5.2.7 shows that $\delta$ is integral over $A$, and since $A$ is assumed normal, $\delta$ is in $A$. This is a contradiction that proves the theorem.
4.6.3. Corollary. Let $Z$ be a proper closed subset of a variety $X$ of dimension d, and let p be a point of $Z$. There is a maximal chain of closed subvarieties $X=X_{0}>X_{1}>\cdots>X_{d}$ with $X_{d}=\{p\}$, such that $X_{d-1}$ isn't contained in $Z$, and therefore $Z$ doesn't contain $X_{i}$ for any $i<d$.
proof. Lemma 4.5 .8 shows that we may assume that $X$ is affine: $X=\operatorname{Spec} A$. Lemma 4.4.7 asserts that $A$ contains an element $\alpha$ that vanishes at $p$ but doesn't vanish identically on any component of $Z$. Then at least one component of $V_{X}(\alpha)$ contains $p$. Let $X_{1}$ be such a component. Krull's Theorem tells us that $X_{1}$ has
dimension $d-1$. Since $\alpha$ doesn't vanish identically on any component of $Z, X_{1}$ doesn't contain any of those components. Therefore $Z_{1}=Z \cap X_{1}$ is a proper closed subvariety of $X_{1}$. We replace $X$ by $X_{1}$. Then the corollary follows by induction on the dimension $d$.

### 4.7 Chevalley's Finiteness Theorem

## (4.7.1) finite morphisms

The concepts of a finite morphism and an integral morphism of affine varieties were defined in Section 4.2. A morphism $Y \xrightarrow{u} X$ of affine varieties $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$ is a finite morphism if the homomorphism $A \xrightarrow{\varphi} B$ that corresponds to $u$ makes $B$ into a finite $A$-module. As was noted before, the difference between a finite morphism and an integral morphism of affine varieties is that for a finite morphism, the homomorphism $\varphi$ needn't be injective. If $\varphi$ is injective, $B$ will be an integral extension of $A$, and $u$ will be an integral morphism. We extend these definitions to varieties that aren't necessarily affine here.

By the restriction of a morphism $Y \xrightarrow{u} X$ to an open subset $X^{\prime}$ of $X$, we mean the induced morphism $Y^{\prime} \rightarrow X^{\prime}$, where $Y^{\prime}$ is the inverse image of $X^{\prime}$.
4.7.2. Definition. A morphism of varieties $Y \xrightarrow{u} X$ is a finite morphism if $X$ can be covered by affine open subsets $X^{i}$ such that the restriction of $u$ to each $X^{i}$ is a finite morphism of affine varieties, as defined in (4.2.4). A morphism $u$ is an integral morphism if $X$ can be overed by affine open sets to which the restriction of $u$ is an integral morphism of affine varieties.
4.7.3. Corollary. An integral morphism is a finite morphism. The composition of finite morphisms is a finite morphism. The inclusion of a closed subvariety into a variety is a finite morphism.

When $X$ is affine, Definition 4.2 .4 and Definition 4.7 .2 both apply. The next proposition shows that the two definitions are equivalent.
4.7.4. Proposition. Let $Y \xrightarrow{u} X$ be a finite or an integral morphism, as defined in $\sqrt{4.7 .2}$, and let $X^{\prime}$ be an affine open subset of $X$. The restriction of $u$ to $X^{\prime}$ is a finite or an integral morphism of affine varieties, as was defined in (4.2.4).
4.7.5. Lemma. (i) Let $A \xrightarrow{\varphi} B$ be a homomorphism of finite-type domains that makes $B$ into a finite $A$ module, and let s be a nonzero element of $A$. Then $B_{s}$ is a finite $A_{s}$-module.
(ii) The restriction of a finite (or an integral) morphism $Y \xrightarrow{u} X$ to an open subset of $X$ is a finite (or an integral) morphism, as in Definition 4.7.2
proof. (i) In the statement, $B_{s}$ denotes the localization of $B$ as $A$-module. This localization can also be obtained by localizing the algebra $B$ with respect to the image $s^{\prime}=\varphi(s)$, provided that it isn't zero. If $s^{\prime}$ is zero, then $s$ annihilates $B$, so $B_{s}=0$. In either case, a set of elements that spans $B$ as $A$-module will span $B_{s}$ as $A_{s}$-module, so $B_{s}$ is a finite $A_{s}$-module.
(ii) Say that $X$ is covered by affine open sets to which the restriction of $u$ is a finite morphism. The localizations of these open sets form a basis for the Zariski topology on $X$. So $X^{\prime}$ can be covered by such localizations. Part (i) shows that the restriction of $u$ to $X^{\prime}$ is a finite morphism.
proof of Proposition 4.7.4 We'll do the case of a finite morphism. The proof isn't difficult, but there are several things to check, and this makes the proof longer than one would like.

## Step 1. Preliminaries.

We are given a morphism $Y \xrightarrow{u} X, X$ is covered by affine open sets $X^{i}$, and the restrictions of $u$ to these open sets are finite morphisms of affine varieties. We are to show that the restriction to any affine open set $X^{\prime}$ is a finite morphism of affine varieties.

The affine open set $X^{\prime}$ is covered by the affine open sets $X^{\prime i}=X^{\prime} \cap X^{i}$, and the restrictions fo $X^{\prime i}$ are finite morphisms (4.7.5 (ii)). So we may replace $X$ by $X^{\prime}$. Since the localizations of an affine variety form a basis for its Zariski topology, we see that what is to be proved is this:

A morphism $Y \xrightarrow{u} X$ is given in which $X=\operatorname{Spec} A$ is affine. There are elements $s_{1}, \ldots, s_{k}$ that generate the unit ideal of $A$, such that for every $i$, the inverse image $Y^{i}$ of $X^{i}=X_{s_{i}}$ if nonempty, is affine, and its
finmorph
prodagain
deffinmorph
finiteexamp
onecoverfinite
coordinate algebra $B_{i}$ is a finite module over the localized algebra $A_{i}=A_{s_{i}}$. We must show that $Y$ is affine, and that its coordinate algebra $B$ is a finite $A$-module.

## Step 2. The algebra $B$ of regular functions on $Y$.

If $Y$ is affine, its coordinate algebra $B$ will be a finite-type domain, and $Y$ will be its spectrum. Since $Y$ isn't assumed to be affine, we don't know very much about $B$ other than that it is a subalgebra of the function field $L$ of $Y$. By hypothesis, the inverse image $Y^{i}$ of $X^{i}$, if nonempty, is affine, the spectrum of a finite $A_{i^{-}}$ algebra $B_{i}$. Since the localizations $X^{i}$ cover $X$, the affine varieties the $Y^{i}$ cover $Y$. We throw out the indices $i$ such that $Y^{i}$ is empty. Then a function is regular on $Y$ if and only if it is regular on each $Y^{i}$, and

$$
B=\bigcap B_{i}
$$

the intersection being in the function field $L$ of $Y$.
Let's denote the images in $B$ of the elements $s_{i}$ of $A$ by the same symbols.
Step 3. For any index $j, B_{j}$ is the localization $B\left[s_{j}^{-1}\right]$ of $B$.
The intersection $Y^{j} \cap Y^{i}$ is an affine variety. It is the set of points of the affine variety $Y^{j}=\operatorname{Spec} B_{j}$ at which $s_{i}$ isn't zero. Its coordinate algebra is the localization $B_{j}\left[s_{i}^{-1}\right]$. Then

$$
B\left[s_{j}^{-1}\right] \stackrel{(1)}{=} \bigcap_{i}\left(B_{i}\left[s_{j}^{-1}\right]\right) \stackrel{(2)}{=} \bigcap_{i} B_{j}\left[s_{i}^{-1}\right] \stackrel{(3)}{=} B_{j}\left[s_{j}^{-1}\right] \stackrel{(3)}{=} B_{j}
$$

where the explanation of the numbered equalities is as follows:
(1) A rational function $\beta$ is in $B_{i}\left[s_{j}^{-1}\right]$ if $s_{j}^{n} \beta$ is in $B_{i}$ for large $n$, and we can use the same exponent $n$ for all $i=1, \ldots, r$. Then $\beta$ is in $\bigcap_{i}\left(B_{i}\left[s_{j}^{-1}\right]\right)$ if and only if $s_{j}^{n} \beta$ is in $\bigcap_{i} B_{i}=B$, i.e., if and only if $\beta$ is in $B\left[s_{j}^{-1}\right]$.
(2) $B_{i}\left[s_{j}^{-1}\right]=B_{j}\left[s_{i}^{-1}\right]$ because $Y^{j} \cap Y^{i}=Y^{i} \cap Y^{j}$.
(3) For all $i, B_{j} \subset B_{j}\left[s_{i}^{-1}\right]$. Moreover, $s_{j}$ doesn't vanish on $Y^{j}$. It is a unit in $B_{j}$, and therefore $B_{j}\left[s_{j}^{-1}\right]=$ $B_{j}$. Then $B_{j} \subset \bigcap_{i} B_{j}\left[s_{i}^{-1}\right] \subset B_{j}\left[s_{j}^{-1}\right]=B_{j}$.

Step 4. $B$ is a finite $A$-module.
We choose a finite set $\left(b_{1}, \ldots, b_{n}\right)$ of elements of $B$ that generates the $A_{i}$-module $B_{i}$ for every $i$. We can do this because we can span the finite $A_{i}$-module $B_{i}=B\left[s_{i}^{-1}\right]$ by finitely many elements of $B$, and there are finitely many algebras $B_{i}$. We show that the set $\left(b_{1}, \ldots, b_{n}\right)$ generates the $A$-module $B$.

Let $x$ be an element of $B$. Since $x$ is in $B_{i}$, it is a combination of the elements $\left(b_{1}, \ldots, b_{n}\right)$ with coefficients in $A_{i}$. Then for large $k, s_{i}^{k} x$ will be a combination of those elements with coefficients in $A$, say

$$
s_{i}^{k} x=\sum_{\nu} a_{i, \nu} b_{\nu}
$$

with $a_{i, \nu} \in A$. We can use the same exponent $k$ for all $i$. Then with $\sum r_{i} s_{i}^{k}=1$,

$$
x=\sum_{i} r_{i} s_{i}^{k} x=\sum_{i} r_{i} \sum_{\nu} a_{i, \nu} b_{\nu}
$$

The right side is a combination of $b$ with coefficients in $A$.

## Step 5. $Y$ is affine.

The algebra $B$ of regular functions on $Y$ is a finite-type domain because it is a finite module over the finitetype domain $A$. Let $\widetilde{Y}=\operatorname{Spec} B$. The fact that $B$ is the algebra of regular functions on $Y$ gives us a morphism $Y \xrightarrow{\epsilon} \widetilde{Y}$ (Corollary 3.5.2. Restricting to the open subset $X^{j}$ of $X$ gives us a morphism $Y^{j} \xrightarrow{\epsilon^{j}} \widetilde{Y}^{j}$ in which $Y^{j}$ and $\widetilde{Y}^{j}$ are both equal to Spec $B_{j}$. Therefore $\epsilon^{j}$ is an isomorphism. Corollary 3.4.19 (ii) shows that $\epsilon$ is an isomorphism. So $Y$ is affine and by Step 4, its coordinate algebra $B$ is a finite $A$-module.

We come to Chevalley's theorem now. Let $\mathbb{P}$ denote the projective space $\mathbb{P}^{n}$ with coordinates $y_{0}, \ldots, y_{n}$.
4.7.6. Chevalley's Finiteness Theorem. Let $X$ be a variety, let $Y$ be a closed subvariety of the product $\mathbb{P} \times X$, and let $\pi$ denote the projection $Y \rightarrow X$. If all fibres of $\pi$ are finite sets, then $\pi$ is a finite morphism.
4.7.7. Corollary. Let $Y$ be a projective variety and let $Y \xrightarrow{\pi} X$ be a morphism whose fibres are finite sets. Then $\pi$ is a finite morphism. In particular, if $Y$ is a projective curve, any nonconstant morphism $Y \xrightarrow{u} X$ is a finite morphism.

This corollary follows from the theorem when one replaces $Y$ by the graph of $\pi$ in $Y \times X$. If $Y$ is embedded as a closed subvariety of $\mathbb{P}$, the graph will be a closed subvariety of $\mathbb{P} \times X$ (Proposition 3.4.28).

In the next lemma, $A$ denotes a finite-type domain, $B$ denotes a quotient of the algebra $A[u]$ of polynomials in $n$ variables $u_{1}, \ldots, u_{n}$ with coefficients in $A$, and $A \xrightarrow{\varphi} B$ denotes the canonical homomorphism. We'll use capital letters for nonhomogeneous polynomials here. If $G(u)$ is a polynomial in $A[u]$, we denote its image in $B$ by $G(u)$ too.
4.7.8. Lemma. Let $k$ be a positive integer. Suppose that, for each $i=1, \ldots, n$, there is a polynomial $G_{i}\left(u_{1}, \ldots, u_{n}\right)$ of degree at most $k-1$ in $n$ variables, with coefficients in $A$, such that $u_{i}^{k}=G_{i}(u)$ in $B$. Then $B$ is a finite $A$-module.
proof. Any monomial in $u_{1}, \ldots, u_{n}$ of degree $d \geq n k$ will be divisible by $u_{i}^{k}$ for at least one $i$. So if $m$ is a monomial of degree $d \geq n k$, the relation $u_{i}^{k}=G_{i}(u)$ shows that, in $B, m$ is equal to a polynomial in $u_{1}, \ldots, u_{n}$ of degree less than $d$, with coefficients in $A$. It follows by induction that the monomials in $u_{1}, \ldots, u_{n}$ of degree at most $n k-1$ span $B$.

Let $y_{0}, \ldots, y_{n}$ be coordinates in $\mathbb{P}^{n}$, and let $A\left[y_{0}, \ldots, y_{n}\right]$ be the algebra of polynomials in $y$ with coefficients in $A$. A homogeneous element of $A[y]$ is an element that is a homogeneous polynomial in $y$ with coefficients in $A$. A homogeneous ideal is an ideal that can be generated by homogeneous polynomials.
4.7.9. Lemma. Let $Y$ be a closed subset of $\mathbb{P} \times X$, where $X=\operatorname{Spec} A$ is affine,
(i) The ideal $\mathcal{I}$ of elements of $A[y]$ that vanish at every point of $Y$ is a homogeneous ideal of $A[y]$.
(ii) If the zero locus of a homogeneous ideal $\mathcal{I}$ of $A[y]$ is empty, then $\mathcal{I}$ contains a power of the irrelevant ideal $\mathcal{M}=\left(y_{0}, \ldots, y_{n}\right)$ of $A[y]$.
proof. (i) Let's write a point of $\mathbb{P} \times X$ as $q=\left(y_{0}, \ldots, y_{n}, x\right)$, with $x$ representing a point of $X$. So $(y, x)=$ $(\lambda y, x)$. Then the proof for the case $A=\mathbb{C}$ that is given in 1.3.2 carries over.
(ii) Let $V$ be the complement of the origin in the affine $n+1$-space with coordinates $y$. Then $V \times X$ maps to $\mathbb{P} \times X$ (see 3.2.4). If the locus of zeros of $\mathcal{I}$ in $\mathbb{P} \times X$ is empty, its locus of zeros in $V \times X$ will be contained in $o \times X, o$ being the origin in $\mathbb{P}$. Then the ideal of $o \times X$, which is generated by the elements $y_{0}, \ldots, y_{n}$, will contain $\mathcal{I}$.
proof of Chevelley's Finiteness Theorem. This proof is adapted from a proof by Schelter. We abbreviate the notation for a product of a variety $Z$ with $X$, denoting $Z \times X$ by $\widetilde{Z}$.

We are given a closed subvariety $Y$ of $\widetilde{\mathbb{P}}=\mathbb{P} \times X$, and the fibres over $X$ are finite sets. We are to prove that the projection $Y \rightarrow X$ is a finite morphism (4.2.4. We may assume that $X$ is affine, say $X=\operatorname{Spec} A$, and by induction on $n$, we may assume that the theorem is true when $\mathbb{P}$ is a projective space of dimension $n-1$.

Case 1. There is a hyperplane $H$ in $\mathbb{P}$ such that $Y$ is disjoint from $\widetilde{H}=H \times X$ in $\widetilde{\mathbb{P}}=\mathbb{P} \times X$.
This is the main case. We adjust coordinates $y_{0}, \ldots, y_{n}$ in $\mathbb{P}$ so that $H$ is the hyperplane at infinity $\left\{y_{0}=0\right\}$. Because $Y$ is a closed subvariety of $\widetilde{\mathbb{P}}$ disjoint from $\widetilde{H}, Y$ is also a closed subvariety of $\widetilde{\mathbb{U}}^{0}=\mathbb{U}^{0} \times X, \mathbb{U}^{0}$ being the standard affine $\left\{y_{0} \neq 0\right\}$. So $Y$ is affine.

Let $\mathcal{P}$ and $\mathcal{Q}$ be the (homogeneous) prime ideals in $A[y]$ that define $Y$ and $\widetilde{H}$, respectively, and let $\mathcal{I}=$ $\mathcal{P}+\mathcal{Q}$. So $\mathcal{Q}$ is the principal ideal of $A[y]$ generated by $y_{0}$. A homogeneous element of $\mathcal{I}$ of degree $k$ has the form $f(y)+y_{0} g(y)$, where $f$ is a homogeneous polynomial of degree $k$, and $g$ is a homogeneous polynomial of degree $k-1$.

The closed subsets $Y$ and $\widetilde{H}$ are disjoint. Since $Y \cap \widetilde{H}$ is empty, the $\operatorname{sum} \mathcal{I}=\mathcal{P}+\mathcal{Q}$ contains a power of the irrelevant ideal $\mathcal{M}=\left(y_{0}, \ldots, y_{n}\right)$. Say that $\mathcal{M}^{k} \subset \mathcal{I}$. Then $y_{i}^{k}$ is in $\mathcal{I}$ for $i=0, \ldots, n$. So we may write

$$
\begin{equation*}
y_{i}^{k}=f_{i}(y)+y_{0} g_{i}(y) \tag{4.7.10}
\end{equation*}
$$

projchevfin

## powerso-

fyenuf
hompoly-
coeffA
fplusg
with $f_{i}$ in $\mathcal{P}$ of degree $k$ and $g_{i}$ in $A[y]$ of degree $k-1$. We can omit the index $i=0$.

$$
\begin{equation*}
u_{i}^{k}=F_{i}(u)+G_{i}(u) \tag{4.7.11}
\end{equation*}
$$

The important point is that the degree of $G_{i}$ is at most $k-1$.
Recall that $Y$ is also a closed subset of $\mathbb{U}^{0}$. Its (nonhomogenous) ideal $P$ in $A[u]$ contains the polynomials $F_{1}, \ldots, F_{n}$, and its coordinate algebra is $B=A[u] / P$. In the quotient algebra $B$, the terms $F_{i}$ drop out, leaving us with equations $u_{i}^{k}=G_{i}(u)$, which are true in $B$. Since $G_{i}$ has degree at most $k-1$, Lemma 4.7.8 tells us that $B$ is a finite $A$-algebra, as was to be shown.

This completes the proof of Case 1.

## Case 2. the general case.

We have taken care of the case in which there exists a hyperplane $H$ such that $Y$ is disjoint from $\widetilde{H}$. The next lemma shows that we can cover the given variety $X$ by open subsets to which this special case applies. Then Lemma 4.7.4 and Proposition 4.7.4 will complete the proof.
4.7.12. Lemma. Let $Y$ be a closed subvariety of $\widetilde{\mathbb{P}}=\mathbb{P}^{n} \times X$, and suppose that the projection $Y \xrightarrow{\pi} X$ has finite fibres. Suppose also that Chevalley's Theorem has been proved for closed subvarieties of $\mathbb{P}^{n-1} \times X$. For every point $p$ of $X$, there is an open neighborhood $X^{\prime}$ of $p$ in $X$, and there is a hyperplane $H$ in $\mathbb{P}$, such that the inverse image $Y^{\prime}=\pi^{-1} X^{\prime}$ is disjoint from $\widetilde{H}$.
proof. Let $p$ be a point of $X$, and let $\widetilde{q}=\left(\widetilde{q}_{1}, \ldots, \widetilde{q}_{r}\right)$ be the finite set of points of $Y$ making up the fibre over $p$. We project $\widetilde{q}$ from $\mathbb{P} \times X$ to $\mathbb{P}$, obtaining a finite set $q=\left(q_{1}, \ldots, q_{r}\right)$ of points of $\mathbb{P}$, and we choose a hyperplane $H$ in $\mathbb{P}$ that avoids this finite set. Then $\widetilde{H}$ avoids the fibre $\widetilde{q}$. Let $W$ denote the closed set $Y \cap \widetilde{H}$. Because the fibres of $Y$ over $X$ are finite, so are the fibres of $W$ over $X$. By hypothesis, Chevalley's Theorem is true for subvarieties of $\mathbb{P}^{n-1} \times X$, and $\widetilde{H}$ is isomorphic to $\mathbb{P}^{n-1} \times X$. It follows that, for every component $W^{\prime}$ of $W$, the morphism $W^{\prime} \rightarrow X$ is a finite morphism, and therefore its image is closed in $X$ (Theorem 4.4.3). Thus the image $Z$ of $W$ is a closed subset of $X$, and it doesn't contain $p$. Then $X^{\prime}=X-Z$ is the required neighborhood of $p$.
figure: ??I'm not sure

### 4.8 Double Planes

dplane affdplanes

Let $A$ be the polynomial algebra $\mathbb{C}[x, y]$, and let $X$ be the affine plane Spec $A$. An affine double plane is a locus of the form $w^{2}=f(x, y)$ in affine 3 -space with coordinates $w, x, y$, where $f$ is a square-free polynomial in $x, y$, as in Example 4.3.8. Let $B=\mathbb{C}[w, x, y] /\left(w^{2}-f\right)$. So the affine double plane is $Y=\operatorname{Spec} B$.

We'll denote by $w, x, y$ both the variables and their residues in $B$. As in Example 4.3.8, $B$ is a normal domain of dimension two, and a free $A$-module with basis $(1, w)$. It has an automorphism $\sigma$ of order 2 , defined by $\sigma(a+b w)=a-b w$.

The fibres of $Y$ over $X$ are the $\sigma$-orbits in $Y$. If $f\left(x_{0}, y_{0}\right) \neq 0$, the fibre consists of two points, and if $f\left(x_{0}, y_{0}\right)=0$, it consists of one point. The reason that $Y$ is called a double plane is that most points of the plane $X$ are covered by two points of $Y$. The branch locus of the covering, which will be denoted by $\Delta$, is the (possibly reducible) curve $\{f=0\}$ in $X$. The fibres over the branch points, the points of $\Delta$, are single points.

We study the closed subvarieties $D$ of $Y$ that lie over a curve $C$ in $X$. These subvarieties will have dimension one, and we call them curves too. If $D$ lies over $C$, and if $D=D \sigma$, then $D$ is the only curve lying over $C$. Otherwise, there will be two curves that lie over $C$, namely $D$ and $D \sigma$. In that case we say that $C$ splits in $Y$.

A curve $C$ in $X$ will be the zero set of a principal prime ideal $P$ of $A$, and if $D$ lies over $C$, it will be the zero set of a prime ideal $Q$ of $B$ that lies over $P$ 4.4.1. However, the prime ideal $Q$ needn't be a principal ideal.
4.8.2. Example. Let $f(x, y)=x^{2}+y^{2}-1$. The double plane $Y=\left\{w^{2}=x^{2}+y^{2}-1\right\}$ is an affine quadric in $\mathbb{A}^{3}$. In the affine plane, its branch locus $\Delta$ is the curve $\left\{x^{2}+y^{2}=1\right\}$.

The line $C_{1}:\{y=0\}$ in $X$ meets the branch locus $\Delta$ transversally at the points $(x, y)=( \pm 1,0)$, and when we set $y=0$ in the equation for $Y$, we obtain the irreducible polynomial $w^{2}-x^{2}+1$. So $y$ generates a prime ideal of $B$. On the other hand, the line $C_{2}:\{y=1\}$ is tangent to $\Delta$ at the point $(0,1)$, and it splits. When we set $y=1$ in the equation for $Y$, we obtain $w^{2}=x^{2}$. The locus $\left\{w^{2}=x^{2}\right\}$ is the union of the two lines $\{w=x\}$ and $\{w=-x\}$ that lie over $C_{1}$. The prime ideals of $B$ that correspond to these lines aren't principal ideals.

## figure circle with two lines

This example illustrates a general principle: If a curve intersects the branch locus transversally, it doesn't split. We explain this now.

## (4.8.3) local analysis

Suppose that a plane curve $C:\{g=0\}$ and the branch locus $\Delta:\{f=0\}$ of a double plane $w^{2}=f$ meet at a point $p$. We adjust coordinates so that $p$ becomes the origin $(0,0)$, and we write

$$
f(x, y)=\sum a_{i j} x^{i} y^{j}=a_{10} x+a_{01} y+a_{20} x^{2}+\cdots
$$

Since $p$ is a point of $\Delta$, the constant coefficient of $f$ is zero. If the two linear coefficients aren't both zero, $p$ will be a smooth point of $\Delta$, and the tangent line to $\Delta$ at $p$ will be the line $\left\{a_{10} x+a_{01} y=0\right\}$. Similarly, writing $g(x, y)=\sum b_{i j} x^{i} y^{j}$, the tangent line to $C$, if defined, is the line $\left\{b_{10} x+b_{01} y=0\right\}$.

Let's suppose that the two tangent lines are defined and distinct - that $\Delta$ and $C$ intersect transversally at $p$. We change coordinates once more, to make the tangent lines the coordinate axes. After adjusting by scalar factors, the polynomials $f$ and $g$ will have the form

$$
f(x, y)=x+u(x, y) \quad \text { and } \quad g(x, y)=y+v(x, y)
$$

where $u$ and $v$ are polynomials all of whose terms have degree at least 2 .
Let $X_{1}=\operatorname{Spec} \mathbb{C}\left[x_{1}, y_{1}\right]$ be another affine plane. We consider the map $X_{1} \rightarrow X$ defined by the substitution $x_{1}=x+u(x, y), y_{1}=y+v(x, y)$. In the classical topology, this map is invertible analytically near the origin, because the Jacobian matrix

$$
\begin{equation*}
\left(\frac{\partial\left(x_{1}, y_{1}\right)}{\partial(x, y)}\right)_{(0,0)} \tag{4.8.4}
\end{equation*}
$$

at $p$ is the identity matrix. When we make this substitution, $\Delta$ becomes the locus $\left\{x_{1}=0\right\}$ and $C$ becomes the locus $\left\{y_{1}=0\right\}$. In this local analytic coordinate system, the equation $w^{2}=f$ that defines the double plane becomes $w^{2}=x_{1}$. When we restrict it to $C$ by setting $y_{1}=0, x_{1}$ becomes a local coordinate function on $C$. The restriction of the equation remains $w^{2}=x_{1}$. So the inverse image $Z$ of $C$ doesn't split analytically near $p$. Therefore it doesn't split algebraically either.

### 4.8.5. Corollary. A curve that meets the branch locus transversally at some point doesn't split.

This isn't a complete analysis. When $C$ and $\Delta$ are tangent at every point of intersection, $C$ may split or not, and which possibility occurs cannot be decided locally in most cases. However, one case in which a local analysis suffices to decide splitting is that $C$ is a line. Let $t$ be a coordinate in a line $C$, so that $C \approx \operatorname{Spec} \mathbb{C}[t]$. Let's assume that $C$ does't intersect $\Delta$ at $t=\infty$. The restriction of the polynomial $f$ to $C$ will give us a polynomial $\bar{f}(t)$ in $t$. A root of $\bar{f}$ corresponds to an intersection of $C$ with $\Delta$, and a multiple root corresponds to an intersection at which $C$ and $\Delta$ are tangent, or at which $\Delta$ is singular. The line $C$ will split if and only if $\bar{f}$ is a square in $\mathbb{C}[t]$, and this will be true if and only if every root of $\bar{f}$ has even multiplicity.

A rational curve is a curve whose function field is a rational function field $\mathbb{C}(t)$ in one variable. One can make a similar analysis for any rational plane curve, a conic for example, but one needs to inspect its points at infinity and its singular points as well as its smooth points at finite distance.
circleexample
(4.8.6) projective double planes

Let $X$ be the projective plane $\mathbb{P}^{2}$, with coordinates $x_{0}, x_{1}, x_{2}$. A projective double plane is a locus of the form

$$
\begin{equation*}
y^{2}=f\left(x_{0}, x_{1}, x_{2}\right) \tag{4.8.7}
\end{equation*}
$$

where $f$ is a square-free, homogeneous polynomial of even degree $2 d$. To regard this as a homogeneous equation, we must assign weight $d$ to the variable $y$ (see 1.7.7. Then, since we have weighted variables, we must work in a weighted projective space $\mathbb{W} \mathbb{P}$ with coordinates $x_{0}, x_{1}, x_{2}, y$, where $x_{i}$ have weight 1 and $y$ has weight $d$. A point of this weighted space $\mathbb{W P}$ is represented by a nonzero vector ( $x_{0}, x_{1}, x_{2}, y$ ) with the equivalence relation that, for all $\lambda \neq 0,\left(x_{0}, x_{1}, x_{2}, y\right) \sim\left(\lambda x_{0}, \lambda x_{1}, \lambda x_{2}, \lambda^{d} y\right)$. The points of the projective double plane $Y$ are the points of $\mathbb{W P P}$ that solve the equation (4.8.7).

The projection $\mathbb{W} \mathbb{P} \rightarrow X$ that sends $(x, y)$ to $x$ is defined at all points except at $(0,0,0,1)$. If $(x, y)$ solves 4.8.7) and if $x=0$, then $y=0$ too. So $(0,0,0,1)$ isn't a point of $Y$. The projection is defined at all points of $Y$. The fibre of the morphism $Y \rightarrow X$ over a point $x$ consists of points $(x, y)$ and $(x,-y)$, which will be equal if and only if $x$ lies on the branch locus of the double plane, the (possibly reducible) plane curve $\Delta:\{f=0\}$ in $X$. The map $\sigma:(x, y) \rightsquigarrow(x,-y)$ is an automorphism of $Y$, and points of $X$ correspond bijectively to $\sigma$-orbits in $Y$.

Since the double plane $Y$ is embedded into a weighted projective space, it isn't presented to us as a projective variety in the usual sense. However, it can be embedded into a projective space in the following way: The projective plane $X$ can be embedded by a Veronese embedding of higher order, using as coordinates the monomials $m=\left(m_{1}, m_{2}, \ldots\right)$ of degree $d$ in the variables $x$. This embeds $X$ into a projective space $\mathbb{P}^{N}$ where $N=\binom{d+2}{2}-1$. When we add a coordinate $y$ of weight $d$, we obtain an embedding of the weighted projective space $\mathbb{W} \mathbb{P}$ into $\mathbb{P}^{N+1}$ that sends the point $(x, y)$ to $(m, y)$. The double plane can be realized as a projective variety by this embedding.

If $Y \rightarrow X$ is a projective double plane, then, as happens with affine double planes, a curve $C$ in $X$ may split in $Y$ or not. If $C$ has a transversal intersection with the branch locus $\Delta$, it will not split. On the other hand, if $C$ is a line, and if $C$ intersects the branch locus $\Delta$ with multiplicity 2 at every intersection point, it will split. For example, when the branch locus $\Delta$ is a generic quartic curve, the lines that split will be the bitangent lines to $\Delta$ (see Section 1.11).

## (4.8.8) homogenizing an affine double plane

To construct a projective double plane from an affine double plane, we write the affine double plane as

$$
\begin{equation*}
w^{2}=F\left(u_{1}, u_{2}\right) \tag{4.8.9}
\end{equation*}
$$

for some nonhomogeneous polynomial $F$. We suppose that $F$ has even degree $2 d$, and we homogenize $F$, setting $u_{i}=x_{i} / x_{0}$. We multiply both sides of this equation by $x_{0}^{2 d}$ and set $y=x_{0}^{d} w$. This produces an equation of the form 4.8.7, where $f$ is the homogenization of $F$.

If $F$ has odd degree $2 d-1$, one needs to multiply $F$ by $x_{0}$ in order to make the substitution $y=x_{0}^{d} w$ permissible. When we do this, the line at infinity $\left\{x_{0}=0\right\}$ becomes a part of the branch locus.
(4.8.10) cubic surfaces and quartic double planes

We use coordinates $x_{0}, x_{1}, x_{2}, z$ for the (unweighted) projective 3 -space $\mathbb{P}^{3}$ here, and $X$ will denote the projective $x$-plane $\mathbb{P}^{2}$. Let $\mathbb{P}^{3} \xrightarrow{\pi} X$ denote the projection that sends $(x, z)$ to $x$. It is defined at all points except at the center of projection $q=(0,0,0,1)$, and its fibres are the lines through $q$, with $q$ omitted.

Let $S$ be a cubic surface in $\mathbb{P}^{3}$, the locus of zeros of an irreducible homogeneous cubic polynomial $g(x, z)$. We'll denote the restriction of $\pi$ to $S$ by the same symbol $\pi$.

Let's suppose that $q$ is a point of $S$. Then the coefficient of $z^{3}$ in $g$ will be zero, and $g$ will be quadratic in $z: \quad g(x, z)=a z^{2}+b z+c$, where the coefficienta $a, b, c$ are homogeneous polynomials in $x$, of degrees $1,2,3$, respectively. The equation for $S$ becomes

$$
\begin{equation*}
a z^{2}+b z+c=0 \tag{4.8.11}
\end{equation*}
$$

The discriminant $f=b^{2}-4 a c$ of $g$ is a homogeneous polynomial of degree 4 in $x$. Let $Y$ be the projective double plane

$$
\begin{equation*}
y^{2}=b^{2}-4 a c \tag{4.8.12}
\end{equation*}
$$

We denote by $V$ the affine space of polynomials $a, b, c$ of degrees $1,2,3$ in $x$, and by $W$ the affine space of homogeneous quartic polynomials in $x$. Sending a polynomial $g$ 4.8.11) to its discriminant $f$ defines a morphism $V \xrightarrow{u} W$ 4.8.11.
4.8.13. Lemma. The image of the morphism $u$ contains all quartic polynomials $f$ such that the divisor $D: f=0$ has at least one bitangent line. Therefore the image of $u$ is dense in $W$.
proof. Given such a quartic polynomial $f$, let $a$ be a linear polynomial such that the line $\ell_{1}:\{a=0\}$ is a bitangent to $D:\{f=0\}$. Then, as noted above, $\ell_{1}$ splits in the double plane $y^{2}=f$. So $f$ is congruent to a square, modulo $a$. Let $b$ be a quadratic polynomial such that $f \equiv b^{2}$, modulo $a$. When we take this polynomial as $b$, we will have $f=b^{2}-4 a c$ for some cubic polynomial $c$.

Conversely, if $g(x, y)=a z^{2}+b z+c$, the line $\ell_{1}:\{a=0\}$ will be a bitangent to $D$ provides that the locus $b=0$ meets $\ell_{1}$ in two distinct points.

It follows from the lemma that, if $g(x, z)=a z^{2}+b z+c$ is a polynomial in which $a, b, c$ are generic homogeneous polynomials in $x$, of degrees $1,2,3$, respectively, the discriminant $b^{2}-4 a c$ will be a generic homogeneous quartic polynomial in $x$.

We suppose given a generic cubic surface $S: a z^{2}+b z+c=0$ and the generic double plane $Y: y^{2}=$ $b^{2}-4 a c$.

### 4.8.14. Theorem. A generic cubic surface $S$ in $\mathbb{P}^{3}$ contains precisely 27 lines.

This theorem follows from next lemma, which relates the 27 lines in $S$ to the 28 bitangents of the generic quartic curve $\Delta:\left\{b^{2}-4 a c=0\right\}$ in the plane $X$. (See 1.11.2.).
4.8.15. Lemma. Let $S$ be a generic cubic surface $a z^{2}+b z+c=0$, projected to the plane $X$ from a generic point $q$ of $S$, let $\Delta:\left\{b^{2}=4 a c=0\right\}$ be the discriminant curve, and let $Y$ be the double plane $y^{2}=b^{2}-4 a c$.
(i) The image $\ell$ in $X$ of a line $L$ in $S$ is a line $\ell$ in $X$ that is bitangent to the quartic curve $\Delta$. Distinct lines in $S$ have distinct images in $X$.
(ii) The line $\ell_{1}:\{a=0\}$ is a bitangent to $\Delta$, and it isn't the image of a line in $S$.
(iii) A bitangent $\ell$ to $\Delta$ that is distinct from $\ell_{1}$ is the image of a line in $S$.
proof. (i) The projection $\mathbb{P}^{3} \rightarrow X$ maps a line $L$ bijectively to a line $\ell$ unless $L$ contains the center of projection $q$, in which case its image will be a point. Because our cubic surface $S$ is generic, it contains finitely many lines (3.6). The generic point $q$ of $S$ won't lie on any of those lines. So the image of a line $L$ in $S$ will be a line $\ell$ in $X$.

A line $\ell$ in $X$ is defined by a homogeneous linear equation in the variables $x$. The same linear equation defines a plane $H$ in $\mathbb{P}^{3}$ that contains $q$, and the intersection $C=S \cap H$ will be a cubic curve in $H$. This curve is essentially the inverse image of $\ell$ via the projection $S \rightarrow X$, though the projection is undefined at $q$. Let's call it the inverse image anyway.

At least one of the irreducible components of $C$ contains $q$, and that component isn't a line. So if $C$ is reducible, it will be a union $Q \cup L$, where $Q$ is a conic that contains $q$ and $L$ is a line in $S$. Thus lines $L$ in $S$ correspond bijectively to lines in $X$ such that the corresponding cubic $C$ is reducible.
(ii) The inverse image $C$ in $S$ of the line $\ell_{1}$ is the locus $a=0$ and $a z^{2}+b z+c=0$, or equivalently, $a=0$ and $b z+c=0$.

Let's adjust coordinates so that $a$ becomes the polynomial $x_{0}$. The locus $\left\{x_{0}=0\right\}$ in $\mathbb{P}^{3}$ is the projective plane $P$ with coordinates $x_{1}, x_{2}, z$, and in $P, C$ is the locus $\bar{g}=0$, where $\bar{g}=\bar{b} z+\bar{c}, \bar{b}$ and $\bar{c}$ being the polynomials obtained from $b$ and $c$ by substituting $x_{0}=0$. In $P$, the point $q$ becomes $(0,0,1)$, and $C$ becomes the cubic curve $\bar{g}=0$. This cubic curve is singular at $q$ because $\bar{g}$ has no term of degree $>1$ in $z$. As we have noted, $C$ doesn't contain a line through $q$. Since its degree is three, $C$ must be an irreducible cubic curve. So $\ell_{1}$ doesn't split.
(iii) Referring to 4.8.11) and 4.8.12, the quadratic formula solves for $z$ in terms of $y$ whenever $a \neq 0$ :

$$
\begin{equation*}
z=\frac{-b+y}{2 a} \quad \text { or } \quad y=2 a z+b \tag{4.8.16}
\end{equation*}
$$

These equations define a bijection $S^{\prime} \longleftrightarrow Y^{\prime}$ between the open subsets $S^{\prime}$ and $Y^{\prime}$ of points of $S$ and $Y$ at which $a \neq 0$. A point of $Y$ at which $a \neq 0$ is one that isn't on the line $\ell_{1}$.

If a line $\ell$ in $X$ is distinct from the line $\ell_{1}$, the intersection $\ell \cap \ell_{1}$ will be a single point $p$. The bijection $S^{\prime} \longleftrightarrow Y^{\prime}$ will be defined at all points that lie over $\ell$ except at the finite set of points over $p$.

If $\ell$ is a bitangent line, it splits in $Y$, and therefore it splits in $S$ too. The cubic curve $C=S \cap H$ will be reducible. It will be the union of a line $L$ and a conic. So every bitangent line distinct from $\ell_{1}$ is the image of a unique line in $S$.

Summing up: The 27 bitangents distinct from the bitangent $\ell_{1}:\{a=0\}$ are images of lines in $S$, but $\ell_{1}$ is not the image of a line in $S$.

## Chapter 5 STRUCTURE OF VARIETIES IN THE ZARISKI TOPOLOGY

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| 5.1 | Modules - a review |
| :--- | :--- |
| 5.2 | Valuations |
| 5.3 | Smooth Curves |
| 5.4 | Constructible sets |
| 5.5 | Closed Sets |
| 5.6 | Fibred Products |
| 5.7 | Projective Varieties are Proper |
| 5.8 | Fibre Dimension |

The ultimate goal of this chapter is to show how algebraic curves control the geometry of higher dimensional varieties.

### 5.1 Modules - a review

snake

We start with a brief review of modules, omitting proofs.
(5.1.1) exact sequences

Let $R$ be a ring. A sequence

$$
\cdots \rightarrow V^{n-1} \xrightarrow{d^{n-1}} V^{n} \xrightarrow{d^{n}} V^{n+1} \xrightarrow{d^{n+1}} \cdots
$$

of homomorphisms of $R$-modules is exact if the image of $d^{k-1}$ is equal to the kernel of $d^{k}$. For example, to say that a sequence $0 \rightarrow V \xrightarrow{d} V^{\prime}$ is exact means that the map $d$ is injective, and to say that a sequence $V \xrightarrow{d} V^{\prime} \rightarrow 0$ is exact, means that $d$ is surjective. Any homomorphism $V \xrightarrow{d} V^{\prime}$ can be embedded into an exact sequence

$$
0 \rightarrow K \rightarrow V \xrightarrow{d} V^{\prime} \rightarrow C \rightarrow 0
$$

where $K$ and $C$ are the kernel and cokernel of $d$, respectively.
A short exact sequence is an exact sequence of the form

$$
0 \rightarrow V \xrightarrow{a} V^{\prime} \xrightarrow{b} V^{\prime \prime} \rightarrow 0 .
$$

To say that this sequence is exact means that the map $a$ is injective, and that $V^{\prime \prime}$ is isomorphic to the quotient module $V^{\prime} / a V$.
5.1.2. Proposition. (functorial property of the kernel and cokernel) Suppose given a (commutative) diagram of $R$-modules

whose rows are exact sequences. Let $K, K^{\prime}, K^{\prime \prime}$ and $C, C^{\prime}, C^{\prime \prime}$ denote the kernels and cokernels of $f, f^{\prime}$, and $f^{\prime \prime}$, respectively.
(i) (kernel is left exact) The kernels form an exact sequence $K \rightarrow K^{\prime} \rightarrow K^{\prime \prime}$. If $u$ is injective, the sequence $0 \rightarrow K \rightarrow K^{\prime} \rightarrow K^{\prime \prime}$ is exact.
(ii) (cokernel is right exact) The cokernels form an exact sequence $C \rightarrow C^{\prime} \rightarrow C^{\prime \prime}$. If $v$ is surjective, the sequence $C \rightarrow C^{\prime} \rightarrow C^{\prime \prime} \rightarrow 0$ is exact.
(iii) (Snake Lemma) There is a canonical homomorphism $K^{\prime \prime} \xrightarrow{d} C$ that combines with the above sequences to form an exact sequence

$$
K \rightarrow K^{\prime} \rightarrow K^{\prime \prime} \xrightarrow{d} C \rightarrow C^{\prime} \rightarrow C^{\prime \prime} .
$$

If $u$ is injective and/or $v$ is surjective, the sequence remains exact with zeros at the appropriate ends.

## (5.1.3) tensor products

Let $U$ and $V$ be modules over a ring $R$. The tensor product $U \otimes_{R} V$ is an $R$-module that is generated by elements $u \otimes v$ called tensors, one for each $u$ in $U$ and $v$ in $V$. Its elements are combinations of tensors with coefficients in $R$.

The defining relations among the tensors are the bilinear relations:
and

$$
\begin{gather*}
\left(u_{1}+u_{2}\right) \otimes v=u_{1} \otimes v+u_{2} \otimes v, \quad u \otimes\left(v_{1}+v_{2}\right)=u \otimes v_{1}+u \otimes v_{2}  \tag{5.1.4}\\
r(u \otimes v)=(r u) \otimes v=u \otimes(r v)
\end{gather*}
$$

for all $u$ in $U, v$ in $V$, and $r$ in $R$. The symbol $\otimes$ is used as a reminder that the tensors are to be manipulated using these relations.

One can absorb a coefficient from $R$ into either one of the factors of a tensor, so every element of $U \otimes_{R} V$ can be written as a finite sum $\sum u_{i} \otimes v_{i}$ with $u_{i}$ in $U$ and $v_{i}$ in $V$.
5.1.5. Example. If $U$ and $V$ are free $R$-modules with bases $\left\{u_{i}\right\}$ and $\left\{v_{j}\right\}$, respectively, then $U \otimes_{R} V$ is a free $R$-module with basis $\left\{u_{i} \otimes v_{j}\right\}$. If $U$ is the space of $m$ dimensional (complex) column vectors, and $V$ is the space of $n$-dimensional row vectors. Then $U \otimes_{\mathbb{C}} V$ identifies naturally with the space of $m \times n$-matrices.

There is an obvious map of sets $U \times V \xrightarrow{\beta} U \otimes_{R} V$ from the product set to the tensor product, that sends $(u, v)$ to $u \otimes v$. This map isn't a homomorphism. The defining relations 5.1.4 show that it is $R$-bilinear, not $R$-linear. It is a universal bilinear map.

The product module $U \times V$ and the tensor product module $U \otimes_{R} V$ are very different. For instance, when $U$ and $V$ are free modules of ranks $r$ and $s, U \times V$ is free of rank $r+s$, while $U \otimes_{R} V$ is free of rank $r s$.
5.1.6. Corollary.et $U, V$, and $W$ be $R$-modules. Homomorphisms of $R$-modules $U \otimes_{R} V \rightarrow W$ correspond bijectively to $R$-bilinear maps $U \times V \rightarrow W$.

This follows from the defining relations.
Thus, any $R$-bilinear map $U \times V \xrightarrow{f} W$ to a module $W$ can be obtained from a module homomorphism $U \otimes_{R} V \xrightarrow{\widetilde{f}} W$ by composition with the bilinear map $\beta$ defined above: $U \times V \xrightarrow{\beta} U \otimes_{R} V \xrightarrow{\widetilde{f}} W$.
5.1.7. Proposition. There are canonical isomorphisms

- $U \otimes_{R} R \approx U$, defined by $u \otimes r \leftrightarrow u r$
- $\left(U \oplus U^{\prime}\right) \otimes_{R} V \approx\left(U \otimes_{R} V\right) \oplus\left(U^{\prime} \otimes_{R} V\right)$, defined by $\left(u_{1}+u_{2}\right) \otimes v \leadsto u_{1} \otimes v+u_{2} \otimes v$
- $U \otimes_{R} V \approx V \otimes_{R} U$, defined by $u \otimes v \nrightarrow v \otimes u$
- $\left(U \otimes_{R} V\right) \otimes_{R} W \approx U \otimes_{R}\left(V \otimes_{R} W\right)$, defined by $(u \otimes v) \otimes w \leftrightarrow u \otimes(v \otimes w)$
5.1.8. Proposition. Tensor product is right exact Let $U \xrightarrow{f} U^{\prime} \xrightarrow{g} U^{\prime \prime} \rightarrow 0$ be an exact sequence of $R$-modules. For any $R$-module $V$, the sequence

$$
U \otimes_{R} V \xrightarrow{f \otimes i d} U^{\prime} \otimes_{R} V \xrightarrow{g \otimes i d} U^{\prime \prime} \otimes_{R} V \rightarrow 0
$$

is exact.

Tensor product isn't left exact. For example, Let $R=\mathbb{C}[x]$. Then $R / x R \approx \mathbb{C}$, so there is an exact sequence $0 \rightarrow R \xrightarrow{x} R \rightarrow \mathbb{C} \rightarrow 0$. When we tensor with $\mathbb{C}$ we get the sequence $0 \rightarrow \mathbb{C} \xrightarrow{0} \mathbb{C} \rightarrow \mathbb{C} \rightarrow 0$, which isn't exact on the left.
proof of Proposition 5.1.8. We are given an exact sequence of $R$-modules $U \xrightarrow{f} U^{\prime} \xrightarrow{g} U^{\prime \prime} \rightarrow 0$ and another $R$-module $V$. We are to prove that the sequence $U \otimes_{R} V \xrightarrow{f \otimes 1} U^{\prime} \otimes_{R} V \xrightarrow{g \otimes 1} U^{\prime \prime} \otimes_{R} V \rightarrow 0$ is exact. It is evident that the composition $(g \otimes 1)(f \otimes 1)$ is zero, and that $g \otimes 1$ is surjective. We must prove that $U^{\prime \prime} \otimes_{R} V$ is the cokernel of $f \otimes_{R} V$. Let $C$ denote that cokernel. Then we have a canonical map $U^{\prime \prime} \otimes_{R} V \xrightarrow{\varphi} C$ that we want to show is an isomorphism. To do this, we construct its inverse function, first defining the inverse on tensors. We form a diagram

in which the rows are exact sequences, and the first two vertical arrows are the canonial bilinear maps. The arrow labelled ?? is defined by lifting an element $u^{\prime \prime}, v$ to $U^{\prime} \times V$. It is bilinear, so it induces a map $U^{\prime \prime} \otimes_{R} V \rightarrow C$ that inverts $\varphi$.
5.1.9. Corollary. Let $U$ and $V$ be modules over a domain $R$ and let s be a nonzero element of $R$. Let $R_{s}, U_{s}, V_{s}$ be the localizations of $R, U, V$, respectively.
(i) There is a canonical isomorphism $U \otimes_{R}\left(R_{s}\right) \approx U_{s}$.
(ii) Tensor product is compatible with localization: $U_{s} \otimes_{R_{s}} V_{s} \approx\left(U \otimes_{R} V\right)_{s}$
proof of Corollary 5.1.9((ii).????
Thus the tensor product $\mathcal{M} \otimes_{\mathcal{O}} \mathcal{N}$ of $\mathcal{O}$-modules $\mathcal{M}, \mathcal{N}$ on a variety is defined. The modules of sections of $\mathcal{M} \otimes_{\mathcal{O}} \mathcal{N}$ on an affine open set $U$ is $\mathcal{M}(U) \otimes_{\mathcal{O}}(U) \mathcal{N}(U)$.

## (5.1.10) extension of scalars in a module

Let $R \xrightarrow{\rho} R^{\prime}$ be a ring homomorphism. Extension of scalars is an operation that constructs an $R^{\prime}$-module from an $R$-module.

Let's write scalar multiplication on the right. So $M$ will be a right $R$-module. Then $M \otimes_{R} R^{\prime}$ becomes a right $R^{\prime}$-module, multiplication by $s \in R^{\prime}$ being $(m \otimes a) s=m \otimes(a s)$. This gives the functor

$$
R-\text { modules } \xrightarrow{\otimes R^{\prime}} R^{\prime}-\text { modules }
$$

called the extension of scalars from $R$ to $R^{\prime}$.
(5.1.11) localization, again
\#\#\#where???\#\#\#
If $s$ is a nonzero element of a domain $A$, the simple localization $A_{s}$, which is often referred to simply as a localization, is the ring obtained by adjoining an inverse of $s$. To work with the inverses of finitely many nonzero elements, one may simply adjoin the inverse of their product. For working with an infinite set of inverses, the concept of a multiplicative system is convenient.

A multiplicative system $S$ in a domain $A$ is a subset that consists of nonzero elements, is closed under multiplication, and contains 1 . If $S$ is a multiplicative system, the ring of $S$-fractions $A S^{-1}$ is the ring obtained by adjoining inverses of all elements of $S$. Its elements are equivalence classes of fractions $a s^{-1}$ with $a$ in $A$ and $s$ in $S$, the equivalence relation and the laws of composition being the usual ones for fractions. The ring $A S^{-1}$ called a localization too.
5.1.12. Examples. (i) The set consisting of the powers of a nonzero element $s$ of a domain $A$ is a multiplicative system. Its ring of fractions is the simple localization $A_{s}=A\left[s^{-1}\right]$.
(ii) The set $S$ of all nonzero elements of a domain $A$ is a multiplicative system. Its ring of fractions is the field of fractions of $A$.
(iii) An ideal $P$ of a domain $A$ is a prime ideal if and only if its complement, the set of elements of $A$ not in $P$, is a multiplicative system.

Let $A \subset B$ be a ring extension, and let $I$ and $J$ be ideals of $A$ and $B$, respectively. Recall that the extension of $I$ is the ideal $I B$ of $B$ generated by $I$, whose elements are finite sums $\sum_{i} z_{i} b_{i}$ with $z_{i}$ in $I$ and $b_{i}$ in $B$. The contraction of $J$ is the intersection $J \cap A$, which is an ideal of $A$.
5.1.13. Proposition. Let $S$ be a multiplicative system in a domain $A$, and let $A^{\prime}$ be the localization $A S^{-1}$.
(i) Let $I$ be an ideal of $A$. The extended ideal $I A^{\prime}$ is the set $I S^{-1}$ whose elements are classes of fractions $x s^{-1}$, with $x$ in $I$ and $s$ in $S$. The extended ideal is the unit ideal if and only if $I$ contains an element of $S$.
(ii) Let $J$ be an ideal of the localization $A^{\prime}$ and let $I$ denote its contraction $J \cap A$. The extended ideal $I A^{\prime}=(J \cap A) A^{\prime}$ is equal to $J$.
(iii) If $Q$ is a prime ideal of $A$ and if $Q \cap S$ is empty, the extended ideal $Q^{\prime}=Q A^{\prime}$ is a prime ideal of $A^{\prime}$, and the contraction $Q^{\prime} \cap A$ is equal to $Q$. If $Q \cap S$ isn't empty, the extended ideal is the unit ideal. Thus prime ideals of $A S^{-1}$ correspond bijectively to prime ideals of $A$ that don't meet $S$.
5.1.14. Corollary. Every localization $A S^{-1}$ of a noetherian domain $A$ is noetherian.

Tis follows from (5.1.13) (ii).

## (5.1.15) a general principle

An important, though elementary, principle for working with fractions is that any finite sequence of computations in a localization $A S^{-1}$ will involve only finitely many denominators, and can therefore be done in a simple localization $A_{s}$, where $s$ is a common denominator for the fractions that occur.

For example, let $A \subset B$ be finite-type domains, and let $S$ be the multiplicative system of nonzero elemets of $A$. Then $A S^{-1}=K$ is the field of fractions of $A$, and $B_{K}=B S^{-1}$ is a finite-type $K$-algebra. The Noether Normalization Theorem tells us that $B_{K}$ is a finite module over a polynomial subring $K\left[y_{1}, \ldots, y_{n}\right]$, There is a nonzero element $s$ in $A$ such that $B_{s}$ is a finite module over the polynomial ring $A_{s}\left[y_{1}, \ldots, y_{n}\right]$.

## (5.1.16) restriction of scalars

## \#\#\# check Chapter 6\#\#\#

If $A \xrightarrow{\rho} B$ is a ring homomorphism, a $B$-module $N$ can be made into an $A$-module by restriction of scalars, scalar multiplication by an element $a$ of $A$ being defined by the formula

$$
\begin{equation*}
a n=\rho(a) n \tag{5.1.17}
\end{equation*}
$$

It is customary to denote a module and the one obtained by restriction of scalars by the same symbol. But if it seems necessary in order to avoid confusion, we may denote a $B$-module $N$ and the $A$-module obtained from it by restriction of scalars by ${ }_{B} N$ and $A_{A} N$, respectively.

## (5.1.18) local rings

A local ring is a noetherian ring that contains just one maximal ideal. We make a few general comments about local rings here though we will be interested mainly in some special ones, the discrete valuation rings that are discussed below.

Let $R$ be a local ring with maximal ideal $M$. The quotient $R / M=k$ is a field that is called the residue field of $R$. . In most of the cases we study, the residue field $k$ will be the field of complex numbers.

An element of $R$ that isn't in $M$ isn't in any maximal ideal, so it is a unit.
The Nakayama Lemma 4.1.3 has a useful version for local rings:
localnakayama
5.1.19. Local Nakayama Lemma. Let $R$ be a local ring with maximal ideal $M$ and residue field $k$. Let $V$ be a finite $R$-module, and let $\bar{V}$ be the quotient $V / M V$, which an also be written as the tensor product $V \otimes_{A} k$. If $\bar{V}=0$, then $V=0$.
proof. If $\bar{V}=0$, then $V=M V$. The usual Nakayama Lemma tells us that $M$ contains an element $z$ such that $1-z$ annihilates $V$. Then $1-z$ isn't in $M$, so it is a unit. A unit annihilates $V$, and therefore $V=0$.

### 5.2 Valuations

A local domain $R$ with maximal ideal $M$ has dimension one if (0) and $M$ are the only prime ideals of $R$, and $(0) \neq M$. In this section, we describe the normal local domains of dimension one. They are the discrete valuation rings that are defined below.

Let $K$ be a field. A discrete valuation v on $K$ is a surjective homomorphism

$$
\begin{equation*}
K^{\times} \xrightarrow{\mathrm{v}} \mathbb{Z}^{+} \tag{5.2.1}
\end{equation*}
$$

from the multiplicative group of nonzero elements of $K$ to the additive group of integers such that, if $a, b$ are elements of $K$ and if $a, b$ and $a+b$ aren't zero, then

- $\mathrm{v}(a+b) \geq \min \{\mathrm{v}(a), \mathrm{v}(b)\}$.

The word "discrete" refers to the fact that $\mathbb{Z}^{+}$has the discrete topology. Other valuations exist. They are interesting, but less important, and we won't use them. To simplify terminology, we refer to a discrete valuation simply as a valuation.

Let $k$ be a positive integer. If v is a valuation and if $\mathrm{v}(a)=k$, then $k$ is called the order of zero of $a$, and if $\mathrm{v}(a)=-k$, then $k$ is called the order of pole of $a$, with respect to the valuation.
5.2.2. Lemma. Let v be a valuation on a field $K$ that contains the complex numbers. Every nonzero complex number has value zero.
proof. This is true because $\mathbb{C}$ contains $n$th roots. If $\gamma$ is an $n$th root of a nonzero complex number $c$, then because v is a homomorphism, $\mathrm{v}(\gamma)=\mathrm{v}(c) / n$. The only integer that is divisible by every positive integer $n$ is zero.

The valuation ring $R$ associated to a valuation v on a field $K$ is the subring of elements of $K$ with nonnegative values, together with zero:

$$
\begin{equation*}
R=\left\{a \in K^{\times} \mid \mathrm{v}(a) \geq 0\right\} \cup\{0\} . \tag{5.2.3}
\end{equation*}
$$

Valuation rings are usually called "discrete valuation rings", but since we have dropped the word discrete from the valuation, we drop it from the valuation ring too.
5.2.4. Proposition. Valuations of the field $\mathbb{C}(t)$ of rational functions in one variable correspond bijectively to points of the projective line $\mathbb{P}_{t}^{1}$. The valuation ring that corresponds to a point $p \neq \infty$ is the local ring of the polynomial ring $\mathbb{C}[t]$ at $p$.
beginning of the proof. Let $K$ denote the field $\mathbb{C}(t)$, and let $a$ be a complex number. To define the valuation v that corresponds to the point $p: t=a$ of $\mathbb{P}^{1}$, we write a nonzero polynomial $f$ as $(t-a)^{k} h$, where $t-a$ doesn't divide $h$, and we define, $\mathrm{v}(f)=k$. Then we define $\mathrm{v}(f / g)=\mathrm{v}(f)-\mathrm{v}(g)$. You will be able to check that with this definition, v becomes a valuation whose valuation ring is the algebra of regular functions at $p$ 2.7.1. This algebra is called the local ring at $p$ (see 5.2.8 below). The valuation that corresponds to the point at infinity of $\mathbb{P}^{1}$ is obtained by working with $t^{-1}$ in place of $t$.

The proof that these are all of the valuations of $\mathbb{C}(t)$ will be given at the end of the section.
5.2.5. Proposition. Let v be a valuation on a field $K$, let $R$ be its valuation ring, and let $x$ be a nonzero element of $K$ with value $\mathrm{v}(x)=1$.
idealsinvalring
(i) The ring $R$ is a normal local domain of dimension one. Its maximal ideal $M$ is the principal ideal $x R$. The elements of $M$ are the elements of $K$ with positive value, together with zero:

$$
M=\left\{a \in K^{\times} \mid \mathrm{v}(a)>0\right\} \cup\{0\}
$$

(ii) The units of $R$ are the elements of $K^{\times}$with value zero. Every element $z$ of the multiplicative group of $K^{\times}$ has the form $z=x^{k} u$, where $u$ is a unit and $k=\mathrm{v}(z)$ can be any integer.
(iii) The proper $R$-submodules of $K$ are the sets $x^{k} R$, where $k$ is a positive or negative integer. The set $x^{k} R$ consists of zero and the elements of $K^{\times}$with value $\geq k$. The nonzero ideals of $R$ are the principal ideals $x^{k} R$ with $k \geq 0$. They are the powers of the maximal ideal.
(iv) There is no ring properly between $R$ and $K$ : If $R^{\prime}$ is a ring and if $R \subset R^{\prime} \subset K$, then either $R=R^{\prime}$ or $R^{\prime}=K$.
proof. We prove (i) last.
(ii) Since v is a homomorphism, $\mathrm{v}\left(u^{-1}\right)=-\mathrm{v}(u)$. So $u$ is a unit of $R$, i.e., $u$ and $u^{-1}$ are both in $R$, if and only if $\mathrm{v}(u)=0$. If $z$ is a nonzero element of $K$ with $\mathrm{v}(z)=k$, then $u=x^{-k} z$ has value zero, so it is a unit, and $z=x^{k} u$.
(iii) The $R$-module $x^{k} R$ consists of the elements of $K$ of value at least $k$. Suppose that a nonzero $R$-submodule $N$ of $K$ contains an element $z$ with value $k$. Then $z=x^{k} u$, where $u$ is a unit, and therefore $N$ contains $x^{k}$, and $x^{k} R \subset N$. If $k$ is the smallest integer such that $N$ contains an element $z$ with value $k$, then $N=x^{k} R$. If there is no minimum value of the elements of $N$, then $N$ contains $x^{k} R$ for every $k$, and $N=K$.
(iv) This follows from (iii). The ring $R^{\prime}$ will be a nonzero $R$-submodule of $K$. since $R^{\prime} \neq K, \quad R^{\prime}=x^{k} R$ for some $k$, and if $R^{\prime}$ contains $R, k \leq 0$. If $k<0$ then $x^{k} R$ isn't closed under multiplication. So $k=0$ and $R^{\prime}=R$.
(i) First, $R$ is noetherian because (iii) tells us that it is a principal ideal domain, and it follows from (ii) that the only prime ideals of $R$ are $\{0\}$ and $M=x R$, where $x$ is an element with value 1 . So $R$ is a local ring of dimension 1 . If the normalization of $R$ were larger than $R$, then according to (iv), it would be equal to $K$. Then $x^{-1}$ would be integral over $R$. It would satisfy a polynomial relation $x^{-r}+a_{1} x^{-(r-1)}+\cdots+a_{r}=0$ with $a_{i}$ in $R$. When one multiplies this relation by $x^{r}$, one sees that 1 would be a multiple of $x$. Then $x$ would be a unit, which it is not.

### 5.2.6. Theorem.

(i) A local domain whose maximal ideal is a nonzero principal ideal is a valuation ring.
(ii) Every normal local domain of dimension 1 is a valuation ring.
proof. (i) Let $R$ be a local domain whose maximal ideal $M$ is a nonzero principal ideal, say $M=x R$, with $x \neq 0$, and let $y$ be a nonzero element of $R$. The integers $k$ such that $x^{k}$ divides $y$ are bounded 4.1.6. Let $x^{k}$ be the largest power that divides $y$. Then $y=u x^{k}$, where $k \geq 0$ and $u$ isn't in $M$. It is a unit. Then any nonzero element $z$ of the fraction field $K$ of $R$ will have the form $z=u x^{r}$ where $u$ is a unit and $r$ is an integer, possibly negative. This is shown by writing the numerator and denominator of a fraction in such a form and dividing.

The valuation whose valuation ring is $R$ is defined by $\mathrm{v}(z)=r$ when $z=u x^{r}$ as above. If $z_{i}=u_{i} x^{r_{i}}$, $i=1,2$, where $u_{i}$ are units and $r_{1} \leq r_{2}$, then $z_{1}+z_{2}=\alpha x^{r_{1}}$, where $\alpha=u_{1}+u_{2} x^{r_{2}-r_{1}}$ is an element of $R$. Therefore $\mathrm{v}\left(z_{1}+z_{2}\right) \geq r_{1}=\min \left\{\mathrm{v}\left(z_{1}\right), \mathrm{v}\left(z_{2}\right)\right\}$. We also have $\mathrm{v}\left(z_{1} z_{2}\right)=\mathrm{v}\left(z_{1}\right)+\mathrm{v}\left(z_{2}\right)$. Thus v is a surjective homomorphism. The requirements for a valuation are satisfied.
(ii) The fact that a valuation ring is a normal, one-dimensional local ring is Proposition 5.2 .5 (i). We show that a normal local domain $R$ of dimension 1 is a valuation ring by showing that its maximal ideal is a principal ideal. The proof is tricky.

Let $z$ be a nonzero element of $M$. Because $R$ is a local ring of dimension $1, M$ is the only prime ideal that contains $z$, so $M$ is the radical of the principal ideal $z R$, and $M^{r} \subset z R$ if $r$ is large. Let $r$ be the smallest integer such that $M^{r} \subset z R$. Then there is an element $y$ in $M^{r-1}$ that isn't in $z R$, but such that $y M \subset z R$.
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We restate this by saying that $w=y / z$ isn't in $R$, but $w M \subset R$. Since $M$ is an ideal, multiplication by an element of $R$ carries $w M$ to $w M$. So $w M$ is an ideal of $R$. Since $M$ is the maximal ideal of the local ring $R$, either $w M \subset M$, or $w M=R$. If $w M \subset M$, the lemma below shows that $w$ is integral over $R$. This can't happen because $R$ is normal and $w$ isn't in $R$. Therefore $w M=R$ and $M=w^{-1} R$. This implies that $w^{-1}$ is in $R$ and that $M$ is a principal ideal.
5.2.7. Lemma. Let I be a nonzero ideal of a noetherian domain $A$, and let $B$ be a domain that contains $A$. An element $w$ of $B$ such that $w I \subset I$ is integral over $A$.
proof. This is the Nakayama Lemma once more. Because $A$ is noetherian, $I$ is finitely generated. Let $v=$ $\left(v_{1}, \ldots, v_{n}\right)^{t}$ be a vector whose entries generate $I$. The hypothesis $w I \subset I$ allows us to write $w v_{i}=\sum p_{i j} v_{j}$ with $p_{i j}$ in $A$, or in matrix notation, $w v=P v$. So $w$ is an eigenvalue of $P$. If $p(t)$ denotes the characteristic polynomial of $P, p(w) v=0$. Since $I \neq 0$, at least one $v_{i}$ is nonzero. Since $A$ is a domain, $p(w) v_{i}=0$ implies that $p(w)=0$. The characteristic polynomial is a monic polynomial with coefficients in $A$, so $w$ is integral over $A$.

## (5.2.8) the local ring at a point

Let $\mathfrak{m}$ be the maximal ideal at a point $p$ of an affine variety $X=\operatorname{Spec} A$, and let $S$ be the complement of $\mathfrak{m}$ in $A$. This is a multiplicative system, and the prime ideals $P$ of the localization $A S^{-1}$ are extensions of the prime ideals $Q$ of $A$ that are contained in $\mathfrak{m}$ : $\quad P=Q S^{-1} 5.1 .13$. Thus $A S^{-1}$ is a local ring whose maximal ideal is $\mathfrak{m} S^{-1}$. This ring is called the local ring of $A$ at $p$, and is often denoted by $A_{p}$.

Any finite set $\alpha_{1}, \ldots, \alpha_{k}$ of elements of the local ring $A_{p}$ at $p$ will be contained in a simple localization $A_{s}$, for some $s$ in $S$. So $A_{p}$ will be in the coordinate algebra of some affine open neighborhood $X_{s}$ of $p$.

For example, let $X=\operatorname{Spec} A$ be the affine line, $A=\mathbb{C}[t]$, and let $p$ be the point $t=0$. The local ring $A_{p}$ is the ring whose elements are fractions $f(t) / g(t)$ with $g(0) \neq 0$.
5.2.9. Lemma. A rational function $\alpha$ on a variety $X$ is regular on $X$ if it is in the local ring of $X$ at every point $p$.

This is true because a function $\alpha$ is in the local ring at $p$ if and only if it is in the coordinate algebra of some affine neighborhood of $p$.
5.2.10. Corollary. Let $X=\operatorname{Spec} A$ be an affine variety.
(i) The coordinate algebra $A$ is the intersection of the local rings $A_{p}$ at points of $X$.

$$
A=\bigcap_{p \in X} A_{p}
$$

(ii) The coordinate algebra $A$ is normal if and only if all of its local rings $A_{p}$ are normal.

See Lemma 4.3 .3 for (ii).
5.2.11. Note. (about the overused word local) A property is true locally on a topological space $X$ if every point $p$ of $X$ has an open neighborhood $U$ such that the property is true on $U$.

In these notes, the words localize and localization refer to the process of adjoining inverses. The localizations $X_{s}$ of an affine variety $X=\operatorname{Spec} A$ form a basis for the topology on $X$. So if some property is true locally on $X$, one can cover $X$ by localizations on which the property is true. There will be elements $s_{1}, \ldots, s_{k}$ of $A$ that generate the unit ideal, such that the property is true on each of the localizations $X_{s_{i}}$.

An $A$-module $M$ is locally free if there are elements $s_{1}, \ldots, s_{k}$ that generate the unit ideal of $A$, such that $M_{s_{i}}$ is a free $A_{s_{i}}$-module for each $i$. The free modules $M_{s_{i}}$ will have the same rank. That rank is the rank of the locally free $A$-module $M$.

An ideal $I$ of $A$ is locally principal if there are elements $s_{i}$ that generate the unit ideal, such that, for every $i, I_{s_{i}}$ is a principal ideal of $A_{s_{i}}$. A locally principal ideal is a locally free module of rank one.
5.2.12. Corollary. Let $M$ be a finite module over a finite-type domain A. If for some point $p$ of $X=\operatorname{Spec} A$ the localized module $M_{p}$ (5.7) is a free module, there is an element s not in $\mathfrak{m}_{p}$ such that $M_{s}$ is free.
proof. This is an example of the general principle (5.1.15).
We finish the proof of Proposition 5.2.4 now, by showing that every valuation v of the function field $K=\mathbb{C}(t)$ of $\mathbb{P}^{1}$ corresponds to a point of $\mathbb{P}^{1}$.

Let $R$ be the valuation ring of v . If $\mathrm{v}(t)<0$, we replace $t$ by $t^{-1}$. So we may assume that $\mathrm{v}(t) \geq 0$. Then $t$ is an element of $R$, and therefore $\mathbb{C}[t] \subset R$. The maximal ideal $M$ of $R$ isn't zero. It contains a nonzero element of $K$, a fraction $f / g$ of polynomials in $t$. The denominator $g$ is in $R$, so the ideal $M$ also contains the nonzero polynomial $f$. Since $M$ is a prime ideal, it contains a monic irreducible factor of $f$, which will have the form $t-a$ for some complex number $a$. When $c \neq a$, then the scalar $c-a$ isn't in $M$, so $t-c$ won't be in $M$. Since $R$ is a local ring, $t-c$ is a unit of $R$ for all $c \neq a$. The localization $R_{0}$ of $\mathbb{C}[t]$ at the point $t=a$ is a valuation ring that is contained in the valuation ring $R(5.2 .4)$. There is no ring properly containing $R_{0}$ except $K$, so $R_{0}=R$.

### 5.3 Smooth Curves

A curve is a variety of dimension 1 . The proper closed subsets of a curve are finite.
5.3.1. Definition. A point $p$ of a curve $X$ is a smooth point if the local ring at $p$ is a valuation ring. Otherwise, $p$ is a singular point. A curve $X$ is smooth if all of its points are smooth.

Let $p$ be a smooth point of a curve $X$, and let $\mathrm{v}_{p}$ be the corresponding valuation. As with any valuation, we say that a rational function $\alpha$ on $X$ has a zero of order $k>0$ at $p$ if $\mathrm{v}_{p}(\alpha)=k$, and that it has a pole of order $k$ at $p$ if $\mathrm{v}_{p}(\alpha)=-k$.
5.3.2. Lemma. (i) An affine curve $X$ is smooth if and only if its coordinate algebra is a normal domain.
(ii) A curve has finitely many singular points.
(iii) The normalization $X^{\#}$ of a curve $X$ is a smooth curve, and the (finite) morphism $X^{\#} \rightarrow X$ becomes an isomorphism when singular points of $X$ and their inverse images are deleted.
proof. (i) This follows from Theorem 5.2.6 and Proposition 4.3.3
(ii),(iii) Any nonempty open subset of a curve $X$ will be the complement of a finite set, so we may replace $X$ by an affine open subset, say $\operatorname{Spec} A$. The normalization $A^{\#}$ of $A$ will be a finite $A$-module, and therefore a finite-type algebra with the same fraction field as $A$, and $\operatorname{Spec} A^{\#}$ will be a smooth curve. Let $\alpha_{1}, \ldots, \alpha_{k}$ be generators for the finite $A$-module $A^{\#}$. They are elements of the fraction field $K$, and can be written as fractions $\alpha_{i}=a_{i} / s$ for some $s$ in $A$. The localizations $A_{s}$ and $A_{s}^{\#}$ are equal, so the open subset $X_{s}=\operatorname{Spec} A_{s}$ of $X$ will be smooth.
5.3.3. Proposition. Let $X$ be a smooth curve with function field $K$. Every point of $\mathbb{P}^{n}$ with values in $K$ defines a morphism $X \rightarrow \mathbb{P}^{n}$.
proof. Let $\mathrm{v}_{p}$ denote the valuation that corresponds to a point $p$ of $X$. A point $\left(\alpha_{0}, \ldots, \alpha_{n}\right)$ of $\mathbb{P}^{n}$ with values in $K$ determines a morphism $X \rightarrow \mathbb{P}^{n}$ if and only if, for every point $p$ of $X$, there is an index $j$ such that the functions $\alpha_{i} / \alpha_{j}$ are regular at $p$ for all $i=0, \ldots, n \sqrt{3.4 .12}$. This will be true when $j$ is chosen so that the order of zero $\mathrm{v}_{p}\left(\alpha_{j}\right)$ of $\alpha_{j}$ at $p$ is the minimal integer among the $\mathrm{v}_{p}\left(\alpha_{i}\right)$, for the indices $i$ such that $\alpha_{i} \neq 0$. $\square$

As the next example shows, this proposition cannot be extended to varieties $X$ of dimension greater than one.
5.3.4. Example. Let $Y$ be the complement of the origin in the affine plane $X=\operatorname{Spec} \mathbb{C}[x, y]$, and let $K=\mathbb{C}(x, y)$ be the function field of $X$. The vector $(x, y)$ defines a point of $\mathbb{P}_{x, y}^{1}$ with values in $K$. This point can be written as $(1, y / x)$ and also as $(x / y, 1)$. So $(x, y)$ defines a morphism to $\mathbb{P}^{1}$ wherever at least one of the functions $x / y$ or $y / x$ is regular, which is true at all points of $Y$. To extend the morphism to $X$, one would need an element $\lambda$ in $K$ such that $\lambda x$ and $\lambda y$ are both regular at $(0,0)$ and not both zero there. There is no such element, so the morphism doesn't extend to $X$.
5.3.5. Proposition. Let $X=$ Spec $A$ be a smooth affine curve with function field $K$. The local rings of $X$ are the valuation rings of $K$ that contain $A$. Therefore the maximal ideals of $A$ are locally principal.
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proof. Since $A$ is a normal domain of dimension one, its local rings are valuation rings that contain $A$ (Theorem 5.2.6. Let $R$ be a valuation ring of $K$ that contains $A$, let v be the associated valuation, and let $M$ be the maximal ideal of $R$. The intersection $M \cap A$ is a prime ideal of $A$. Since $A$ has dimension 1, the zero ideal is the only prime ideal of $A$ that isn't a maximal ideal. We can clear the denominator of an element of $M$, multiplying by an element of $R$, to obtain an element of $A$ while staying in $M$. So $M \cap A$ isn't the zero ideal. It is the maximal ideal $\mathfrak{m}_{p}$ of $A$ at a point $p$ of $X$. The elements of $A$ that aren't in $\mathfrak{m}_{p}$ aren't in $M$ either, so they are invertible in $R$. So the local ring $A_{p}$ at $p$, which is a valuation ring, is contained in $R$, and therefore it is equal to $R \quad 5.2 .5$ (iii).
5.3.6. Proposition. Let $X^{\prime}$ and $X$ be smooth curves with the same function field $K$.
(i) A morphism $X^{\prime} \xrightarrow{f} X$ that is the identity on the function field $K$ maps $X^{\prime}$ isomorphically to an open subvariety of $X$.
(ii) If $X$ is projective, $X^{\prime}$ is isomorphic to an open subvariety of $X$.
(iii) If $X^{\prime}$ and $X$ are both projective, they are isomorphic.
(iv) If $X$ is projective, every valuation ring of $K$ is the local ring at a point of $X$.
proof. (i) Let $p$ be the image in $X$ of a point $q$ of $X^{\prime}$, let $U$ be an affine open neighborhood of $p$, and let $V$ be an affine open neighborhood of $q$ in $X^{\prime}$ that is contained in the inverse image of $U$. Say $U=\operatorname{Spec} A$ and $V=\operatorname{Spec} B$. The morphism $f$ gives us a homomorphism $A \rightarrow B$, and since $q$ maps to $p$, this homomorphism extends to an inclusion of local rings $A_{p} \subset B_{q}$. They are valuation rings with the same field of fractions, so they are equal. Since $B$ is a finite-type algebra, there is an element $s$ in $A$, with $s(q) \neq 0$, such that $A_{s}=B_{s}$. Then the open subsets $\operatorname{Spec} A_{s}$ of $X$ and Spec $B_{s}$ of $X^{\prime}$ are the same. Since the point $q$ is arbitrary, $X^{\prime}$ is covered by open subvarieties of $X$. So $X^{\prime}$ is an open subvariety of $X$.
(ii) The projective embedding $X \subset \mathbb{P}^{n}$ is defined by a point $\left(\alpha_{0}, \ldots, \alpha_{n}\right)$ with values in $K$, and that same point defines a morphism $X^{\prime} \rightarrow \mathbb{P}^{n}$. If $f\left(x_{0}, \ldots, x_{n}\right)=0$ is a set of defining equations of $X$ in $\mathbb{P}^{n}$, then $f(\alpha)=0$ in $K$, and therefore $f$ vanishes on $X^{\prime}$ too. So the image of $X^{\prime}$ is contained in the zero locus of $f$, which is $X$. Then (i) shows that $X^{\prime}$ is an open subvariety of $X$.
(iii) This follows from (ii).
(iv) The local rings of $X$ are normal and of dimension one. They are valuation rings of $K$. We prove the converse. Let $\beta=\left(\beta_{0}, \ldots, \beta_{n}\right)$ be the point with values in $K$ that defines the projective embedding of $X$, let $R$ be a valuation ring of $K$, and let v be the corresponding valuation. We order the coordinates so that $\mathrm{v}\left(\beta_{0}\right)$ is minimal. Then the ratios $\gamma_{j}=\beta_{j} / \beta_{0}$ will be in $R$. The coordinate algebra $A_{0}$ of the affine variety $X^{0}=X \cap \mathbb{U}^{0}$ is generated by the coordinate functions $\gamma_{j}$, so $A_{0} \subset R$. Prposition 5.3.5 tells us that $R$ is the local ring of $X^{0}$ at some point.
5.3.7. Proposition. Let $X=\operatorname{Spec} A$ be an affine curve, let $\mathfrak{m}$ and v be the maximal ideal and valuation, respectively, at a smooth point $p$ of $X$, and let t be an element of $A$ with value $\mathrm{v}(t)=1$. The valuation ring $R$ of v is the local ring of $A$ at $p$. Let $M$ be its maximal ideal.
(i) The power $\mathfrak{m}^{k}$ consists of the elements of $A$ whose values are at least $k$. If I is an ideal of $A$ whose radical is $\mathfrak{m}$, then $I=\mathfrak{m}^{k}$ for some $k>0$.
(ii) The algebras $A / \mathfrak{m}^{n}$ and $R / M^{n}$ are isomorphic to the truncated polynomial ring $\mathbb{C}[t] /\left(t^{n}\right)$.
(iii) If $X$ is a smooth affine curve, every nonzero ideal $I$ of $A$ is a product $\mathfrak{m}_{1}^{e_{1}} \cdots \mathfrak{m}_{k}^{e_{k}}$ of powers of maximal ideals.
proof. (i) The nonzero ideals of $R$ are powers of $M$. So $M^{k}$ is the set of eleents of $R$ with value $\geq k$.
Let $I$ be an ideal of $A$ whose radical is $\mathfrak{m}$, and let $k$ be the minimal value $\mathrm{v}(x)$ of the nonzero elements $x$ of $I$. We will show that $I$ is the set of all elements of $A$ with value $\geq k$, i.e., that $I=M^{k} \cap A$. Since we can apply the same reasoning to $\mathfrak{m}^{k}$, it will follow that $I=\mathfrak{m}^{k}$.

We must show that if an element $y$ if $A$ has value $\mathrm{v}(y) \geq k$, then it is in $I$. Let $x$ be an element of $I$ with value $k$. Then $x$ divides $y$ in $R$, say $y / x=u$, with $u$ in $R$. The element $u$ will be a fraction $a / s$ with $s$ and $a$ in $A$, $s$ not in $\mathfrak{m}$, and $s y=a x$. The element $s$ will vanish at a finite set of points $q_{1}, \ldots, q_{r}$, but not at $p$. We choose an element $z$ of $A$ that vanishes at $p$ but not at any of the points $q_{1}, \ldots, q_{r}$. Then $z$ is in $\mathfrak{m}$, and since the radical of $I$ is $\mathfrak{m}$, some power of $z$ is in $I$. We replace $z$ by that power, so that $z$ is in $I$. By our choice, $z$ and
$s$ have no common zeros in $X$. They generate the unit ideal of $A$, say $1=c s+d z$ with $c$ and $d$ in $A$. Then $y=c s y+d z y=c a x+d z y$. Since $x$ and $z$ are in $I$, so is $y$.
(ii) Since $p$ is a smooth point, the local ring of $A$ at $p$ is the valuation ring $R$. Let $P$ be the subring $\mathbb{C}[t]$ of $A$, and let $\bar{P}_{k}=P /(t)^{k}, \bar{A}_{k}=A / \mathfrak{m}^{k}$, and $\bar{R}_{k}=R / M^{k}$. Since $\mathfrak{m}$ isn't the zero ideal, $\mathfrak{m}^{k-1}<\mathfrak{m}^{k}$ (Corollary 4.1.6. It follows from (i) that $t \mathfrak{m}^{k-1}=\mathfrak{m}^{k}$. Therefore $\mathfrak{m}^{k-1} / \mathfrak{m}^{k}$ has $\mathbb{C}$-dimension 1 . The map labelled $g_{k-1}$ in the diagram below is bijective.


By induction on $k$, we may assume that the map labelled $f_{k-1}$ is bijective, and then $f_{k}$ is bijective. The same argument shows that $\bar{P}_{k}$ and $\bar{R}_{k}$ are isomorphic
(iii) Let $I$ be a nonzero ideal of $A$. Because $X$ has dimension one, the locus of zeros of $I$ is a finite set $\left\{p_{1}, \ldots, p_{k}\right\}$. Therefore the radical of $I$ is the intersection $\mathfrak{m}_{1} \cap \cdots \cap \mathfrak{m}_{k}$ of the maximal ideals $\mathfrak{m}_{j}$ at $p_{j}$, which, by the Chinese Remainder Theorem, is the product ideal $\mathfrak{m}_{1} \cdots \mathfrak{m}_{k}$, Moreover, $I$ contains a power of that product, say $I \supset \mathfrak{m}_{1}^{N} \cdots \mathfrak{m}_{k}^{N}$. Let $J=\mathfrak{m}_{1}^{N} \cdots \mathfrak{m}_{k}^{N}$. The quotient algebra $A / J$ is the product $B_{1} \times \cdots \times B_{k}$, with $B_{j}=A / \mathfrak{m}_{j}^{N}$, and $A / I$ is a quotient of $A / J$. Proposition 2.1.7 tells us that $A / I$ is a product $\bar{A}_{1} \times \cdots \times \bar{A}_{k}$, where $\bar{B}_{j}$ is a quotient of $B_{j}$. By part (ii), each $B_{j}$ is a truncated polynomial ring. So the quotients $\bar{A}_{j}$ must also be truncated polynomial rings. Then the kernel $I$ of the map $A \rightarrow \bar{A}_{1} \times \cdots \times \bar{A}_{k}$ is a product of powers of the maximal ideals $\mathfrak{m}_{j}$.

## (5.3.8) isolated points

5.3.9. Proposition. In the classical topology, a curve, smooth or not, contains no isolated point. .

This was proved before for plane curves (Proposition 1.3.18).

### 5.3.10. Lemma.

(i) Let $Y^{\prime}$ be an open subvariety of a variety $Y$. A point $q$ of $Y^{\prime}$ is an isolated point of $Y$ if and only if it is an isolated point of $Y^{\prime}$.
(ii) Let $Y^{\prime} \xrightarrow{u^{\prime}} Y$ be a nonconstant morphism of curves, let $q^{\prime}$ be a point of $Y^{\prime}$, and let $q$ be its image in $Y$. If $q$ is an isolated point of $Y$, then $q^{\prime}$ is an isolated point of $Y^{\prime}$.
proof. (i) A point $q$ of $Y$ is isolated if $\{q\}$ is an open subset of $Y$. If $\{q\}$ is open in $Y^{\prime}$ and $Y^{\prime}$ is open in $Y$, then $\{q\}$ is open in $Y$, and if $\{q\}$ is open in $Y$, it is open in $Y^{\prime}$.
(ii) Because $Y^{\prime}$ has dimension one, the fibre over $q$ will be a finite set, say $\left\{q^{\prime}\right\} \cup F$, where $F$ is finite. Let $Y^{\prime \prime}$ denote the (open) complement $Y^{\prime}-F$ of $F$ in $Y^{\prime}$, and let $u^{\prime \prime}$ be the restriction of $u^{\prime}$ to $Y^{\prime \prime}$. The fibre of $Y^{\prime \prime}$ over $q$ is the point $q^{\prime}$. If $\{q\}$ is open in $Y$, then because $u^{\prime \prime}$ is continuous, $\left\{q^{\prime}\right\}$ will be open in $Y^{\prime \prime}$, and therefore open in $Y^{\prime}$.
proof of Proposition 5.3.9. Let $q$ be a point of a curve $Y$. Part (i) of Lemma 5.3.10 allows us to replace $Y$ by an affine neighborhood of $q$. Let $Y^{\prime}$ be the normalization of $Y$. Part (ii) of that lemma allows us to replace $Y$ by $Y^{\prime}$. So we may assume that $Y$ is a smooth affine curve, say $Y=\operatorname{Spec} B$. We can still replace $Y$ by an open neighborhood of $q$, so we may assume that the maximal ideal $\mathfrak{m}_{q}$ at $q$ is a principal ideal.

Say that $B=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] /\left(f_{1}, \ldots, f_{k}\right)$, that $q$ is the origin $(0, \ldots, 0)$ in $\mathbb{A}_{x}^{n}$, and that $\mathfrak{m}_{q}$ is generated by the residue of a polynomial $f_{0}$ in $B$. Then $f_{0}, \ldots, f_{k}$ generate the maximal ideal $M$ of $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ at the origin, which is also generated by $x_{1}, \ldots, x_{n}$. Let's write $f_{i}=\sum_{1}^{n} c_{i j} x_{j}+O(2)$, where $O(2)$ denotes an undetermined polynomial, all of whose terms have degree $\geq 2$ in $x$. The coefficient $c_{i j}$ is the partial derivative $\frac{\partial f_{i}}{\partial x_{j}}$, evaluated at $q$. So if $J$ denotes $(k+1) \times n$ Jacobian matrix $\left(\frac{\partial f_{i}}{\partial x_{j}}\right)$, then $\left(f_{0}, \ldots, f_{k}\right)^{t}=J\left(x_{1}, \ldots, x_{n}\right)^{t}+O(2)$. Since $f_{0}, \ldots, f_{k}$ generate $M$, there is a matrix $P$ with polynomial entries such that $P f^{t}=x^{t}$. Then $x^{t}=$ $P J x^{t}+O(2)$. If $P_{0}$ and $J_{0}$ are the constant terms of $P$ and $J, P_{0} J_{0}$ will be the identity matrix. So $J_{0}$ has rank $n$.

Let $J_{0}^{\prime}$ be the matrix obtained by deleting the column with index 0 from $J_{0}$. This matrix has rank at least $n-1$, and we may assume that the submatrix with indices $1 \leq i, j \leq n-1$ is invertible. The Implicit Function Theorem says that the equations $f_{1}, \ldots, f_{n-1}$ can be solved for the variables $x_{1}, \ldots, x_{n-1}$ as analytic functions of $x_{n}$, for small $x_{n}$. The locus $Z$ of zeros of $f_{1}, \ldots, f_{n-1}$ is locally homeomorphic to the affine line (1.4.18, and it contains $Y$. Since $Y$ has dimension 1, the component of $Z$ that contains $q$ must be equal to $Y$. So $Y$ is locally homeomorphic to $\mathbb{A}^{1}$, which has no isolated point. Therefore $q$ isn't an isolated point of $Y$.

### 5.4 Constructible Sets

In this section, $X$ will denote a noetherian topological space. Every closed subset of $X$ is a finite union irreducible closed sets 2.2.14).

The intersection $L=C \cap U$ of a closed set $C$ and an open set $U$ is called a locally closed set. For example, open sets and closed sets are locally closed.
5.4.1. Lemma. The following conditions on a subset $L$ of $A$ are equivalent.

- L is locally closed.
- L is a closed subset of an open subset $U$ of $X$.
- L is an open subset of a closed subset $C$ of $X$.

A constructible set is a set that is the union of finitely many locally closed sets.

### 5.4.2. Examples.

(i) A subset $S$ of a curve $X$ is constructible if and only if it is either a finite set or the complement of a finite set. Thus $S$ is constructible if and only if it is either closed or open.
(ii) Let $C$ be the line $\{y=0\}$ in the affine plane $X=\operatorname{Spec} \mathbb{C}[x, y]$, let $U=X-C$ be its open complement, and let $p=(0,0)$. The union $U \cup\{p\}$ is constructible, but not locally closed.

We will use the following notation: $L$ will denote a locally closed set, $C$ will denote a closed set, and $U$ willl denote an open set.
5.4.3. Theorem. The set $\mathbb{S}$ of constructible subsets of a noetherian topological space $X$ is the smallest family of subsets that contains the open sets and is closed under the three operations of finite union, finite intersection, and complementation.
proof. Let $\mathbb{S}_{1}$ denote the family of subsets obtained from the open sets by the three operations mentioned in the statement. Open sets are constructible, and using those three operations, one can make any constructible set from the open sets. So $\mathbb{S} \subset \mathbb{S}_{1}$. To show that $\mathbb{S}=\mathbb{S}_{1}$, we show that the family of constructible sets is closed under the three operations.

It is obvious that a finite union of constructible sets is constructible. The intersection of two locally closed sets $L_{1}=C_{1} \cap U_{1}$ and $L_{2}=C_{2} \cap U_{2}$ is locally closed because $L_{1} \cap L_{2}=\left(C_{1} \cap C_{2}\right) \cap\left(U_{1} \cap U_{2}\right)$. If $S=L_{1} \cup \cdots \cup L_{k}$ and $S^{\prime}=L_{1}^{\prime} \cup \cdots \cup L_{r}^{\prime}$ are constructible sets, the intersection $S \cap S^{\prime}$ is the union of the locally closed intersections $\left(L_{i} \cap L_{j}^{\prime}\right)$, so it is constructible.

Let $S$ be the constructible set $L_{1} \cup \cdots \cup L_{k}$. Its complement $S^{c}$ is the intersection of the complements $L_{i}^{c}$ of $L_{i}$ : $S^{c}=L_{1}^{c} \cap \cdots \cap L_{k}^{c}$. We have shown that intersections of constructible sets are constructible. So to show that the complement $S^{c}$ is constructible, it suffices to show that the complement of a locally closed set is constructible. Let $L$ be the locally closed set $C \cap U$, and let $C^{c}$ and $U^{c}$ be the complements of $C$ and $U$, respectively. Then $C^{c}$ is open and $U^{c}$ is closed. The complement $L^{c}$ of $L$ is the union $C^{c} \cup U^{c}$ of constructible sets, so it is constructible.
5.4.4. Proposition. In a noetherian topological space $X$, every constructible subset is a union $L_{1} \cup \cdots \cup L_{k}$ of locally closed sets: $L_{i}=C_{i} \cap U_{i}$, in which the closed sets $C_{i}$ are irreducible and distinct.
proof. Let $L=C \cap U$ be a locally closed set, and let $C=C_{1} \cup \cdots \cup C_{r}$ be the decomposition of $C$ into irreducible components. Then $L=\left(C_{1} \cap U\right) \cup \cdots \cup\left(C_{r} \cap U\right)$, which is constructible. So every constructible set $S$ is a union of locally closed sets $L_{i}=C_{i} \cap U_{i}$ in which the $C_{i}$ are irreducible. Next, suppose that two of the irreducible closed sets are equal, say $C_{1}=C_{2}$. Then $L_{1} \cup L_{2}=\left(C_{1} \cap U_{1}\right) \cup\left(C_{1} \cap U_{2}\right)=C_{1} \cap\left(U_{1} \cup U_{2}\right)$ is locally closed. So we can find an expression in which the closed sets are also distinct.

### 5.4.5. Lemma.

(i) Let $X_{1}$ be a closed subset of a variety $X$, and let $X_{2}$ be its open complement. A subset $S$ of $X$ is constructible if and only if $S \cap X_{1}$ and $S \cap X_{2}$ are constructible.
(ii) Let $X^{\prime}$ be an open or a closed subvariety of a variety $X$.
a) If $S$ is a constructible subset of $X$, then $S^{\prime}=S \cap X^{\prime}$ is a constructible subset of $X^{\prime}$.
b) A subset $S^{\prime}$ of $X^{\prime}$ is a constructible subset of $X^{\prime}$ if and onlt if it is a constructible subset of $X$.
proof. (i) This follows from Theorem 5.4.3.
(iia) It suffices to prove that the intersection $L^{\prime}=L \cap X^{\prime}$ of a locally closed subset $L$ of $X$ is a locally closed subset of $X^{\prime}$. If $L=C \cap U$, then $C^{\prime}=C \cap X^{\prime}$ is closed in $X^{\prime}$, and $U^{\prime}=U \cap X^{\prime}$ is open in $X^{\prime}$. So $L^{\prime}=C^{\prime} \cap U^{\prime}$ is locally closed.
(iib) It follows from (iia) that if a subset $S^{\prime}$ of $X^{\prime}$ is contructible in $X$, then it is constructible in $X^{\prime}$. To show that a constructible subset of $X^{\prime}$ is contructible in $X$, it suffices to show that a locally closed subset $L^{\prime}=C^{\prime} \cap U^{\prime}$ of $X^{\prime}$ is locally closed in $X$. If $X^{\prime}$ is closed in $X$, then $C^{\prime}$ is closed in $X$, and $U^{\prime}=X \cap U$ for some open subset $U$ of $X$. Since $C^{\prime} \subset X^{\prime}, L^{\prime}=C^{\prime} \cap U^{\prime}=C^{\prime} \cap X^{\prime} \cap U=C^{\prime} \cap U$, which is locally closed in $X$.

Suppose that $X^{\prime}$ is open in $X$. Then $U^{\prime}$ is open in $X$. If $C$ is the closure of $C^{\prime}$ in $X$, then $L^{\prime}=C \cap U^{\prime}=$ $C \cap X^{\prime} \cap U^{\prime}=C^{\prime} \cap U^{\prime}$. Again, $L^{\prime}$ is locally closed in $X$.

The next theorem illustrates a general fact, that sets arising in algebraic geometry are often constructible.
5.4.6. Theorem. Let $Y \xrightarrow{f} X$ be a morphism of varieties. The inverse image of a constructible subset of $X$ is a constructible subset of $Y$. The image of a constructible subset of $Y$ is a constructible subset of $X$.
proof. The fact that a morphism is continuous implies that the inverse image of a constructible set is constructible. To prove that the image of a constructible set is constructible, one keeps reducing the problem until there is nothing is left.

Let $S$ be a constructible subset of $Y$. Lemma 5.4.5 and Noetherian induction allow us to assume that the theorem is true when $S$ is contained in a proper closed subvariety of $Y$, and also when its image $f(S)$ is contained in a proper closed subvariety of $X$.

Suppose that $X$ is the union a proper closed subvariety $X_{1}$ and its open complement $X_{2}$. The inverse image $Y_{1}=f^{-1}\left(X_{1}\right)$ will be closed in $Y$, and its open complement will be the inverse image $Y_{2}=f^{-1}\left(X_{2}\right)$. A constructible subset $S$ of $Y$ is the union of the constructible sets $S_{1}=S \cap Y_{1}$ and $S_{2}=S \cap Y_{2}$, and $f(S)=f\left(S_{1}\right) \cup f\left(S_{2}\right)$. It suffices to show that $f\left(S_{1}\right)$ and $f\left(S_{2}\right)$ are constructible, and to show this, it suffices to show that $f\left(S_{i}\right)$ is a constructible subset of $X_{i}$ for $i=1,2(5.4 .5)$ (iib). Moreover, noetherian induction applies to $X_{1}$. So we need only show that $f\left(S_{2}\right)$ is a constructible subset of $X_{2}$. This means that we can replace $X$ by $X_{2}$, which is an arbitrary nonempty open subset, and $Y$ by its inverse image.

Next, suppose that $Y$ is the union of a proper closed subvariety $Y_{1}$ and its open complement $Y_{2}$, and let $S_{i}=S \cap Y_{i}$. It suffices to show that $S_{i}$ is a constructible subset of $Y_{i}, i=1,2$, and induction applies to $Y_{1}$. So we may replace $Y$ by any nonempty open subvariety.

Summing up, we can replace $X$ by any nonempty open subset $X^{\prime}$, and $Y$ by any nonempty open subset $Y^{\prime}$ that is contained in the inverse image of $X^{\prime}$. We can do this finitely often.

Since $S$ is a finite union of locally closed sets, it suffices to treat the case that $S$ is locally closed. Moreover, we may suppose that $S=C \cap U$, where $C$ is irreducible. Then $Y$ is the union of the closed subset $C=Y_{1}$ and its complement $Y_{2}$. Since $S \cap Y_{2}=\emptyset$, it suffices to treat $Y_{1}$. We may replace $Y$ by $C$. So we may assume that $S=Y \cap U=U$, and we may replace $Y$ by $U$. We are thus reduced to the case that $S=Y$.

We may still replace $X$ and $Y$ by nonempty open subsets, so we may assume that they are affine, say $Y=\operatorname{Spec} B$ and $X=\operatorname{Spec} A$. Then the morphism $Y \rightarrow X$ corresponds to an algebra homomorphism $A \xrightarrow{\varphi} B$. If the kernel of $\varphi$ were nonzero, the image of $Y$ would be contained in a proper closed subset of $X$ to which induction would apply. So we may assume that $\varphi$ is injective.

Proposition 4.2.10 tells us that, for suitable nonzero $s$ in $A, B_{s}$ will be a finite module over a polynomial subring $A_{s}\left[y_{1}, \ldots, y_{k}\right]$. Then the maps $Y_{s} \rightarrow \operatorname{Spec} A_{s}[y]$ and Spec $A_{s}[y] \rightarrow X_{s}$ are both surjective, so $Y_{s}$ maps surjectively to $X_{s}$. When we replace $X$ and $Y$ by $X_{s}$ and $Y_{s}$, the map $Y \rightarrow X$ becomes surjective, and we are done.

### 5.5 Closed Sets

(5.5.1) $C$ is a smooth affine curve, $q$ is a point of $C$, and $C^{\prime}$ is the complement of $q$ in $C$.

The closure of $C^{\prime}$ will be $C$, and we think of $q$ as a limit point. Theorem5.5.3. which is below, asserts that a constructible subset of a variety is closed if it contains all such limit points.

The next theorem tells us that there are enough curves to do the job.
5.5.2. Theorem. (enough curves) Let $Y$ be a constructible subset of a variety $X$, and let $p$ be a point of its closure $\bar{Y}$. There exists a morphism $C \xrightarrow{f} X$ from a smooth affine curve to $X$, and a point $q$ of $C$ with $f(q)=p$, such that the image of $C^{\prime}=C-\{q\}$ is contained in $Y$.
proof. We use Krull's Theorem to slice $Y$ down to dimension 1. If $X=p$, then $Y=p$ too. In this case, we may take for $f$ the constant morphism from any curve $C$ to $p$. So we may assume that $X$ has dimension at least one. Next, we may replace $X$ by an affine open subset $X^{\prime}$ that contains $p$, and $Y$ by $Y^{\prime}=Y \cap X^{\prime}$. The closure $\bar{Y}^{\prime}$ of $Y^{\prime}$ in $X^{\prime}$ will be the intersection $\bar{Y} \cap X^{\prime}$, and it will contain $p$. So we may assume that $X$ is affine, say $X=\operatorname{Spec} A$.

Since $Y$ is constructible, it is a union $L_{1} \cup \cdots \cup L_{k}$ of locally closed sets, say $L_{i}=Z_{i} \cap U_{i}$ where $Z_{i}$ are irreducible closed sets and $U_{i}$ are open sets. (We use $Z_{i}$ in place of $C_{i}$ here to avoid confusion with a curve.) The closure of $Y$ is the union $Z_{1} \cup \cdots \cup Z_{k}$, and $p$ will be in at least one of those closed sets, say $p \in Z_{i}$. We replace $X$ by $Z_{i}$ and $Y$ by $L_{i}$. This reduces us to the case that $Y$ is a nonempty open subset of $X$.

Suppose that the dimension $n$ of $X$ is at least two. Let $D=X-Y$ be the (closed) complement of the open set $Y$. The components of $D$ have dimension at most $n-1$. We choose an element $\alpha$ of the coordinate algebra $A$ of $X$ that is zero at $p$ and isn't identically zero on any component of $D$ except $p$ itself, if $p$ happens to be a component. Krull's Theorem tells us that every component of the zero locus of $\alpha$ has dimension $n-1$, and at least one of those components, call it $V$, contains $p$. If $V$ were contained in $D$, it would be a component of $D$ because $\operatorname{dim} V=n-1$ and $\operatorname{dim} D \leq n-1$. By our choice of $\alpha$, this isn't the case. So $V \not \subset D$, and therefore $V \cap Y \neq \emptyset$. Because $V$ is irreducible and $Y$ is open, $W=V \cap Y$ is a dense open subset of $V$, and $p$ is a point of its closure $V$. We replace $X$ by $V$ and $Y$ by $W$. The dimension of $X$ is thereby reduced to $n-1$.

Thus it suffices to treat the case that $X$ has dimension one. In this case, $X$ will be a curve that contains $p$ and $Y$ will be a nonempty open subset of $X$. The normalization of $X$ will be a smooth curve $x_{1}$ that comes with an integral and therefore surjective morphism to $Y$. Finitely many points of $X_{1}$ will map to $p$. We choose for $C$ an affine open subvariety of $X_{1}$ that contains just one of those points, and we call that point $q$.
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Limits of sequences are often used to analyze subsets of a topological space. In the classical topology, a subset $Y$ of $\mathbb{C}^{n}$ is closed if, whenever a sequence of points in $Y$ has a limit in $\mathbb{C}^{n}$, the limit is in $Y$. In algebraic geometry curves can be used as substitutes.

We use the following notation: sise $Y$ is
5.5.3. Theorem (curve criterion for a closed set) Let $Y$ be a constructible subset of a variety $X$. The following conditions are equivalent:
(a) $Y$ is closed.
(b) For every morphism $C \xrightarrow{f} X$ from a smooth affine curve to $X$, the inverse image $f^{-1} Y$ is closed in $C$.
(c) Let $q$ be a point of a smooth affine curve $C$, let $C^{\prime}=C-\{q\}$, and let $C \xrightarrow{f} X$ be a morphism. If $f\left(C^{\prime}\right) \subset Y$, then $f(C) \subset Y$.

The hypothesis that $Y$ be constructible is necessary. For example, in the affine line $X$, the set $Z$ of points with integer coordinates isn't constructible, but it satisfies the curve criterion. Any morphism $C^{\prime} \rightarrow X$ whose image is in $Z$ will map $C^{\prime}$ to a single point, and therefore it will extend to $C$.
proof. The implications $(\mathbf{a}) \Rightarrow(\mathbf{b}) \Rightarrow(\mathbf{c})$ are obvious. We prove the contrapositive of the implication $(\mathbf{c}) \Rightarrow$ (a). Suppose that $Y$ isn't closed. We choose a point $p$ of the closure $\bar{Y}$ that isn't in $Y$, and we apply Theorem 5.5.2. There exists a morphism $C \xrightarrow{f} X$ from a smooth curve to $X$ and a point $q$ of $C$ such that $f(q)=p$ and $f\left(C^{\prime}\right) \subset Y$. Since $q \notin Y$, this morphism shows that (c) doesn't hold either.
5.5.4. Theorem. A constructible subset $Y$ of a variety $X$ is closed in the Zariski topology if and only if it is closed in the classical topology.
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proof. A Zariski closed set is closed in the classical topology because the classical topology is finer than the Zariski topology.

Suppose that $Y$ is closed in the classical topology. Let $q$ be a point of the Zariski closure $\bar{Y}$ of $Y$, and let $C \xrightarrow{f} X$ be a morphism from a smooth affine curve to $X$ that maps the complement $C^{\prime}$ of $q$ to $Y$. Let $Y_{1}=f^{-1} Y$. Then $Y_{1}$ contains $C^{\prime}$, so $Y_{1}$ is either $C^{\prime}$ or $C$. A morphism is a continuous map in the classical topology. Since $Y$ is closed in the classical topology, $Y_{1}$ is closed in $C$. If $Y_{1}$ were equal to $C^{\prime}$, then $\{q\}$ would be open as well as closed. It would be an isolated point of $C$. Since a curve contains no isolated point, the closure is $C$. Therefore the curve criterion (5.5.3) is satisfied, and $Y$ is closed in the Zariski topology.

### 5.6 Fibred Products

## (5.6.1) the mapping property of a product

The product $X \times Y$ of two sets $X$ and $Y$ has a mapping property that is easy to verify: Maps from a set $T$ to the product set $X \times Y$, correspond bijectively to pairs of maps $T \xrightarrow{f} X$ and $T \xrightarrow{g} Y$. The map $T \xrightarrow{(f, g)} X \times Y$ defined by the pair of maps $f, g$ sends a point $t$ to the point pair $(f(t), g(t))$.

Let $X \times Y \xrightarrow{\pi_{1}} X$ and $X \times Y \xrightarrow{\pi_{2}} Y$ denote the projection maps. If $T \xrightarrow{h} X \times Y$ is a map to the product, the corresponding maps to $X$ and $Y$ are the compositions with the projections: $T \xrightarrow{\pi_{1} h} X$ and $T \xrightarrow{\pi_{2} h} Y$ :

The analogous statements are true for morphisms of varieties.
5.6.2. Proposition. Let $X$ and $Y$ be varieties, and let $X \times Y$ be the product variety.
(i) The projections $X \times Y \xrightarrow{\pi_{1}} X$ and $X \times Y \xrightarrow{\pi_{2}} Y$ are morphisms.
(ii) Morphisms from a variety $T$ to the product variety $X \times Y$ correspond bijectively to pairs of morphisms $T \rightarrow X$ and $T \rightarrow Y$, the correspondence being the same as for maps of sets.
(iii) If $X \xrightarrow{f} Z$ and $Y \xrightarrow{g} W$ are morphisms of varieties, the product map $X \times Y \xrightarrow{f \times g} Z \times W$ defined by $[f \times g](x, y)=(f(x), g(y))$ is a morphism.
(5.6.3) fibred products of sets

If $X \xrightarrow{f} Z$ and $Y \xrightarrow{g} Z$ are maps of sets, the fibred product $X \times{ }_{Z} Y$ is the subset of the product $X \times Y$ consisting of pairs of points $x, y$ such that $f(x)=g(y)$. It fits into a diagram
in which $\pi_{1}$ and $\pi_{2}$ are the projections. Many important subsets of a product can be described as fibred products. If a map $Y \rightarrow Z$ is given, and if $p \rightarrow Z$ is the inclusion of a point into $Z$, then $p \times{ }_{Z} Y$ is the fibre of $Y$ over $p$. The diagonal in $X \times X$ is the fibred product $X \times_{X} X$.

The reason for the term "fibred product" is that the fibre of $X \times_{Z} Y$ over a point $x$ of $X$ maps bijectively to the fibre of $Y$ over the image $z=f(x)$, and that the fibre of $X \times{ }_{Z} Y$ over a point $y$ of $Y$ maps bijectively to the fibre of $X$ over the image $g(y)$.

## (5.6.5) fibred products of varieties


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Since we are working with varieties, not schemes, we have a small problem: A fibred product of varieties will be a scheme, but it needn't be a variety.
5.6.6. Example. Let $X=\operatorname{Spec} \mathbb{C}[x], Y=\operatorname{Spec} \mathbb{C}[y]$ and $Z=\operatorname{Spec} \mathbb{C}[z]$ be affine lines, let $X \xrightarrow{f} Z$ and $X \xrightarrow{g} Z$ be the maps defined by $z=x^{2}$ and $z=y^{2}$, respectively. The fibred product $X \times{ }_{Z} Y$ is the closed subset of the affine $x, y$-plane consisting of points $(x, y)$ such that $x^{2}=y^{2}$. It is the union of the two lines $x=y$ and $x=-y$.

The next proposition will be enough for our purposes.
5.6.7. Proposition. Let $X \xrightarrow{f} Z$ and $Y \xrightarrow{g} Z$ be morphisms of varieties. The fibred product $X \times{ }_{Z} Y$ is a closed subset of the product variety $X \times Y$.
5.6.8. Lemma. Let $u$ and $v$ be morphisms of varieties $X \rightarrow Z$. The subset $W$ of points $x$ of $X$ such that $u(x)=v(x)$ is closed in $X$.
proof. In $X \times Z$, let $W^{\prime}$ be the intersection of the graphs of $u$ and $v: W^{\prime}=\Gamma_{u} \cap \Gamma_{v}$. A point $(x, z)$ is in $W^{\prime}$ if $u x=z$ and $v x=z$. This is an intersection of closed sets, so it is closed in $\Gamma_{u}$ (and in $\Gamma_{v}$ ). The projection $\Gamma_{u} \rightarrow X$, which is an isomorphism, carries $W^{\prime}$ to $W$, so $W$ is closed in $X$.
poof of Proposition 5.6.7. The graph $\Gamma_{f}$ of a morphism $X \xrightarrow{f} Z$ of varieties is a closed subvariety of $X \times Z$ isomorphic to $X$ (Proposition 3.4.28). With reference to Diagram5.6.4, $X \times{ }_{Z} Y$ is the subset of the product $X \times Y$ of points at which the maps $f \pi_{X}$ and $g \pi_{Y}$ to $Z$ are equal. The lemma shows that it is closed in $X \times Y$.

### 5.7 Projective Varieties are Proper

As has been noted before, an important property of projective space is that, in the classical topology, it is a compact space. A variety isn't compact in the Zariski topology unless it is a single point. However, in the Zariski topology, projective varieties have a property closely related to compactness: They are proper.

Before defining the concept of a proper variety, we explain an analogous property of compact spaces.
5.7.1. Proposition. Let $X$ be a compact space, let $Z$ be a Hausdorff space, and let $V$ be a closed subset of $Z \times X$. The image of $V$ via projection to $Z$ is a closed subset of $Z$.
proof. Let $W$ be the image of $V$ in $Z$. We show that if a sequence of points $z_{i}$ of $W$ has a limit $\underline{z}$ in $Z$, then $\underline{z}$ is in $W$. For each $i$, we choose a point $p_{i}$ of $V$ that lies over $z_{i}$. So $p_{i}$ is a pair $\left(z_{i}, x_{i}\right), x_{i}$ being a point of $X$. Since $X$ is compact, there is a subsequence of the sequence $x_{i}$ that has a limit $\underline{x}$ in $X$. Passing to subsequences, we may suppose that $x_{i}$ has limit $\underline{x}$. Then $p_{i}$ will have the limit $\underline{p}=(\underline{z}, \underline{x})$. Since $V$ is closed, $\underline{p}$ is in $V$, and therefore $\underline{z}$ is in its image $W$.
5.7.2. Definition. A variety $X$ is proper if has the following property: Let $Z \times X$ be the product with another variety $Z$, let $\pi$ denote the projection $Z \times X \longrightarrow Z$, and let $V$ be a closed subvariety of $Z \times X$. The image $W=\pi_{Z}(V)$ is a closed subvariety of $Z$ :


If $X$ is proper, then because every closed set is a finite union of closed subvarieties, the image of any closed subset of $Z \times X$ will be ca closed subset of $Z$,
5.7.4. Proposition. Let $X$ be a proper variety, let $V$ be a closed subvariety of $X$, and let $X \xrightarrow{f} Y$ be a morphism. The image $f(V)$ of $V$ is a closed subvariety of $Y$.
proof. In $X \times Y$, the graph $\Gamma_{f}$ of $f$ is a closed subset isomorphic to $X$, and via this isomorphism, $V$ corresponds to a subset $V^{\prime}$ of $\Gamma_{f}$ that is closed in $\Gamma_{f}$ and in $X \times Y$. The points of $V^{\prime}$ sre pairs $(x, y)$ such that $x \in V$ and $y=f(x)$. Since $X$ is proper, the image of $V^{\prime}$ via projection to $Y$, which is $f(V)$, is closed.

### 5.7.5. Theorem. Projective varieties are proper.

This is the most important application of the use of curves to characterize closed sets.
proof. Let $X$ be a projective variety. With notation as in Definition 5.7.2, suppose we are given a closed subvariety $V$ of the product $Z \times X$. We must show that its image $W$ is a closed subvariety of $Z$. If the image is a closed set, it will be irreducible. So it suffices to show that $W$ is closed, and to do this, it suffices to show that $W$ is closed in the classical topology (Theorem5.5.4). Theorem 5.4.6tells us that $W$ is a constructible set, and since $X$ is closed in projective space, it is compact in the classical topology. Proposition 5.7.1 tells us that $W$ is closed in the classical topology.
5.7.6. Note. Of course, this is an algebraic theorem, and one would prefer an algebraic proof. To make an algebraic proof, one could attempt to use the curve criterion, proceeding as follows: Given a closed subset $V$ of $Z \times X$ with image $W$ and a point $z$ in the closure of $W$, one chooses a map $C \xrightarrow{f} Z$ from an affine curve $C$ to $Z$ such that $f(q)=p$ and $f\left(C^{\prime}\right) \subset W$. Then one tries to lift this map, defining a morphism $C \xrightarrow{g} Z \times X$ such that $g\left(C^{\prime}\right) \subset V$ and $f=\pi \circ g$. Since $V$ is closed, it would contain $g(q)$, and therefore $f(q)=\pi g(q)$ would be in $\pi(V)=W$. Unfortunately, to find $g$, it may be necessary to replace $C$ by a curve that covers $C$. It isn't very difficult to make this method work, but it takes longer. That is why we resorted to the classical topology.

The next examples show how Theorem5.7.5 can be used.
5.7.7. Example. (singular curves) We parametrize the plane curves of a given degree $d$. The number of monomials $x_{0}^{i} x_{1}^{j} x_{2}^{k}$ of degree $d=i+j+k$ is the binomial coefficient $\binom{d+2}{2}$. We order those monomials arbitrarily, and label them as $m_{0}, \ldots, m_{r}$, with $r=\binom{d+2}{2}-1$. A homogeneous polynomial of degree $d$ will be a combination $\sum z_{i} m_{i}$ of monomials with complex coefficients $z_{i}$, so the homogeneous polynomials $f$ of degree $d$ in $x$, taken up to scalar factors, are parametrized by the projective space of dimension $r$ with coordinates $z$. Let's denote that projective space by $Z$. Points of $Z$ correspond bijectively to divisors of degree $d$ in the projective plane 1.3.12.

The product variety $Z \times \mathbb{P}^{2}$ represents pairs $(D, p)$, where $D$ is a divisor of degree $d$ and $p$ is a point of $\mathbb{P}^{2}$. A variable homogeneous polynomial of degree $d$ in $x$ will be a bihomogeneous polynomial $f(z, x)$ of degree 1 in $z$ and degree $d$ in $x$. So the locus $\Gamma:\{f(z, x)=0\}$ in $Z \times \mathbb{P}^{2}$ is a closed set. Its points are pairs $(D, p)$ such that $D$ is the divisor of $f$ and $p$ is a point of $D$.

Let $\Sigma$ be the set of pairs $(D, p)$ such that $p$ is a singular point of $D$. This is also a closed set. It is defined by the system of equations $f_{0}(z, x)=f_{1}(z, x)=f_{2}(z, x)=0$, where $f_{i}$ are the partial derivatives $\frac{\partial f}{\partial x_{i}}$. (Euler's Formula shows that then $f(x, z)=0$.) The partial derivatives $f_{i}$ are bihomogeneous, of degree 1 in $z$ and degree $d-1$ in $x$.

The next proposition isn't especially easy to prove directly, but the proof becomes easy when one uses the fact that projective space is proper.
5.7.8. Proposition The singular divisors of degree d, the divisors containing at least one singular point, form a closed subset $S$ of the projective space $Z$ of all divisors of degree $d$.
proof. The points of $S$ are the images of points of the set $\Sigma$ via projection to $Z$. Theorem 5.7.5tells us that the image of $\Sigma$ is closed.
5.7.9. Example. (surfaces that contain a line) We go back to the discussion of lines in a surface, as in 3.6. Let $\mathbb{S}$ denote the projective space that parametrizes surfaces of degree $d$ in $\mathbb{P}^{3}$, as before.
5.7.10. Proposition In $\mathbb{P}^{3}$, the surfaces of degree d that contain a line form a closed subset of the space $\mathbb{S}$.
proof. Let $\mathbb{G}$ be the Grassmanian $G(2,4)$ of lines in $\mathbb{P}^{3}$, and let $\Xi$ be the subset of $\mathbb{G} \times \mathbb{S}$ of pairs of pairs $[\ell],[S]$ such that $\ell \subset S$. Lemma 3.6.17 tells us that $\Xi$ is a closed subset of $\mathbb{G} \times \mathbb{S}$. Therefore its image in $\mathbb{S}$ is closed.

### 5.8 Fibre Dimension

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A function $Y \xrightarrow{\delta} \mathbb{Z}$ from a variety to the integers is constructible if, for every integer $n$, the set of points of $Y$ such that $\delta(p)=n$ is constructible, and $\delta$ is upper semicontinuous if for every $n$, the set of points such that $\delta(p) \geq n$ is closed. For brevity, we refer to an upper semicontinuous function as semicontinuous, though the term is ambiguous. A function might be lower semicontinuous.

A function $\delta$ on a curve $C$ is semicontinuous if and only if for every integer $n$, there is a nonempty open subset $C^{\prime}$ of $C$ such that $\delta(p)=n$ for all points $p$ of $C^{\prime}$ and $\delta(p) \geq n$ for all points not in $C^{\prime}$.

The next curve criterion for semicontinuous functions follows from the criterion for closed sets.
5.8.1. Proposition. (curve criterion for semicontinuity) Let $Y$ be a variety. A function $Y \xrightarrow{\delta} \mathbb{Z}$ is semicontinuous if and only if it is a constructible function, and for every morphism $C \xrightarrow{f} Y$ from a smooth curve $C$ to $Y$, the composition $\delta \circ f$ is a semicontinuous function on $C$.

Let $Y \xrightarrow{f} X$ be a morphism of varieties, let $q$ be a point of $Y$, and let $Y_{p}$ be the fibre of $f$ over $p=f(q)$. The fibre dimension $\delta(q)$ of $f$ at $q$ is the maximum among the dimensions of the components of the fibre that contain $q$.

Note. One could also define the fibre dimension of a point $p$ of $X$ to be the dimension of the fibre over $p$. This would be simpler if all components of a fibre have the same dimension. However, it is possible that a fibre contains components whose dimensions are distinct, and if so, then the fibre dimension defined here is more precise.
5.8.2. Theorem. (semicontinuity of fibre dimension) Let $Y \xrightarrow{u} X$ be a morphism of varieties, and let $\delta(q)$ denote the fibre dimension at a point $q$ of $Y$.
(i) Suppose that $X$ is a smooth curve, that $Y$ has dimension $n$, and that $u$ does not map $Y$ to a single point. Then $\delta$ is constant: Every nonempty fibre has constant dimension $n-1$.
(ii) Suppose that the image of $Y$ contains a nonempty open subset of $X$, and let the dimensions of $X$ and $Y$ be $m$ and $n$, respectively. There is a nonempty open subset $X^{\prime}$ of $X$ such that $\delta(q)=n-m$ for every point $q$ in the inverse image of $X^{\prime}$.
(iii) $\delta$ is a semicontinuous function on $Y$.

The proof of this theorem is left as a long exercise. When you have done it, you will have understood the chapter.

## Chapter 6 MODULES

july 10
6.1 The Structure Sheaf
$6.2 \mathcal{O}$-Modules
6.3 The Sheaf Property
6.4 Some $\mathcal{O}$-Modules
6.5 Direct Image
6.6 Support
6.7 Twisting
6.8 Proof of Theorem 6.3.2

This chapter explains how modules on a variety are defined. For this, we review a few facts about localization. Recall that, if $s$ is a nonzero element of a domain $A$, the symbol $A_{s}$ stands for the localization $A\left[s^{-1}\right]$, and if $X=\operatorname{Spec} A$, then $X_{s}=\operatorname{Spec} A_{s}$.

- Let $X=\operatorname{Spec} A$ be an affine variety. The intersection of two localizations $X_{s}=\operatorname{Spec} A_{s}$ and $X_{t}=$ Spec $A_{t}$ is the localization $X_{s t}=\operatorname{Spec} A_{s t}$.
- Let $W \subset V \subset U$ be affine open subsets of a variety $X$. If $V$ is a localization of $U$ and $W$ is a localization of $V$, then $W$ is a localization of $U$ (2.6.2).
- The affine open subsets of a variety $X$ form a basis for the topology on a variety $X$. The localizations of an affine variety form a basis for its topology (??).
- If $U$ and $V$ are affine open subsets of $X$, the open sets $W$ that are localizations of $U$ and localizations of $V$, form a basis for the topology on $U \cap V$. 2.6.2.


### 6.1 The Structure Sheaf.

We introduce two categories associated to a variety $X$. The first is the category (opens). Its objects are the open subsets of $X$, and its morphisms are inclusions: If $U$ and $V$ are open sets and if $V \subset U$, there is a unique morphism $V \rightarrow U$ in (opens), and if $V \not \subset U$ there is no morphism $V \rightarrow U$.

We also introduce a subcategory (affines) of the category (opens). Its objects are the affine open subsets of $X$, and its morphisms are localizations. A morphism $V \rightarrow U$ in (opens) is a morphism in (affines) if $U$ is affine and $V$ is a localization of $U$, i.e., if $V$ is an open subset of the form $U_{s}$, where $s$ is a nonzero element of the coordinate algebra of $U$.

The structure sheaf $\mathcal{O}_{X}$ on a variety $X$ is the functor

$$
\begin{equation*}
\text { (affines }^{\circ} \xrightarrow{\mathcal{O}_{X}} \text { (algebras) } \tag{6.1.1}
\end{equation*}
$$

from affine open sets to algebras, that sends an affine open set $U=\operatorname{Spec} A$ to its coordinate algebra. The coordinate algebra of $U$ is then denoted by $\mathcal{O}_{X}(U)$.

As has been noted before, inclusions $V \rightarrow U$ of affine open subsets needn't be localizations. We focus attention on localizations because the relation between the coordinate algebras of an affine variety and a localization is easy to understand. However, the structure sheaf extends with little difficulty to the category (opens), (See Corollary 6.1.2 below.)

A brief review about regular functions: The function field of a variety $X$ is the field of fractions of the coordinate algebra of any one of its affine open subsets, and a rational function on $X$ is an element of the function field. A rational function $f$ is regular on an affine open set $U=\operatorname{Spec} A$ if it is an element $A$, and $f$ is regular on a nonempty open set $U$ if $U$ can be covered by affine open sets on which it is regular. Thus the function field of a variety $X$ contains the regular functions on every nonempty open subset, and the regular functions are governed by the regular functions on affine open subsets.

An affine variety is determined by its regular functions, but the regular functions don't suffice to determine a variety that isn't affine. For instance, the only rational functions that are regular everywhere on the projective line $\mathbb{P}^{1}$ are the constant functions, which are useless. We will be interested in regular functions on non-affine open sets, especially in functions that are regular on the whole variety, but one should always work with the affine open sets, where the definition of a regular function is clear.

Let $U$ and $V$ be open subsets of a variety $X$, with $V \subset U$. If a rational function is regular on $U$, it is also regular on $V$. Thus if $V \subset U$ is an inclusion of affine open subsets, say $U=\operatorname{Spec} A$ and $V=\operatorname{Spec} B$, then $A \subset B$. However, it won't be clear how to construct $B$ from $A$ unless $B$ is a localization. If $V=U_{s}$, then $B=A\left[s^{-1}\right]$. When $B$ isn't a localization of $A$, the exact relationship between $A$ and $B$ remains obscure.

We extend the notation introduced for affine open sets to all open sets, denoting the algebra of regular functions on any open set $U$ by $\mathcal{O}_{X}(U)$.
6.1.2. Corollary. Let $X$ be a variety. By defining $\mathcal{O}_{X}(U)$ to be the algebra of regular functions on an open subset $U$, the structure sheaf $\mathcal{O}_{X}$ on $X$ extends to a functor

$$
\text { (opens) })^{\circ} \xrightarrow{\mathcal{O}_{X}} \text { (algebras) }
$$

If it is clear which variety is being studied, we may write $\mathcal{O}$ for $\mathcal{O}_{X}$.
The regular functions on $U$, the elements of $\mathcal{O}(U)$, are called sections of the structure sheaf $\mathcal{O}_{X}$ on $U$.
When $V \rightarrow U$ is a morphism in (opens), $\mathcal{O}_{X}(U)$ is contained in $\mathcal{O}_{X}(V)$. This gives us the homomorphism, an inclusion,

$$
\mathcal{O}_{X}(U) \rightarrow \mathcal{O}_{X}(V)
$$

that makes $\mathcal{O}_{X}$ into a functor. Note that arrows are reversed by $\mathcal{O}_{X}$. If $V \rightarrow U$, then $\mathcal{O}_{X}(U) \rightarrow \mathcal{O}_{X}(V)$. A functor that reverses arrows is a contravariant functor. The superscript $\circ$ in 6.1.1) and 6.1.2 is a customary notation to indicate that a functor is contravariant.
6.1.3. Proposition The (extended) structure sheaf has the following sheaf property:

- Let $Y$ be an open subset of $X$, and let $U^{i}=\operatorname{Spec} A_{i}$ be affine open subsets that cover $Y$. Then

$$
\mathcal{O}_{X}(Y)=\bigcap \mathcal{O}_{X}\left(U^{i}\right) \quad\left(=\bigcap A_{i}\right)
$$

This sheaf property is especially simple because regular functions are elements of the function field. It is more complicated for $\mathcal{O}$-modules, which will be defined in the next section.

By definition, if $f$ is a regular function on $X$, there is a covering by affine open sets $U^{i}$ such that $f$ is regular on each of them, i.e., that $f$ is in $\bigcap \mathcal{O}\left(U^{i}\right)$. Therefore the next lemma proves the proposition.
6.1.4. Lemma. Let $Y$ be an open subset of a variety $X$. The intersection $\bigcap \mathcal{O}_{X}\left(U^{i}\right)$ is the same for every affine open covering $\left\{U^{i}\right\}$ of $Y$.

We prove the lemma first in the case of a covering of an affine open set by localizations.
6.1.5. Sublemma. Let $U=\operatorname{Spec} A$ be an affine variety, and let $\left\{U^{i}\right\}$ be a covering of $U$ by localizations, say $U^{i}=\operatorname{Spec} A_{s_{i}}$. Then $A=\bigcap A_{s_{i}}$, i.e., $\mathcal{O}(U)=\bigcap \mathcal{O}\left(U^{i}\right)$.
extendOone
proof. The fact that $A$ is a subset of $\bigcap A_{s_{i}}$ is clear. We prove the opposite inclusion.
A finite subset of the set $\left\{U^{i}\right\}$ will cover $U$, so we may assume that the index set is finite. Let $\alpha$ be an element of $\bigcap A_{s_{i}}$. So $\alpha=s_{i}^{-r} a_{i}$, or $s_{i}^{r} \alpha=a_{i}$ for some $a_{i}$ in $A$ and some integer $r$, and we can use the same $r$ for every $i$. Because $\left\{U^{i}\right\}$ covers $U$, the elements $s_{i}$ generate the unit ideal in $A$, and so do their powers $s_{i}^{r}$. There are elements $b_{i}$ in $A$ such that $\sum b_{i} s_{i}^{r}=1$. Then $\alpha=\sum b_{i} s_{i}^{r} \alpha=\sum b_{i} a_{i}$ is in $A$.
proof of Lemma 6.1.4. Say that $Y$ is covered by affine open sets $\left\{U^{i}\right\}$ and also by affine open sets $\left\{V^{j}\right\}$. We cover the intersections $U^{i} \cap V^{j}$ by open sets $W^{i j \nu}$ that are localizations of $U^{i}$ and also localizations of $V^{j}$. Fixing $i$ and letting $j$ and $\nu$ vary, the set $\left\{W^{i j \nu}\right\}_{j, \nu}$ will be a covering of $U^{i}$ by localizations, and the sublemma shows that $\mathcal{O}\left(U^{i}\right)=\bigcap_{j, \nu} \mathcal{O}\left(W^{i j \nu}\right)$. Then $\bigcap_{i} \mathcal{O}\left(U^{i}\right)=\bigcap_{i, j, \nu} \mathcal{O}\left(W^{i j \nu}\right)$. Similarly, $\bigcap_{j} \mathcal{O}\left(V^{j}\right)=$ $\bigcap_{i, j, \nu} \mathcal{O}\left(W^{i j \nu}\right)$.

## 6.2 $\mathcal{O}$-Modules

An $\mathcal{O}_{X}$-module on a variety $X$ associates a module to every affine open subset.
6.2.1. Definition. An $\mathcal{O}$-module $\mathcal{M}$ on a variety $X$ is a (contravariant) functor

$$
\text { (affines }^{\circ} \xrightarrow{\mathcal{M}}(\text { modules })
$$

such that $\mathcal{M}(U)$ is an $\mathcal{O}(U)$-module for every affine open set $U$, and such that, if $s$ is a nonzero element of $\mathcal{O}(U)$, the module $\mathcal{M}\left(U_{s}\right)$ is the localization of $\mathcal{M}(U)$ :

$$
\mathcal{M}\left(U_{s}\right)=\mathcal{M}(U)_{s}
$$

A section of an $\mathcal{O}$-module $\mathcal{M}$ on an affine open set $U$ is an element of $\mathcal{M}(U)$. An $\mathcal{O}$-module $\mathcal{M}$ is called a finite $\mathcal{O}$-module. If $\mathcal{M}(U)$ is a finite $\mathcal{O}(U)$-module for every affine open set $U$.

A homomorphism $\mathcal{M} \xrightarrow{\varphi} \mathcal{N}$ of $\mathcal{O}$-modules consists of homomorphisms of $\mathcal{O}(U)$-modules

$$
\mathcal{M}(U) \xrightarrow{\varphi(U)} \mathcal{N}(U)
$$

for each affine open subset $U$ of $X$ such that, when $s$ is a nonzero element of $\mathcal{O}(U)$, the homomorphism $\varphi\left(U_{s}\right)$ is the localization of $\varphi(U)$.

A sequence of homomorphisms

$$
\begin{equation*}
\mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{P} \tag{6.2.2}
\end{equation*}
$$

of $\mathcal{O}$-modules on a variety $X$ is exact if the sequence of sections $\mathcal{M}(U) \rightarrow \mathcal{N}(U) \rightarrow \mathcal{P}(U)$ is exact for every affine open subset $U$ of $X$.

At first glance, this definition of an $\mathcal{O}$-module will seem too complicated for comfort. However, when a module has a natural definition, the data involved in the definition take care of themselves. This will become clear as we go along.

Note. When stating that $\mathcal{M}\left(U_{s}\right)$ is the localization of $\mathcal{M}(U)$, it would be more correct to say that $\mathcal{M}\left(U_{s}\right)$ and $\mathcal{M}(U)_{s}$ are canonically isomorphic. Let's not worry about this.

One example of an $\mathcal{O}$-module is the free module $\mathcal{O}^{k}$. The sections of the free module on an affine open set $U$ are the elements of the free $\mathcal{O}(U)$-module $\mathcal{O}(U)^{k}$. In particular, $\mathcal{O}=\mathcal{O}^{1}$ is an $\mathcal{O}$-module.

The kernel, image, and cokernel of a homomorphism $\mathcal{M} \xrightarrow{\varphi} \mathcal{N}$ are among the operations that can be made on $\mathcal{O}$-modules. The kernel $\mathcal{K}$ of $\varphi$ is the $\mathcal{O}$-module defined by $\mathcal{K}(U)=\operatorname{ker}(\mathcal{M}(U) \xrightarrow{\varphi(U)} \mathcal{N}(U))$ for every affine open set $U$, and the image and cokernel are defined analogously. The reason that we work with localizatons is that many operations, including these, are compatible with localization.

### 6.3 The Sheaf Property

In this section, we extend an $\mathcal{O}$-module $\mathcal{M}$ on a variety $X$ to a functor (opens) ${ }^{\circ} \xrightarrow{\widetilde{\mathcal{M}}}$ (modules) on all open subsets of $X$, such that $\widetilde{\mathcal{M}}(Y)$ is an $\mathcal{O}(Y)$-module for every open subset $Y$, and when $U$ is an affine open set, $\widetilde{\mathcal{M}}(U)=\mathcal{M}(U)$.

The tilde is used for clarity here. When we have finished with the discussion, we will use the same notation for a functor on (affines) and for its extension to (opens).
6.3.1. Terminology. If (opens) ${ }^{\circ} \xrightarrow{\widetilde{\mathcal{M}}}$ (modules) is a functor and $U$ is an open subset, an element of $\widetilde{\mathcal{M}}(U)$ is a section of $\widetilde{\mathcal{M}}$ on $U$. If $V \xrightarrow{j} U$ is an inclusion of open subsets, the associated homomorphism $\widetilde{\mathcal{M}}(U) \rightarrow$ $\widetilde{\mathcal{M}}(V)$ is the restriction from $U$ to $V$. The restriction to $V$ of a section $m$ may be denoted by $j^{\circ} m$. However, the operation of restriction occurs very often. Because of this, we usually abbreviate, using the same symbol $m$ for a section and for its restriction. Also, if an open set $V$ is contained in two open sets $U$ and $U^{\prime}$, and if $m, m^{\prime}$ are sections of $\widetilde{\mathcal{M}}$ on $U$ and $U^{\prime}$, respectively, we may say that $m$ and $m^{\prime}$ are equal on $V$ if their restrictions to $V$ are equal.
6.3.2. Theorem. An $\mathcal{O}$-module $\mathcal{M}$ extends uniquely to a functor

$$
\text { (opens) })^{\circ} \xrightarrow{\widetilde{\mathcal{M}}}(\text { modules })
$$

that has the sheaf property described below. Moreover, for every open set $U, \widetilde{\mathcal{M}}(U)$ is an $\mathcal{O}(U)$-module, and for every inclusion $V \rightarrow U$ of nonempty open sets, the map $\widetilde{\mathcal{M}}(U) \rightarrow \widetilde{\mathcal{M}}(V)$ is compatible with scalar multiplication in this sense:

Let $m$ be a section of $\widetilde{\mathcal{M}}$ on $U$, let $\alpha$ be a regular function on $U$, and let $m^{\prime}$ and $\alpha^{\prime}$ denote the restrictions of $m$ and $\alpha$ to $V$. The restriction of $\alpha m$ is $\alpha^{\prime} m^{\prime}$.

Thought the proof of this theorem isn't expecially difficult, but it is lengthy because there are several things to check. In order not to break up the discussion, we have put the proof into Section 6.8 at the end of the chapter.

## (6.3.3) the sheaf property

The sheaf property is the key requirement that determines the extension of an $\mathcal{O}$-module $\mathcal{M}$ to a functor $\widetilde{\mathcal{M}}$ on (opens).

Let $Y$ be an open subset of $X$, and let $\left\{U^{i}\right\}$ be a covering of $Y$ by affine open sets. The intersections $U^{i j}=U^{i} \cap U^{j}$ are also affine open sets, so $\mathcal{M}\left(U^{i}\right)$ and $\mathcal{M}\left(U^{i j}\right)$ are defined. The sheaf property asserts that an element $m$ of $\widetilde{\mathcal{M}}(Y)$ corresponds to a set of elements $m_{i}$ in $\mathcal{M}\left(U^{i}\right)$ such that the restrictions of $m_{j}$ and $m_{i}$ to $U^{i j}$ are equal.

If the affine open subsets $U^{i}$ are indexed by $i=1, \ldots, n$, the sheaf property asserts that an element of $\widetilde{\mathcal{M}}(Y)$ is determined by a vector $\left(m_{1}, \ldots, m_{n}\right)$ with $m_{i}$ in $\mathcal{M}\left(U^{i}\right)$, such that the restrictions of $m_{i}$ and $m_{j}$ to $U^{i j}$ are equal. This means that $\widetilde{\mathcal{M}}(Y)$ is the kernel of the map

$$
\begin{equation*}
\prod_{i} \mathcal{M}\left(U^{i}\right) \xrightarrow{\beta} \prod_{i, j} \mathcal{M}\left(U^{i j}\right) \tag{6.3.4}
\end{equation*}
$$

sheafker
that sends the vector $\left(m_{1}, \ldots, m_{n}\right)$ to the $n \times n$ matrix $\left(z_{i j}\right)$, where $z_{i j}$ is the difference $m_{j}-m_{i}$ of the restrictions of $m_{j}$ and $m_{i}$ to $U^{i j}$. The analogous description is true when the index set is infinite.

In short, the sheaf property tells us that sections of $\widetilde{\mathcal{M}}$ are determined locally: A section on an open set $Y$ is determined by its restrictions to the open subsets $U^{i}$ of any affine covering of $Y$.

Note. The morphisms $U^{i j} \rightarrow U^{i}$ needn't be localizations, and if not the restriction maps $\mathcal{M}\left(U^{i}\right) \rightarrow \mathcal{M}\left(U^{i j}\right)$ aren't a part of the structure of an $\mathcal{O}$-module. We need a definition of the restriction map for an arbitrary inclusion $V \rightarrow U$ of affine open subsets. This point will be taken care of by the proof of Theorem6.3.2 (See Step 2 in Section 6.8) We don't need to worry about it here.

We drop the tilde now, and denote by $\mathcal{M}$ also the extension of an $\mathcal{O}$-module to all open sets. The sheaf property for $\mathcal{M}$ is the statement that, when $\left\{U^{i}\right\}$ is an affine open covering of an open set $U$, the sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{M}(U) \xrightarrow{\alpha} \prod_{i} \mathcal{M}\left(U^{i}\right) \xrightarrow{\beta} \prod_{i, j} \mathcal{M}\left(U^{i j}\right) \tag{6.3.5}
\end{equation*}
$$

sheafproptwo
is exact, where alpha is the product of the restriction maps, and $\beta$ is the map described in 6.3.4.
The next corollary follows from Theorem 6.3.2
injoncover

$$
\begin{equation*}
Y \longleftarrow \mathbb{U}_{0} \leftleftarrows \mathbb{U}_{1} \tag{6.3.9}
\end{equation*}
$$

When we apply a functor (opens) $\xrightarrow{\mathcal{M}}$ (modules) to this diagram, we obtain a sequence
defbeta
6.3.6. Corollary. Let $\left\{U^{i}\right\}$ be an affine open covering of a variety $X$.
(i) An $\mathcal{O}$-module $\mathcal{M}$ is zero if and only if $\mathcal{M}\left(U^{i}\right)=0$ for every $i$.
(ii) A homomorphism $\mathcal{M} \xrightarrow{\varphi} \mathcal{N}$ of $\mathcal{O}$-modules is injective, surjective, or bijective if and only if the maps $\mathcal{M}\left(U^{i}\right) \xrightarrow{\varphi\left(U^{i}\right)} \mathcal{N}\left(U^{i}\right)$ are injective, surjective, or bijective, respectively, for every $i$.
proof. (i) Let $V$ be an open subset of $X$. We can cover the intersections $V \cap U^{i}$ by affine open sets $V^{i \nu}$ that are localizations of $U^{i}$, and these sets, taken together, cover $V$. If $\mathcal{M}\left(U^{i}\right)=0$, then the localizations $\mathcal{M}\left(V^{i \nu}\right)$ are zero too. The sheaf property shows that the map $\mathcal{M}(V) \rightarrow \prod \mathcal{M}\left(V^{i \nu}\right)$ is injective, and therefore that $\mathcal{M}(V)=0$.
(ii) This follows from (i) because a homomorphism $\varphi$ is injective or surjective if and only if its kernel or its cokernel is zero.

## (6.3.7) families of open sets

It is convenient to have a more compact notation for the sheaf property. For this, one can use symbols to represent families of open sets. Say that $\mathbb{U}$ and $\mathbb{V}$ represent families of open sets $\left\{U^{i}\right\}$ and $\left\{V^{\nu}\right\}$, respectively. A morphism of families $\mathbb{V} \rightarrow \mathbb{U}$ consists of a morphism from each $V^{\nu}$ to one of the subsets $U^{i}$. Such a morphism will be given by a map $\nu \rightsquigarrow i_{\nu}$ of index sets, such that $V^{\nu} \subset U^{i_{\nu}}$.

There may be more than one morphism $\mathbb{V} \rightarrow \mathbb{U}$, because a subset $V^{\nu}$ may be contained in more than one of the subsets $U^{i}$. To define a morphism, one must make a choice among the subsets $U^{i}$ that contain $V^{\nu}$. For example, let $\mathbb{U}=\left\{U^{i}\right\}$ be a family of open sets, and let $V$ be another open set. There is a morphism $V \rightarrow \mathbb{U}$ that sends $V$ to $U^{i}$ whenever $V \subset U^{i}$. In the other direction, there is a unique morphism $\mathbb{U} \rightarrow V$ provided that $U^{i} \subset V$ for all $i$.

We extend a functor (opens) ${ }^{\circ} \xrightarrow{\mathcal{M}}$ (modules) to families $\mathbb{U}=\left\{U^{i}\right\}$ by defining

$$
\begin{equation*}
\mathcal{M}(\mathbb{U})=\prod \mathcal{M}\left(U^{i}\right) \tag{6.3.8}
\end{equation*}
$$

Then a morphism of families $\mathbb{V} \xrightarrow{f} \mathbb{U}$ defines a map $\mathcal{M}(\mathbb{V}) \stackrel{f^{\circ}}{\longleftarrow} \mathcal{M}(\mathbb{U})$ in a way that is fairly obvious, though notation for it is clumsy. Say that $f$ is given by a map $\nu \rightsquigarrow i_{\nu}$ of index sets, with $V^{\nu} \rightarrow U^{i_{\nu}}$. A section of $\mathcal{M}$ on $\mathbb{U}$, an element of $\mathcal{M}(\mathbb{U})$, can be thought of as a vector $\left(u_{i}\right)$ with $u_{i} \in \mathcal{M}\left(U^{i}\right)$, and a section of $\mathcal{M}(\mathbb{V})$ as a vector $\left(v_{\nu}\right)$ with $v_{\nu} \in \mathcal{M}\left(V^{\nu}\right)$. If $v_{\nu}$ denotes the restriction of $u_{i_{\nu}}$ to $V^{\nu}$, the map $f^{\circ}$ sends $\left(u_{i_{\nu}}\right) \rightarrow\left(v_{\nu}\right)$.

We write the sheaf property in terms of families of open sets: Let $\mathbb{U}_{0}=\left\{U^{i}\right\}$ be an affine open covering of an open set $Y$, and let $\mathbb{U}_{1}$ denote the family $\left\{U^{i j}\right\}$ of intersections: $U^{i j}=U^{i} \cap U^{j}$. The intersections are also affine. and there are two sets of inclusions

$$
U^{i j} \subset U^{i} \quad \text { and } \quad U^{i j} \subset U^{j}
$$

They define two morphisms of families $\mathbb{U}_{1} \xrightarrow{d_{0}, d_{1}} \mathbb{U}_{0}$ of affine open sets: $U^{i j} \xrightarrow{d_{0}} U^{j}$ and $U i j \xrightarrow{d_{1}} U^{i}$. We also have a morphism $\mathbb{U}_{0} \rightarrow Y$, and Tte two composed morphisms $\mathbb{U}_{1} \xrightarrow{d_{i}} \mathbb{U}_{0} \rightarrow Y$ are equal. These maps form what we all a covering diagram

$$
\begin{equation*}
0 \rightarrow \mathcal{M}(Y) \xrightarrow{\alpha_{\mathbb{U}}} \mathcal{M}\left(\mathbb{U}_{0}\right) \xrightarrow{\beta_{\mathrm{U}}} \mathcal{M}\left(\mathbb{U}_{1}\right) \tag{6.3.10}
\end{equation*}
$$

where $\alpha_{\mathbb{U}}$ is the restriction map and $\beta_{\mathbb{U}}$ is the difference $\mathcal{M}\left(d_{0}\right)-\mathcal{M}\left(d_{1}\right)$ of the maps induced by the two morphisms $\mathbb{U}_{1} \rightrightarrows \mathbb{U}_{0}$. The sheaf property for the covering $\mathbb{U}_{0}$ of $Y$ 6.3.5) is the assertion that this sequence is exact, which means that $\alpha_{\mathbb{U}}$ is injective, and that its image is the kernel of $\beta_{\mathbb{U}}$.
6.3.11. Note. One can suppose that the open sets $U^{i}$ that make a covering are distinct. However, the intersections won't be distinct, because $U^{i j}=U^{j i}$ and $U^{i i}=U^{i}$. These coincidences lead to redundancy in the statement 6.3.10 of the sheaf property. If the indices are $i=1, \ldots, k$, we only need to look at intersections $U^{i j}$ with $i<j$. The product $\mathcal{M}\left(\mathbb{U}_{1}\right)=\prod_{i, j} \mathcal{M}\left(U^{i j}\right)$ that appears in the sheaf property can be replaced by the product $\prod_{i<j} \mathcal{M}\left(U^{i j}\right)$ with increasing pairs of indices. For instance, suppose that an open set $Y$ is covered by two affine open sets $U$ and $V$. The sheaf property is the exact sequence

$$
0 \rightarrow \mathcal{M}(Y) \xrightarrow{\alpha}(\mathcal{M}(U) \times \mathcal{M}(V)) \xrightarrow{\beta}(\mathcal{M}(U \cap U) \times \mathcal{M}(U \cap V) \times \mathcal{M}(V \cap U) \times \mathcal{M}(V \cap V))
$$

is equivalent with the exact sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{M}(Y) \rightarrow(\mathcal{M}(U) \times \mathcal{M}(V)) \xrightarrow{+,-} \mathcal{M}(U \cap V) \tag{6.3.12}
\end{equation*}
$$

### 6.3.13. Example.

Let $A$ denote the polynomial ring $\mathbb{C}[x, y]$, and let $V$ be the complement of a point $p$ in affine space $X=$ $\operatorname{Spec} A$. We cover $V$ by two localizations: $X_{x}=\operatorname{Spec} A\left[x^{-1}\right]$ and $X_{y}=\operatorname{Spec} A\left[y^{-1}\right]$. A regular function on $V$ will be regular on $X_{x}$ and on $X_{y}$, so it will be in the intersection of their coordinate algebras. The intersection $A\left[x^{-1}\right] \cap A\left[y^{-1}\right]$ is $A$. This tells us that the sections of the structure sheaf $\mathcal{O}_{X}$ on $V$ are the elements of $A$. They are the same as the sections on $X$.

We have been working with nonempty open sets. The next lemma takes care of the empty set.
6.3.14. Lemma. The only section of an $\mathcal{O}$-module $\mathcal{M}$ on the empty set is the zero section: $\mathcal{M}(\emptyset)=\{0\}$. In particular, $\mathcal{O}(\emptyset)$ is the zero ring.
proof. This follows from the sheaf property. The empty set is covered by the empty covering, the covering indexed by the empty set. Therefore $\mathcal{M}(\emptyset)$ is contained in an empty product. We want both the empty product and $\mathcal{M}(\emptyset)$ to be modules, and we have no choice but to set them equal to $\{0\}$.

If you find this reasoning pedantic, you can take $\mathcal{M}(\emptyset)=\{0\}$ as an axiom.

## (6.3.15) Interlude: a useful diagram

We consider a commutative diagram of abelian groups of the form

6.3.16. Lemma. (i) Suppose that the rows of the diagram are exact. If $b$ and $c$ are bijective, so is $a$.
(ii) Suppose that the bottom row of the diagram is exact. If a is bijective and $b$ and $c$ are injective, the top row is exact.

It is customary to leave the proofs of such statements to the reader. But since this sort of reasoning may be new, we'll give the proof of part (ii). Here $a, b, c$ and $f^{\prime}$ are assumed to be injective, so $f$ is injective and $g f=0$. It remains to show that $\operatorname{ker} g=\operatorname{im} f$. Let $x$ be an element of ker $g$, so that $g x=0$, and let $x^{\prime}=b x$. Then $g^{\prime} x^{\prime}=g^{\prime} b x=c g x=0$. Since the bottom row is exact, $x^{\prime}=f^{\prime} y^{\prime}$ for some unique element $y^{\prime}$ of $A^{\prime}$. Since $a$ is bijective, $y^{\prime}=a y$ for some element $y$ of $A$. Then bfy $=f^{\prime} a y=f^{\prime} y^{\prime}=x^{\prime}=b x$. Since $b$ is injective, $f y=x$.

## (6.3.17) the coherence property

In addition to the sheaf property, an $\mathcal{O}$-module on a variety $X$ has a property called coherence.
cohprop
6.3.18. Proposition. (the coherence property) Let $Y$ be an open subset of a variety $X$, let $s$ be a nonzero regular function on $Y$, and let $\mathcal{M}$ be an $\mathcal{O}_{X}$-module. Then $\mathcal{M}\left(Y_{s}\right)$ is the localization $\mathcal{M}(Y)_{s}$ of $\mathcal{M}(Y)$.

Compatibility with localization is a requirement for an $\mathcal{O}$-module when $Y$ is affine. The coherence property is the extension to all open subsets.
proof of Proposition 6.3.18 Let $\mathbb{U}_{0}=\left\{U^{i}\right\}$ be a family of affine open sets that covers an open set $Y$. The intersections $U^{i j}$ will be affine open sets too. We inspect the covering diagram $Y \leftarrow \mathbb{U}_{0} \leftleftarrows \mathbb{U}_{1}$. If $s$ is a nonzero regular function on $Y$, the localization of this diagram forms a covering diagram $Y_{s} \leftarrow \mathbb{U}_{0, s} \leftleftarrows \mathbb{U}_{1, s}$, in which $\mathbb{U}_{0, s}=\left\{U_{s}^{i}\right\}$ is an affine covering of $Y_{s}$. Therefore $\mathcal{M}\left(\mathbb{U}_{0}\right)_{s} \approx \mathcal{M}\left(\mathbb{U}_{0, s}\right)$. The sheaf property for the two covering diagrms gives us exact sequences

$$
0 \rightarrow \mathcal{M}(Y) \rightarrow \mathcal{M}\left(\mathbb{U}_{0}\right) \rightarrow \mathcal{M}\left(\mathbb{U}_{1}\right) \quad \text { and } \quad 0 \rightarrow \mathcal{M}\left(Y_{s}\right) \rightarrow \mathcal{M}\left(\mathbb{U}_{0, s}\right) \rightarrow \mathcal{M}\left(\mathbb{U}_{1, s}\right)
$$

and since $s$ is invertible in the second sequence, the localization of the first sequence maps to the second one:


The bottom row is exact, and since localization is an exact operation, the top row of the diagram is exact too. Since $\mathbb{U}_{0}$ and $\mathbb{U}_{1}$ are families of affine open sets, the vertical arrows $b$ and $c$ are bijections. Therefore $a$ is a bijection. This is the coherence property.

### 6.4 Some $\mathcal{O}$-Modules

### 6.4.1. modules on a point

Let's denote a point, the affine variety Spec $\mathbb{C}$, by $p$. The point has only one nonempty open set: the whole space $p$, and $\mathcal{O}_{p}(p)=\mathbb{C}$. Let $\mathcal{M}$ be an $\mathcal{O}_{p}$-module. The space of global sections $\mathcal{M}(p)$ is an $\mathcal{O}_{p}(p)$-module, a complex vector space. To define $\mathcal{M}$, that vector space can be assigned arbitrarily. One may say that a module on the point is a sscomplex vector space.

### 6.4.2. the residue field module $\kappa_{p}$.

Let $p$ be a point of a variety $X$. A residue field module $\kappa_{p}$ is defined as follows: If $U$ is an affine open subset of $X$ that contains $p$, then $\mathcal{O}(U)$ has a residue field $k(p)$ at $p$, and $\kappa_{p}(U)=k(p)$. If $U$ doesn't contain $p$, then $\kappa_{p}(U)=0$.

### 6.4.3. torsion modules.

An $\mathcal{O}$-module $\mathcal{M}$ is a torsion module if $\mathcal{M}(U)$ is a torsion $\mathcal{O}(U)$-module for every affine open set $U$ (see (2.6.11).

### 6.4.4. ideals.

An ideal $\mathcal{I}$ of the structure sheaf is an $\mathcal{O}$-submodule of $\mathcal{O}$.
Let $p$ be a point of a variety $X$. The maximal ideal at $p$, which we denote by $\mathfrak{m}_{p}$, is an ideal. If an affine open subset $U$ contains $p$, its coordinate algebra $\mathcal{O}(U)$ will have a maximal ideal consisting of the elements that vanish at $p$. That maximal ideal is the module of sections $\mathfrak{m}_{p}(U)$ on $U$. If $U$ doesn't contain $p$, then $\mathfrak{m}_{p}(U)=\mathcal{O}(U)$.
\#\#where should this go??\#\#
When $\mathcal{I}$ is an ideal of $\mathcal{O}$, we denote by $V_{X}(\mathcal{I})$ the closed set of points $p$ such that $\mathcal{I} \subset \mathfrak{m}_{p}$ - such that all elements of $\mathcal{I}$ vanish.

### 6.4.5. examples of homomorphisms

(i) There is a homomorphism of $\mathcal{O}$-modules $\mathcal{O} \rightarrow \kappa_{p}$ whose kernel is the maximal ideal $\mathfrak{m}_{p}$.
(ii) Homomorphisms $\mathcal{O}^{n} \rightarrow \mathcal{O}^{m}$ of free $\mathcal{O}$-modules correspond to $m \times n$-matrices of global sections of $\mathcal{O}$.
(iii) Let $\mathcal{M}$ be an $\mathcal{O}$-module. The $\mathcal{O}$-module homomorphisms $\mathcal{O} \xrightarrow{\varphi} \mathcal{M}$ correspond bijectively to global sections of $\mathcal{M}$. This is analogous to the fact that, when $M$ is a module over a ring $A$, homomorphisms
$A \rightarrow M$ correspond to elements of $M$. To be explicit: If $m$ is a global section of $\mathcal{M}$, the homomorphism $\mathcal{O}(U) \xrightarrow{\varphi} \mathcal{M}(U)$ is multiplication by the restriction of $m$ to $U$. It sends a regular function $f$ on $U$ to $f m$.
(iv) If $f$ is a global section of $\mathcal{O}$, scalar multiplication by $f$ defines a homomorphism $\mathcal{M} \xrightarrow{f} \mathcal{M}$.

### 6.4.6. kernel

kerthm
leftexactsection
lexsect
is exact.
proof. We choose a covering diagram $Y \longleftarrow \mathbb{U}_{0} \leftleftarrows \mathbb{U}_{1}$, and inspect the diagram

where the vertical maps are the maps $\beta_{\mathbb{U}}$ described in 6.3.10. The rows are exact because $\mathbb{U}_{0}$ and $\mathbb{U}_{1}$ are families of affines, and the sheaf property asserts that the kernels of the vertical maps form the sequence 6.4.8. That sequence is exact because taking kernels is a left exact operation.

The section functor isn't right exact. When $\mathcal{M} \rightarrow \mathcal{N}$ is a surjective homomorphism of $\mathcal{O}$-modules and $Y$ is a non-affine open set, the map $\mathcal{M}(Y) \rightarrow \mathcal{N}(Y)$ may fail to be surjective. There is an example below. Cohomology, which will be discussed in the next chapter, is a substitute for right exactness.

### 6.4.9. modules on the projective line

The projective line $\mathbb{P}^{1}$ is covered by the standard open sets $\mathbb{U}^{0}$ and $\mathbb{U}^{1}$, and the intersection $\mathbb{U}^{01}=\mathbb{U}^{0} \cap \mathbb{U}^{1}$ is a localization of $\mathbb{U}^{0}$ and of $\mathbb{U}^{1}$. The coordinate algebras of these affine open sets are $\operatorname{co}\left(\mathbb{U}^{0}\right)=A_{0}=\mathbb{C}[u]$ and $\mathcal{O}\left(\mathbb{U}^{1}\right)=A_{1}=\mathbb{C}[v]$, with $v=u^{-1}$, and $\mathcal{O}\left(\mathbb{U}^{01}\right)=A_{01}=\mathbb{C}\left[u, u^{-1}\right]$. The algebra $A_{01}$ is the Laurent polynomial ring whose elements are (finite) combinations of powers of $u$, negative powers included. The form 6.3.12 of the sheaf property asserts that a global section of $\mathcal{O}$ is determined by polynomials $f(u)$ in $A_{0}$ and $g(v)$ in $A_{1}$ such that $f(u)=g\left(u^{-1}\right)$ in $A_{01}$. The only such polynomials $f$ and $g$ are the constants. The only rational functions that are regular everywhere on $\mathbb{P}^{1}$ are the constants. I think we knew this.

If $\mathcal{M}$ is an $\mathcal{O}$-module, $\mathcal{M}\left(\mathbb{U}^{0}\right)=M_{0}$ and $\mathcal{M}\left(\mathbb{U}^{1}\right)=M_{1}$ will be modules over the algebras $A_{0}$ and $A_{1}$, and the $A_{01}$-module $\mathcal{M}\left(\mathbb{U}^{01}\right)=M_{01}$ can be obtained by localizing $M_{0}$ and also by localizing $M_{1}$ : $M_{0}\left[u^{-1}\right] \approx M_{01} \approx M_{1}\left[v^{-1}\right]$. As 6.3.12 tells us, a global section of $\mathcal{M}$ is determined by a pair of elements $m_{1}, m_{2}$ in $M_{1}, M_{2}$ that become equal in the common localization $M_{01}$. In fact, these data determine the module $\mathcal{M}$.
6.4.10. Lemma. With notation as above, let $M_{0}, M_{1}$, and $M_{01}$ be modules over $A_{0}, A_{1}$, and $A_{01}$, respectively, and let $M_{0}\left[u^{-1}\right] \xrightarrow{\varphi_{0}} M_{01}$ and $M_{1}\left[v^{-1}\right] \xrightarrow{\varphi_{1}} M_{01}$ be $A_{01}$-isomorphisms. There is an $\mathcal{O}_{X}$-module $\mathcal{M}$, unique up to isomorphism, that gives this data: $\mathcal{M}\left(\mathbb{U}^{0}\right)$ and $\mathcal{M}\left(\mathbb{U}^{1}\right)$ are isomorhic to $M_{0}$ and $M_{1}$, and the diagram belpow ommutes.

tensprodmodule
tensprodmap
residue-fieldmodule
idealsheafatp

The proof is at the end of this section.
Suppose that $M_{0}$ and $M_{1}$ are free modules of rank $r$ over $A_{0}$ and $A_{1}$,. Then $M_{01}$ will be a free $A_{01}$-module of rank $r$. A basis $\mathbf{B}_{0}$ of the free $A_{0}$-module $M_{0}$ will also be a basis of the $A_{01}$-module $M_{01}$, and a basis $\mathbf{B}_{1}$ of $M_{1}$ will be a basis of $M_{01}$. When regarded as bases of $M_{01}, \mathbf{B}_{\mathbf{0}}$ and $\mathbf{B}_{\mathbf{1}}$ will be related by an $r \times r$ invertible $A_{01}$-matrix $P$, and that matrix determines $\mathcal{M}$ up to isomorphism. When $r=1, P$ will be an invertible $1 \times 1$ matrix in the Laurent polynomial ring $A_{01}-$ a unit of that ring. The units in $A_{01}$ are scalar multiples of powers of $u$. Since the scalar can be absorbed into one of the bases, an $\mathcal{O}$-module of rank 1 is determined, up to isomorphism, by a power of $u$. It is one of the twisting modules that will be described in Section 6.7 .

The Birkhoff-Grothendieck Theorem, which will be proved in Chapter 8 , describes the $\mathcal{O}$-modules on the projective line whose sections on $\mathbb{U}^{0}$ and on $\mathbb{U}^{1}$ are free. They are direct sums of free $\mathcal{O}$-modules of rank one. This means that by changing the bases $\mathbf{B}_{0}$ and $\mathbf{B}_{1}$, one can diagonalize the matrix $P$. Such changes of basis will be given by an invertible $A_{0}$-matrix $Q_{0}$ and an invertible $A_{1}$-matrix $Q_{1}$, respectively. In down-toEarth terms, the Birkhoff-Grothendieck Theorem asserts that, for any invertible $A_{01}$-matrix $P$, there exist two matrices: an invertible $A_{0}$-matrix $Q_{0}$ and an invertible $A_{1}$-matrix $Q_{1}$, such that $Q_{0}^{-1} P Q_{1}$ is diagonal. This can be proved by matrix operations.

### 6.4.11. tensor products

Tensor products are compatible with localization. If $M$ and $N$ are modules over a domain $A$ and $s$ is a nonzero element of $A$, the canonical map $\left(M \otimes_{A} N\right)_{s} \rightarrow M_{s} \otimes_{A_{s}} N_{s}$ is an isomorphism. Therefore the tensor product $\mathcal{M} \otimes_{\mathcal{O}} \mathcal{N}$ of $\mathcal{O}$-modules $\mathcal{M}$ and $\mathcal{N}$ is defined. On an affine open set $U,\left[\mathcal{M} \otimes_{\mathcal{O}} \mathcal{N}\right]$ is the tensor product $\mathcal{M}(U) \otimes_{\mathcal{O}(U)} \mathcal{N}(U)$.

Let $\mathcal{M}$ and $\mathcal{N}$ be $\mathcal{O}$-modules, let $\mathcal{M} \otimes_{\mathcal{O}} \mathcal{N}$ be the tensor product module, and let $V$ be an arbitrary open subset of $X$. There is a canonical map

$$
\begin{equation*}
\mathcal{M}(V) \otimes_{\mathcal{O}(V)} \mathcal{N}(V) \rightarrow\left[\mathcal{M} \otimes_{\mathcal{O}} \mathcal{N}\right](V) \tag{6.4.12}
\end{equation*}
$$

By definition of the tensor product module, this map exists and is an equality when $V$ is affine. For arbitrary $V$, we cover by a family $\mathbb{U}_{0}$ of affine open sets. The family $\mathbb{U}_{1}$ of intersections also consists of affine open sets. We form a diagram


The composition of the two arrows in the top row is zero, the bottom row is exact, and the vertical maps $b$ and $c$ are equalities. The canonical map $a$ is induced by the diagram. It is bijective when $V$ is affine, but when $V$ isn't affine, it needn't be either injective or surjective.
6.4.13. Examples. (i) Let $p$ and $q$ be distinct points of the projective line $X$, and let $\kappa_{p}$ and $\kappa_{q}$ be the residure field modules on $X$. Then $\kappa_{p}(X)=\kappa_{q}(X)=\mathbb{C}$, so $\kappa_{p}(X) \otimes_{\mathcal{O}(X)} \kappa_{q}(X) \approx \mathbb{C} \otimes_{\mathbb{C}} \mathbb{C}=\mathbb{C}$. But $\kappa_{p} \otimes_{\mathcal{O}} \kappa_{q}=0$. The canonical map 6.4.12 is the zero map. It isn't injective.
(ii) Let $p$ a point of a variety $X$, and let $\mathfrak{m}_{p}$ and $\kappa_{p}$ be the maximal ideal and residue field modules at $p$. There is an exact sequence of $\mathcal{O}$-modules

$$
\begin{equation*}
0 \rightarrow \mathfrak{m}_{p} \rightarrow \mathcal{O} \xrightarrow{\pi_{p}} \kappa_{p} \rightarrow 0 \tag{6.4.14}
\end{equation*}
$$

In this case, the sequence of global sections is exact.
(iii) Let $p_{0}$ and $p_{1}$ be the points $(1,0)$ and $(0,1)$ of the projective line $\mathbb{P}^{1}$. We form a homomorphism

$$
\mathfrak{m}_{p_{0}} \times \mathfrak{m}_{p_{1}} \xrightarrow{\varphi} \mathcal{O}
$$

$\varphi$ being the map $(a, b) \mapsto b-a$. On the open set $\mathbb{U}^{0}, \mathfrak{m}_{p_{1}} \rightarrow \mathcal{O}$ is bijective and therefore surjective. Similarly, $\mathfrak{m}_{p_{0}} \rightarrow \mathcal{O}$ is surjective on $\mathbb{U}^{1}$. Therefore $\varphi$ is surjective. The only global sections of $\mathfrak{m}_{p}, \mathfrak{m}_{q}$, and $\mathfrak{m}_{p_{0}} \times \mathfrak{m}_{p_{1}}$ are the zero sections, while $\mathcal{O}$ has the nonzero global section 1 . The map $\varphi$ isn't surjective on global sections.

### 6.4.15. the function field module

Let $F$ be the function field of a variety $X$. The function field module $\mathcal{F}$ is defined as follows: Its module of sections any nonempty open set $U$ is the field $F$. This is an $\mathcal{O}$-module. It is called a constant $\mathcal{O}$-module because the modules of sections $\mathcal{F}(U)$ are the same for ever $U$. It isn't a finite module unless $X$ is a point.

Tensoring with the function field module: Let $\mathcal{M}$ be an $\mathcal{O}$-module on a variety $X$, and let $\mathcal{F}$ be the function field module. Then $\mathcal{M} \otimes_{\mathcal{O}} \mathcal{F}$ is a constant $\mathcal{O}$-module whose sections on any affine open set $U$ form an $F$-vector space. But if $\mathcal{M}$ is a torsion module, $\mathcal{M} \otimes_{\mathcal{O}} \mathcal{F}$ will be zero.
(6.4.16) limits of $\mathcal{O}$-modules
6.4.17. A directed set $M_{\bullet}$ is a sequence of maps of sets $M_{0} \rightarrow M_{1} \rightarrow M_{2} \rightarrow \ldots$. Its limit $\underset{\longrightarrow}{\lim } M_{\bullet}$ is the set of equivalence classes on the union $\bigcup M_{k}$, the equivalence relation being that elements $m$ in $\vec{M}_{i}$ and $m^{\prime}$ in $M_{j}$ are equivalent if they have the same image in $M_{n}$ when $n$ is sufficiently large. An element of $\xrightarrow[\longrightarrow]{\lim } M_{\bullet}$ will be represented by an element of $M_{i}$ for some $i$.
6.4.18. Example. Let $R=\mathbb{C}[x]$ and let $\mathfrak{m}$ be the maximal ideal $x R$. Repeated multiplication by $x$ defines a directed set

$$
R \xrightarrow{x} R \xrightarrow{x} R \xrightarrow{x} R \ldots
$$

whose limit is isomorphic to the Laurent Polynomial Ring $R\left[x^{-1}\right]=\mathbb{C}\left[x, x^{-1}\right]$. Proving this is an exercise.
A directed set of $\mathcal{O}$-modules on a variety $X$ is a sequence $\mathcal{M}_{\bullet}=\left\{\mathcal{M}_{0} \rightarrow \mathcal{M}_{1} \rightarrow \mathcal{M}_{2} \rightarrow \cdots\right\}$ of homomorphisms of $\mathcal{O}$-modules. So, for every affine open set $U$, the $\mathcal{O}(U)$-modules $\mathcal{M}_{n}(U)$ form a directed set, as defined in 6.4.17. The direct limit $\lim _{\longrightarrow} \mathcal{M}_{\bullet}$ is defined simply, by taking the limit for each affine open set: $\left[\underset{\longrightarrow}{\lim } \mathcal{M}_{\bullet}\right](U)=\underset{\longrightarrow}{\lim }\left[\mathcal{M}_{\bullet}(U)\right]$. This limit operation is compatible with localization, so $\underset{\longrightarrow}{\lim } \mathcal{M}_{\bullet}$ is an $\mathcal{O}$-module. In fact, the equality $\left[\lim _{\longrightarrow} \mathcal{M}_{\bullet}\right](U)=\underset{\longrightarrow}{\lim }\left[\mathcal{M}_{\bullet}(U)\right]$ is true for every open set.
6.4.19. Lemma. (i) The limit operation is exact. If $\mathcal{M}_{\bullet} \rightarrow \mathcal{N}_{\bullet} \rightarrow \mathcal{P}_{\bullet}$ is an exact sequence of directed sets of $\mathcal{O}$-modules, the limits form an exact sequence.
(ii) Tensor products are compatible with limits: If $\mathcal{N}_{\bullet}$ is a directed set of $\mathcal{O}$-modules and $\mathcal{M}$ is another $\mathcal{O}$-module, then $\underset{\longrightarrow}{\lim }\left[\mathcal{M} \otimes_{\mathcal{O}} \mathcal{N}_{\bullet}\right] \approx \mathcal{M} \otimes_{\mathcal{O}}\left[\underline{\longrightarrow} \mathcal{N}_{\bullet}\right]$.
6.4.20. Proposition. Let $X=\operatorname{Spec} A$ be an affine variety. Sending an $\mathcal{O}$-module $\mathcal{M}$ to the $A$-module $\mathcal{M}(X)$ of its global sections defines a bijective correspondence between $\mathcal{O}$-modules and $A$-modules.
proof. We must invert the functor $\mathcal{O}$-(modules) $\rightarrow A$-(modules) that sends $\mathcal{M}$ to $\mathcal{M}(X)$. Given an $A$-module $M$, the corresponding $\mathcal{O}$-module $\mathcal{M}$ is defined as follows: Let $U=\operatorname{Spec} B$ be an affine open subset of $X$. The inclusion $U \subset X$ corresponds to an algebra homomorphism $A \rightarrow B$. We define $\mathcal{M}(U)$ to be the $B$ module $B \otimes_{A} M$. If $s$ is a nonzero element of $B$, then $B_{s} \otimes_{A} M$ is the localization $\left(B \otimes_{A} M\right)_{s}$ of $B \otimes_{A} M$. So $\mathcal{M}$ is an $\mathcal{O}$-module.

### 6.4.21. Example.

This example shows that, when an open subset isn't affine, defining $\mathcal{M}(V)=B \otimes_{A} M$, as in Proposition 6.4.20 may be wrong. Let $X$ be the affine plane $\operatorname{Spec} A, \quad A=\mathbb{C}[x, y]$, let $V$ be the complement of the origin in $X$, and let $M$ be the $A$-module $A / y A$. This module can be identified with $\mathbb{C}[x]$, which becomes an $A$-module when scalar multiplication by $y$ is defined to be zero. Here $\mathcal{O}(V)=\mathcal{O}(X)=A$ 6.3.13). If we followed the method used for affine open sets, we would set $\mathcal{M}(V)=A \otimes_{A} M=\mathbb{C}[x]$.

To identify $\mathcal{M}(V)$ correctly, we cover $V$ by the two affine open sets $V_{x}=\operatorname{Spec} A\left[x^{-1}\right]$ and $V_{y}=$ Spec $A\left[y^{-1}\right]$. Then $\mathcal{M}\left(V_{x}\right)=M\left[x^{-1}\right]$ while $\mathcal{M}\left(V_{y}\right)=0$. The sheaf property of $\mathcal{M}$ shows that $\mathcal{M}(V) \approx$ $\mathcal{M}\left(V_{x}\right)=M\left[x^{-1}\right]=\mathbb{C}\left[x, x^{-1}\right]$.
proof of Lemma 6.4.10 With notation as in the proposition, we suppose given the modules $M_{0}, M_{1}$ and isomorphisms $M_{0}\left[u^{-1}\right] \rightarrow M_{1}\left[v^{-1}\right]$, and we are to show that this data comes from an $\mathcal{O}$-module $\mathcal{M}$. The
previous proposition shows that $M_{i}$ defines modules $\mathcal{M}_{i}$ on $\mathbb{U}^{i}$ whose restrictions to $\mathbb{U}^{01}$ are isomorphic. Let's denote these modules by $\mathcal{M}$. Then $\mathcal{M}$ is defined on open sets that are contained in $\mathbb{U}^{0}$ or in $\mathbb{U}^{1}$.

Let $V$ be an arbitrary open set $V$, and let $V^{i}=V \cap \mathbb{U}^{i}$ we define $\mathcal{M}(V)$ to be the kernel of the map $\mathcal{M}\left(V^{0} \times \mathcal{M}\left(V^{1}\right)\right) \rightarrow \mathcal{M}\left(V^{01}\right)$. It is clear that with this definition, $\mathcal{M}$ is a functor. We must verify the sheaf property, and the notation gets confusing. We suppose given an open covering $\left\{V^{\nu}\right\}$ of $V$, and to avoid confusion with $V^{0}$ and $V^{1}$, we label the corresponding covering diagram as $V @ \lll \mathbb{W}_{0} \leftleftarrows \mathbb{W}_{1}$. So $\mathbb{W}_{0}=\left\{V^{\nu}\right\}$. We form a diagram

in which the first asterisk stands for $\mathcal{M}\left(\mathbb{W}_{0} \cap \mathbb{U}^{0}\right) \times \mathcal{M}\left(\mathbb{W}_{0} \cap \mathbb{U}^{1}\right)$, etc. The columns are exact by our definition of $\mathcal{M}$, and the second and third columns are exact because the open sets involved are contained in $\mathbb{U}^{0}$ or $\mathbb{U}^{1}$. It follows that the top row is exact. This is the sheaf property.

### 6.5 Direct Image

Let $Y \xrightarrow{f} X$ be a morphism of varieties, and let $\mathcal{N}$ be an $\mathcal{O}_{Y^{-}}$-module. The direct image $f_{*} \mathcal{N}$ is an $\mathcal{O}_{X^{-}}$ module that is defined as follows: The sections of $f_{*} \mathcal{N}$ on an affine open subset $U$ of $X$ are the sections of $\mathcal{N}$ on the inverse image $V=f^{-1} U$ in $Y$. For example, the direct image $f_{*} \mathcal{O}_{Y}$ of the structure sheaf $\mathcal{O}_{Y}$ is the functor

$$
\mathcal{O}_{Y}-\text { modules } \xrightarrow{f_{*}} \mathcal{O}_{X} \text {-modules }
$$

defined by $\left[f_{*} \mathcal{O}_{Y}\right](U)=\mathcal{O}_{Y}\left(f^{-1} U\right)$.
The direct image generalizes restriction of scalars in modules over rings. If $A \xrightarrow{\varphi} B$ is an algebra homomorphism and $N$ is a $B$-module, one can restrict scalars to make $N$ into an $A$-module. Scalar multiplication by an element $a$ of $A$ on the restricted module $N$ is defined to be scalar multiplication by its image $\varphi(a)$. For clarity, we will sometimes denote the given $B$-module by $N_{B}$ and the $A$-module obtained by restriction of scalars by $N_{A}$. The additive groups $N_{B}$ and $N_{A}$ are the same, but the scalars change.

Let $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$, and let $Y \xrightarrow{f} X$ be the morphism of affine varieties defined by an algebra homomorphism $A \xrightarrow{\varphi} B$. An $\mathcal{O}_{Y}$-module $\mathcal{N}$ is determined by a $B$-module $N=N_{B}$. The direct image $f_{*} \mathcal{N}$ is the $\mathcal{O}_{X}$-module determined by the $A$-module $N_{A}$.
6.5.1. Lemma. Let $Y \xrightarrow{f} X$ be a morphism of varieties. The direct image $f_{*} \mathcal{N}$ of an $\mathcal{O}_{Y}$-module $\mathcal{N}$ is an $\mathcal{O}_{X}$-module. Moreover, for all open subsets $U$ of $X$, not only for affine open subsets,

$$
\left[f_{*} \mathcal{N}\right](U)=\mathcal{N}\left(f^{-1} U\right)
$$

proof. Let $U^{\prime} \rightarrow U$ be an inclusion of affine open subsets of $X$, and let $V=f^{-1} U$ and $V^{\prime}=f^{-1} U^{\prime}$. These inverse images are open subsets of $Y$, but they aren't necessarily affine. The inclusion $V^{\prime} \rightarrow V$ gives us a homomorphism $\mathcal{N}(V) \rightarrow \mathcal{N}\left(V^{\prime}\right)$, and therefore a homomorphism $f_{*} \mathcal{N}(U) \rightarrow f_{*} \mathcal{N}\left(U^{\prime}\right)$. So $f_{*} \mathcal{N}$ is a functor. Its $\mathcal{O}_{X}$-module structure is explained as follows: Composition with $f$ defines a homomorphism $\mathcal{O}_{X}(U) \rightarrow \mathcal{O}_{Y}(V)$, and $\mathcal{N}(V)$ is an $\mathcal{O}_{Y}(V)$-module. Restriction of scalars to $\mathcal{O}_{X}(U)$ makes $\left[f_{*} \mathcal{N}\right](U)=$ $\mathcal{N}(V)$ into an $\mathcal{O}_{X}(U)$-module.

To show that $f_{*} \mathcal{N}$ is an $\mathcal{O}_{X}$-module, we must show that if $s$ is a nonzero element of $\mathcal{O}_{X}(U)$, then $\left[f_{*} \mathcal{N}\right]\left(U_{s}\right)$ is obtained by localizing $\left[f_{*} \mathcal{N}\right](U)$. Let $s^{\prime}$ be the image of $s$ in $\mathcal{O}_{V}(V)$. Scalar multiplication
by $s$ on $\left[f_{*} \mathcal{N}\right](U)$ is given by restriction of scalars, so it is the same as scalar multiplication by $s^{\prime}$ on $\mathcal{N}(V)$. If $s^{\prime} \neq 0$, the localization $V_{s^{\prime}}$ is the inverse image of $U_{s}$. So $\left[f_{*} \mathcal{N}\right]\left(U_{s}\right)=\mathcal{N}\left(V_{s^{\prime}}\right)$. The coherence property 6.3.17) tells us that $\mathcal{N}\left(V_{s^{\prime}}\right)=\mathcal{N}(V)_{s^{\prime}}$. Then $\left[f_{*} \mathcal{N}\right]\left(U_{s}\right)=\mathcal{N}\left(V_{s^{\prime}}\right)=\mathcal{N}(V)_{s^{\prime}}=\left[\left[f_{*} \mathcal{N}\right](U)\right]_{s}$.

If $s^{\prime}=0$, then $\mathcal{N}(V)_{s^{\prime}}=0$. In this case, because scalar multiplication is defined by restricting scalars, $s$ annihilates $\left[f_{*} \mathcal{N}\right](U)$, and therefore $\left[f_{*} \mathcal{N}\right](U)_{s}=0$ too.
6.5.2. Lemma. Direct images are compatible with limits: If $\mathcal{M}_{\bullet}$ is a directed set of $\mathcal{O}$-modules, then $\lim _{\longrightarrow}\left(f_{*} \mathcal{M}_{\bullet}\right) \approx f_{*}\left(\lim _{\longrightarrow} \mathcal{M}_{\bullet}\right)$.

## extension by zero

When $Y \xrightarrow{i} X$ is the inclusion of a closed subvariety into a variety $X$, and $\mathcal{N}$ is an $\mathcal{O}_{Y}$-module, the direct image $i_{*} \mathcal{N}$ is also called the extension by zero of $\mathcal{N}$. If $U$ is an open subset of $X$ then, because $i$ is an inclusion map, $i^{-1} U=U \cap Y$. Therefore

$$
\left[i_{*} \mathcal{N}\right](U)=\mathcal{N}(U \cap Y)
$$

The term "extension by zero" refers to the fact that, when an open set $U$ of $X$ doesn't meet $Y$, the intersection $U \cap Y$ is empty, and the module of sections of $\left[i_{*} \mathcal{N}\right](U)$ is zero. So $i_{*} \mathcal{N}$ is zero outside of the closed set $Y$.

### 6.5.4. Examples.

(i) Let $p \xrightarrow{i} X$ be the inclusion of a point into a variety. When we view the residue field $k(p)$ as an $\mathcal{O}$-module on $p$, its extension by zero $i_{*} k(p)$ is the residue field module $\kappa_{p}$.
(ii) Let $Y \xrightarrow{i} X$ be the inclusion of a closed subvariety, and let $\mathcal{I}$ be the ideal of $Y$ in $\mathcal{O}_{Y}$. The extension by zero of the structure sheaf on $Y$ fits into an exact sequence of $\mathcal{O}_{X}$-modules

$$
0 \rightarrow \mathcal{I} \rightarrow \mathcal{O}_{X} \rightarrow i_{*} \mathcal{O}_{Y} \rightarrow 0
$$

So the extension by zero $i_{*} \mathcal{O}_{Y}$ is isomorphic to the quotient module $\mathcal{O}_{X} / \mathcal{I}$.
6.5.5. Proposition. Let $Y \xrightarrow{i} X$ be the inclusion of a closed subvariety $Y$ into a variety $X$, and let $\mathcal{I}$ be the ideal of $Y$. Let $\mathbb{M}$ denote the subcategory of the category of $\mathcal{O}_{X}$-modules that are annihilated by $\mathcal{I}$. Extension by zero defines an equivalence of categories

$$
\left(\mathcal{O}_{Y}-\text { modules }\right) \xrightarrow{i_{*}} \mathbb{M}
$$

proof. Let $f$ be a section of $\mathcal{O}_{X}$ on an affine open set $U$, let $\bar{f}$ be its restriction to $U \cap Y$, and let $\alpha$ be an element of $\left[i_{*} \mathcal{N}\right](U)(=\mathcal{N}(U \cap Y))$. So multiplication by $f$ is defined by restriction of scalars: $f \alpha=\bar{f} \alpha$. If $f$ is in $\mathcal{I}(U)$, then $\bar{f}=0$ and therefore $f \alpha=\bar{f} \alpha=0$. So the extension by zero of an $\mathcal{O}_{Y}$-module is annihilated by $\mathcal{I}$. The direct image $i_{*} \mathcal{N}$ is an object of $\mathbb{M}$.

We construct a quasi-inverse to the direct image. Starting with an $\mathcal{O}_{X}$-module $\mathcal{M}$ that is annihilated by $\mathcal{I}$, we construct an $\mathcal{O}_{Y}$-module $\mathcal{N}$ such that $i_{*} \mathcal{N}$ is isomorphic to $\mathcal{M}$.

Let $Y^{\prime}$ be an open subset of $Y$. The topology on $Y$ is induced from the topology on $X$, so $Y^{\prime}=X_{1} \cap Y$ for some open subset $X_{1}$ of $X$. We try to set $\mathcal{N}\left(Y^{\prime}\right)=\mathcal{M}\left(X_{1}\right)$. To show that this is well-defined, we show that if $X_{2}$ is another open subset of $X$ such that $Y^{\prime}=X_{2} \cap Y$, then $\mathcal{M}\left(X_{2}\right)$ is isomorphic to $\mathcal{M}\left(X_{1}\right)$. Let $X_{3}=X_{1} \cap X_{2}$. Then it is also true that $Y^{\prime}=X_{3} \cap Y$. Since $X_{3} \subset X_{1}$, we have a map $\mathcal{M}\left(X_{1}\right) \rightarrow \mathcal{M}\left(X_{3}\right)$, and It suffices to show that this map is an isomorphism. The same reasoning will give us an isomorphism $\mathcal{M}\left(X_{2}\right) \rightarrow \mathcal{M}\left(X_{3}\right)$.

The complement $U=X_{1}-Y^{\prime}$ of $Y^{\prime}$ in $X_{1}$ is an open subset of $X_{1}$ and of $X$, and $U \cap Y=\emptyset$. We cover $U$ by a set $\left\{U^{i}\right\}$ of affine open sets. Then $X_{1}$ is covered by the open sets $\left\{U^{i}\right\}$ together with $X_{3}$. The restriction of $\mathcal{I}$ to each of the sets $U^{i}$ is the unit ideal, and since $\mathcal{I}$ annihilates $\mathcal{M}, \mathcal{M}\left(U^{i}\right)=0$. The sheaf property shows that $\mathcal{M}\left(X_{1}\right)$ is isomorphic to $\mathcal{M}\left(X_{3}\right)$. The rest of the proof is boring.
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toopen
exam-pledirectimage

## fstarexact

dirimstandaffine

Let $Y \xrightarrow{j} X$ be the inclusion of an open subvariety $Y$ into a variety $X$.
First, restriction from $X$ to $Y$. Since open subsets of $Y$ are also open subsets of $X$, we can restrict an $\mathcal{O}$-module $\mathcal{M}$ from $X$ to $Y$. By definition, the sections of the restricted module on a subset $U$ of $Y$ are simply the elements of $\mathcal{M}(U)$. For example, the restriction of the structure sheaf $\mathcal{O}_{X}$ to $Y$ is the structure sheaf $\mathcal{O}_{Y}$. We use subscript notation, writing $\mathcal{M}_{Y}$ for the restriction of an $\mathcal{O}_{X}$-module $\mathcal{M}$ to $Y$, then denoting the given module $\mathcal{M}$ by $\mathcal{M}_{X}$ for clarity. If $U$ is an open subset of $Y$,

$$
\begin{equation*}
\mathcal{M}_{Y}(U)=\mathcal{M}_{X}(U) \tag{6.5.7}
\end{equation*}
$$

Now the direct image: Let $Y \xrightarrow{j} X$ be the inclusion of an open subvariety $Y$, and let $\mathcal{N}$ be an $\mathcal{O}_{Y^{-}}$ module. The inverse image of an open subset $U$ of $X$ is the intersection $Y \cap U$. By definition, the direct image is

$$
\left[j_{*} \mathcal{N}\right](U)=\mathcal{N}(Y \cap U)
$$

For example, $\left[j_{*} \mathcal{O}_{Y}\right](U)$ is the algebra of rational functions on $X$ that are regular on $Y \cap U$. They needn't be regular on $U$.
6.5.8. Example. Let $X_{s} \xrightarrow{j} X$ be the inclusion of a localization into an affine variety $X=\operatorname{Spec} A$. Modules on $X$ correspond to their global sections, which are $A$-modules. Similarly, modules on $X_{s}$ correspond to $A_{s}$-modules. We can restrict the $\mathcal{O}_{X}$-module $\mathcal{M}_{X}$ that corresponds to an $A$-module $M$ to the open set $X_{s}$, obtaining the $\mathcal{O}_{X_{s}}$-module $\mathcal{M}_{X_{s}}$ that corresponds to the $A_{S}$-module $M_{s}$. Then $M_{s}$ is also the module of global sections of $j_{*} \mathcal{M}_{X_{s}}$ on $X$ :

$$
\left[j_{*} \mathcal{M}_{X_{s}}\right](X)=\mathcal{M}_{x_{s}}\left(X_{s}\right)=M_{s}
$$

The localization $M_{s}$ is made into an $A$-module by restriction of scalars.
The reversal of arrows when one passes from affine varieties to their coordinate algebras seems especially confusing here.
6.5.9. Proposition. Let $Y \xrightarrow{j} X$ be the inclusion of an open subvariety $Y$ into a variety $X$.
(i) The restriction $\mathcal{O}_{X}$-modules $\rightarrow \mathcal{O}_{Y}$-modules is an exact operation.
(ii) If $Y$ is an affine open subvariety of $X$, the direct image functor $j_{*}$ is exact.
(iii) Let $\mathcal{M}=\mathcal{M}_{X}$ be an $\mathcal{O}_{X}$-module. There is a canonical homomorphism $\mathcal{M}_{X} \rightarrow j_{*}\left[\mathcal{M}_{Y}\right]$.
proof. (ii) Let $U$ be an affine open subset of $X$, and let $\mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{P}$ be an exact sequence of $\mathcal{O}_{Y}$-modules. The sequence $j_{*} \mathcal{M}(U) \rightarrow j_{*} \mathcal{N}(U) \rightarrow j_{*} \mathcal{P}(U)$ is the same as the sequence $\mathcal{M}(U \cap Y) \rightarrow \mathcal{N}(U \cap Y) \rightarrow$ $\mathcal{P}(U \cap Y)$, though the scalars have changed. Since $U$ and $Y$ are affine, $U \cap Y$ is affine. By definition of exactness, this sequence is exact.
(iii) Let $U$ be open in $X$. Then $j_{*} \mathcal{M}_{Y}(U)=\mathcal{M}(U \cap Y)$. Since $U \cap Y \subset U, \mathcal{M}(U)$ maps to $\mathcal{M}(U \cap Y)$.
6.5.10. Example. Let $X=\mathbb{P}^{n}$ and let $j$ denote the inclusion $\mathbb{U}^{0} \subset X$ of the standard affine open subset into $X$. The direct image $j_{*} \mathcal{O}_{\mathbb{U}^{0}}$ is the algebra of rational functions that are allowed to have poles on the hyperplane at infinity.

The inverse image of an open subset $W$ of $X$ is its intersection with $\mathbb{U}^{0}: j^{-1} W=W \cap \mathbb{U}^{0}$. The sections of the direct image $j_{*} \mathcal{O}_{\mathbb{U}^{0}}$ on an open subset $W$ of $X$ are the regular functions on $W \cap \mathbb{U}^{0}$ :

$$
\left[j_{*} \mathcal{O}_{\mathbb{U}^{0}}\right](W)=\mathcal{O}_{\mathbb{U}^{0}}\left(W \cap \mathbb{U}^{0}\right)=\mathcal{O}_{X}\left(W \cap \mathbb{U}^{0}\right)
$$

Say that we write a rational function $\alpha$ as a fraction $g / h$ of relatively prime polynomials. Then $\alpha$ is an element of $\mathcal{O}_{X}(W)$ if $h$ doesn't vanish at any point of $W$, and $\alpha$ is a section of $\left[j_{*} \mathcal{O}_{\mathbb{U}^{0}}\right](W)=\mathcal{O}_{X}\left(W \cap \mathbb{U}^{0}\right)$ if $h$ doesn't vanish on $W \cap \mathbb{U}^{0}$. Arbitrary powers of $x_{0}$ can appear in the denominator $h$ when $\alpha$ is a section of $j_{*} \mathcal{O}_{\mathbb{U}^{0}}$.

### 6.6 Support

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(6.6.1) annihilators
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Let $A$ be a ring, and let $m$ be an element of an $A$-module $M$. The annihilator $I$ of an element $m$ of $M$ is the set of elements $\alpha$ of $A$ such that $\alpha m=0$. It is an ideal of $A$ that is often denoted by ann $(m)$.

The annihilator of an $A$-module $M$ is the set of elements of $A$ such that $a M=0$. This annihilator is also an ideal.
6.6.2. Lemma. Let $I$ be the annihilator of an element $m$ of $M$, and let $s$ be a nonzero element of $A$. The annihilator of the image of $m$ in the localized module $M_{s}$ is the localized ideal $I_{s}$

This lemma allows us to extend the concept of annihilator to sections of a finite $\mathcal{O}$-module on a variety $X$.
6.6.3. Support. Let $M$ be a finite module over a finite-type domain $A$ and $X=\operatorname{Spec} A$. The support of $M$ is the locus $C=V_{X}(I)$ of zeros of its annihilator $I$ in $X$. It is the set of points $p$ of $X$ such that $I \subset \mathfrak{m}_{p}$. The support of a finite module is a closed subset of $X$. If $s$ is a nonzero element of $A, \operatorname{dtThe}$ support of the localized module $M_{s}$ is the intersection $C_{s}=C \cap X_{s}$.

Let $\mathcal{M}$ be a finite $\mathcal{O}$-module on a variety $X$, and let $\mathcal{I}$ be its annihilator. The support of $\mathcal{M}$ is the closed subset $V_{X}(\mathcal{I})$ of points such that $\mathcal{I} \subset \mathfrak{m}_{p}$. For example, the support of the residue field module $\kappa_{p}$ is the point $p$. The support of the maximal ideal $\mathfrak{m}_{p}$ at $p$ is the whole variety $X$.

## $\mathcal{O}$-modules with support of dimension zero

6.6.5. Proposition. Let $\mathcal{M}$ be a finite $\mathcal{O}$-module on a variety $X$.
(i) Suppose that the support of $\mathcal{M}$ is a single point $p$, let $M=\mathcal{M}(X)$, and let $U$ be an affine open subset of $X$. If $U$ contains $p$, then $\mathcal{M}(U)=M$, and if $U$ doesn't contain $p$, then $\mathcal{M}(U)=0$.
(ii) (Chinese Remainder Theorem) If the support of $\mathcal{M}$ is a finite set $\left\{p_{1}, \ldots, p_{k}\right\}$, then $\mathcal{M}$ is the direct sum $\mathcal{M}_{1} \oplus \cdots \oplus \mathcal{M}_{k}$ of $\mathcal{O}$-modules supported at the points $p_{i}$.
proof. (i) Let $\mathcal{I}$ be the annihilator of $\mathcal{M}$. The locus $V_{X}(\mathcal{I})$ is $p$. If $p$ isn't contained in $U$, then when we restrict $\mathcal{M}$ to $U$, we obtain an $\mathcal{O}_{U}$-module whose support is empty. Therefore the restriction to $U$ is the zero module.

Next, suppose that $p$ is contained in $U$, and let $V$ denote the complement of $p$ in $X$. We cover $X$ by a set $\left\{U^{i}\right\}$ of affine open sets with $U=U^{1}$, and such that $U^{i} \subset V$ if $i>1$. By what has been shown, $\mathcal{M}\left(U^{i}\right)=0$ if $i>0$ and $\mathcal{M}\left(U^{i j}\right)=0$ if $j \neq i$. The sheaf axiom for this covering shows that $\mathcal{M}(X) \approx \mathcal{M}(U)$.
(ii) This follows from the ordinary Chinese Remainder Theorem.

### 6.7 Twisting

The twisting modules that we define here are among the most important modules on projective space.
Let $X$ denote the projective space $\mathbb{P}^{n}$ with coordinates $x_{0}, \ldots, x_{n}$. As before, a homogeneous fraction of degree $d$ is a fraction $g / h$ of homogeneous polynomials with $\operatorname{deg} g-\operatorname{deg} h=d$. When $g$ and $h$ are relatively prime, the fraction $g / h$ is regular on an open subset $V$ of $X$ if $h$ isn't zero at any point of $V$.

The definition of the twisting modules is this: The sections of $\mathcal{O}(d)$ on an open subset $V$ of $\mathbb{P}^{n}$ are the homogeneous fractions of degree $d$ that are regular on $V$.

### 6.7.1. Proposition.

(i) Let $V$ be an affine open subset of $\mathbb{P}^{n}$ that is contained in the standard affine open set $\mathbb{U}^{0}$. The sections of
(ii) The twisting module $\mathcal{O}(d)$ is an $\mathcal{O}$-module.
proof. (i) Let $V$ be an open set contained in $\mathbb{U}^{0}$, and let $\alpha$ be a section of $\mathcal{O}(d)$ on $V$. Then $f=\alpha x_{0}^{-d}$ has degree zero. It is a rational function. Since $V \subset \mathbb{U}^{0}, x_{0}$ doesn't vanish at any point of $V$. Since $\alpha$ is regular on $V, f$ is a regular function on $V$, and $\alpha=f x_{0}^{d}$.
(ii) It is clear that $\mathcal{O}(d)$ is a contravariant functor. We verify compatibility with localization. Let $V=\operatorname{Spec} A$ be an affine open subset of $X$ and let $s$ be a nonzero element of $A$. We must show that $[\mathcal{O}(d)]\left(V_{s}\right)$ is the localization of $[\mathcal{O}(d)](V)$, and it is true that $[\mathcal{O}(d)](V)$ is a subset of $[\mathcal{O}(d)]\left(V_{s}\right)$. What has to be shown is that if $\beta$ is a section of $\mathcal{O}(d)$ on $V_{s}$, then $s^{k} \beta$ is a section on $V$, if $k$ is sufficiently large.

We cover $V$ by the affine open sets $V^{i}=V \cap \mathbb{U}^{i}$. To show that $s^{k} \beta$ is a section on $V$, it suffices to show that it is a section on $V \cap \mathbb{U}^{i}$ for every $i$. This is the sheaf property. We apply (i) to the open subset $V_{s}^{0}$ of $V^{0}$. Since $V_{s}^{0}$ is contained in $\mathbb{U}^{0}, \beta$ can be written (uniquely) in the form $f x_{0}^{d}$, where $f$ is a rational function that is regular on $V_{s}^{0}$. We know already that $\mathcal{O}$ has the coherence property. Therefore $s^{k} f$ is a regular function on $V^{0}$ if $k$ is large, and then $s^{k} \alpha=s^{k} f x_{0}^{d}$ is a section of $\mathcal{O}(d)$ on $V^{0}$. The analogous statement is true for every index $i$.

As part (i) of the proposition shows, $\mathcal{O}(d)$ is quite similar to the structure sheaf. However, $\mathcal{O}(d)$ is only locally free. Its sections on the standard open set $\mathbb{U}^{1}$ form a free $\mathcal{O}\left(\mathbb{U}^{1}\right)$-module with basis $x_{1}^{d}$. That basis is related to the basis $x_{0}^{d}$ on $\mathbb{U}^{0}$ by the factor $\left(x_{0} / x_{1}\right)^{d}$, a rational function that isn't invertible on on $\mathbb{U}^{0}$ or on $\mathbb{U}^{1}$.
6.7.2. Proposition. When $d \geq 0$, the global sections of the twisting module $\mathcal{O}(d)$ on $\mathbb{P}^{n}$ are the homogeneous polynomials of degree $d$. When $d<0$, the only global section of $\mathcal{O}(d)$ is zero.
proof. A nonzero global section $u$ of $\mathcal{O}(d)$ will restrict to a section on the standard affine open set $\mathbb{U}^{0}$. Since elements of $\mathcal{O}\left(\mathbb{U}^{0}\right)$ are homogeneous fractions of degree zero whose denominators are powers of $x_{0}$, and since $[\mathcal{O}(d)]\left(\mathbb{U}^{0}\right)$ is a free module over $\mathcal{O}\left(\mathbb{U}^{0}\right)$ with basis $x_{0}^{d}, \quad u$ will have the form $g / x_{0}^{m}$ for some some homogeneous polynomial $g$ and some $m$. Similarly, restriction to $\mathbb{U}^{1}$ shows that $u$ has the form $g_{1} / x_{1}^{n}$. It follows that $m=n=0$ and that $u=g$. Since $u$ has degree $d, g$ will have degree $d$.

### 6.7.3. Examples.

The product $u v$ of homogeneous fractions of degrees $r$ and $s$ is a homogeneous fraction of degree $r+s$, and if $u$ and $v$ are regular on an open set $V$, so is their product $u v$. So multiplication defines a homomorphism of $\mathcal{O}$-modules

$$
\begin{equation*}
\mathcal{O}(r) \times \mathcal{O}(s) \rightarrow \mathcal{O}(r+s) \tag{6.7.4}
\end{equation*}
$$

Multiplication by a homogeneous polynomial $f$ of degree $d$ defines an injective homomorphism

$$
\begin{equation*}
\mathcal{O}(k) \xrightarrow{f} \mathcal{O}(k+d) . \tag{6.7.5}
\end{equation*}
$$

When $k=-d$, this becomes a homomorphism $\mathcal{O}(-d) \xrightarrow{f} \mathcal{O}$.

The twisting modules $\mathcal{O}(n)$ have a second interpretation. They are isomorphic to the modules that we denote by $\mathcal{O}(n H)$, of rational functions on projective space with poles of order at most $n$ on the hyperplane $H:\left\{x_{0}=0\right\}$ at infinity.

By definition, the nonzero sections of $\mathcal{O}(n H)$ on an open set $V$ are the rational functions $f$ such that $x_{0}^{n} f$ is a section of $\mathcal{O}(n)$ on $V$. Thus multiplication by $x_{0}^{n}$ defines an isomorphism

$$
\begin{equation*}
\mathcal{O}(n H) \xrightarrow{x_{0}^{n}} \mathcal{O}(n) \tag{6.7.6}
\end{equation*}
$$

If $f$ is a section of $\mathcal{O}(n H)$ on an open set $V$, and if we write $f$ as a homogeneous fraction $g / h$ of degree zero, with $g, h$ relatively prime, the denominator $h$ may have $x_{0}^{k}$, with $k \leq n$, as factor. The other factors of $h$ cannot vanish anywhere on $V$. If $f=g / h$ is a global section of $\mathcal{O}(n H)$, then $h=c x_{0}^{k}$, with $c \in \mathbb{C}$ and $k \leq n$. A global section of $\mathcal{O}(n H)$ can be represented as a fraction $g / x_{0}^{k}$.

Since $x_{0}$ doesn't vanish at any point of the standard affine open set $\mathbb{U}^{0}$, the sections of $\mathcal{O}(n H)$ on an open subset $V$ of $\mathbb{U}^{0}$ are simply the regular functions on $V$. The restrictions of $\mathcal{O}(n H)$ and $\mathcal{O}$ to $\mathbb{U}^{0}$ are equal. Using the subsctript notation (6.5.6 for restriction to an open set,

$$
\begin{equation*}
\mathcal{O}(n H)_{\mathbb{U}^{0}}=\mathcal{O}_{\mathbb{U}^{0}} \tag{6.7.7}
\end{equation*}
$$

Let $V$ be an open subset of one of the other standard affine open sets, say of $\mathbb{U}^{1}$. The ideal of $H \cap \mathbb{U}^{1}$ in $\mathbb{U}^{1}$ is principal, generated by $v_{0}=x_{0} / x_{1}$, and $v_{0}$ generates the ideal of $H \cap V$ in $V$ too. If $f$ is a rational function, then because $x_{1}$ doesn't vanish on $\mathbb{U}^{1}$, the function $f v_{0}^{n}$ will be regular on $V$ if and only if the homogeneous fraction $f x_{0}^{n}$ is regular there. So $f$ will be a section of $\mathcal{O}(n H)$ on $V$ if and only if $f v_{0}^{n}$ is a regular function. Because $v_{0}$ generates the ideal of $H$ in $V$, we say that such a function $f$ has a pole of order at most $n$ on $H$.

The isomorphic $\mathcal{O}$-modules $\mathcal{O}(n)$ and $\mathcal{O}(n H)$ are interchangeable. The twisting module $\mathcal{O}(n)$ is often better because its definition is independent of coordinates. On the other hand, $\mathcal{O}(n H)$ can be convenient because its restriction to $\mathbb{U}^{0}$ is the structure sheaf $\mathcal{O}_{\mathbb{U}^{0}}$.
6.7.8. Proposition. Let $Y$ be the zero locus of an irreducible homogeneous polynomial $f$ of degree $d$ ( $a$ hypersurface of degree $d$ in $\mathbb{P}^{n}$ ), let $\mathcal{I}$ be the ideal of $Y$, and let $\mathcal{O}(-d)$ be the twisting module on $X$. Multiplication by $f$ defines an isomorphism $\mathcal{O}(-d) \xrightarrow{f} \mathcal{I}$.
proof. If $\alpha$ is a section of $\mathcal{O}(-d)$ on an open set $V$, then $f \alpha$ will be a rational function that is regular on $V$ and that vanishes on $Y \cap V$. Therefore the image of the multiplication map $\mathcal{O}(-d) \xrightarrow{f} \mathcal{O}$ is contained in $\mathcal{I}$. This map is injective because $\mathbb{C}\left[x_{0}, \ldots, x_{n}\right]$ is a domain. To show that it is an isomorphism, it suffices to show that its restrictions to the standard affine open sets $\mathbb{U}^{i}$ are isomorphisms 6.3.6. As usual, we work with $\mathbb{U}^{0}$.

We may suppose that $f$ isn't a scalar multiple of any $x_{i}$. Then $Y \cap \mathbb{U}^{0}$ will be a dense open subset of $Y$. The sections of $\mathcal{O}$ on $\mathbb{U}^{0}$ are the homogeneous fractions $g / x_{0}^{k}$ of degree zero. Such a fraction is a section of $\mathcal{I}$ on $\mathbb{U}^{0}$ if and only if $g$ vanishes on $Y \cap \mathbb{U}^{0}$. If so, then since $Y \cap \mathbb{U}^{0}$ is dense in $Y, g$ will vanish on $Y$, and therefore it will be divisible by $f: g=f q$. The sections of $\mathcal{I}$ on $\mathbb{U}^{0}$ have the form $f q / x_{0}^{k}$. They are in the image of the map $\mathcal{O}(-d) \rightarrow \mathcal{I}$.

The proposition has an interesting corollary:
6.7.9. Corollary. When regarded as $\mathcal{O}$-modules, the ideals of all hypersurfaces of degree $d$ are isomorphic.

## (6.7.10) twisting a module

6.7.11. Definition Let $\mathcal{M}$ be an $\mathcal{O}$-module on projective space $\mathbb{P}^{d}$, and let $\mathcal{O}(n)$ be the twisting module. The $n$th $t w i s t$ of $\mathcal{M}$ is defined to be the tensor product $\mathcal{M}(n)=\mathcal{M} \otimes_{\mathcal{O}} \mathcal{O}(n)$. Similarly, $\mathcal{M}(n H)=\mathcal{M} \otimes_{\mathcal{O}} \mathcal{O}(n H)$. Twisting is a functor on $\mathcal{O}$-modules.

If $X$ is a closed subvariety of $\mathbb{P}^{d}$ and $\mathcal{M}$ is an $\mathcal{O}_{X}$-module, $\mathcal{M}(n)$ and $\mathcal{M}(n H)$ are obtained by twisting the extension of $\mathcal{M}$ by zero. (See the equivalence of categories (6.5.5).

Since $x_{0}^{n}$ is a basis of $\mathcal{O}(n)$ on $\mathbb{U}^{0}$, a section of $\mathcal{M}(n)$ on an open subset $V$ of $\mathbb{U}^{0}$ can be written in the form $\mu=m \otimes f x_{0}^{n}$, where $f$ is a regular function on $V$ and $m$ is a section of $\mathcal{M}$ on $V$ 6.7.1). The function $f$ can be moved over to $m$, so $\mu$ can also be written in the form $\mu=m \otimes x_{0}^{n}$, and this expression is unique.
6.7.12. The modules $\mathcal{O}(n)$ and $\mathcal{O}(n H)$ form directed sets that are related by a diagram


In this diagram, the vertical arrows are bijections and the horizontal arrows are injective. The limit of the upper directed set is the module whose sections are allowed to have arbitrary poles on $H$. This is also the module $j_{*} \mathcal{O}_{\mathbb{U}^{0}}$, where $j$ denotes the inclusion of the standard affine open set $\mathbb{U}^{0}$ into $X$ (see 6.5 .9 (iii)):

$$
\begin{equation*}
\xrightarrow[\longrightarrow]{\lim } \mathcal{O}(n H)=j_{*} \mathcal{O}_{\mathbb{U}^{0}} \tag{6.7.14}
\end{equation*}
$$

twistmodule deftwistm

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The next diagram is obtained by tensoring Diagram6.7.13 with $\mathcal{M}$.

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Here, the vertical maps are bijective, but because $\mathcal{M}$ may have torsion, the horizontal maps needn't be injective. However, since tensor products are compatible with limits,

$$
\begin{equation*}
\xrightarrow[\longrightarrow]{\lim } \mathcal{M}(n H)=\underset{\longrightarrow}{\lim } \mathcal{M} \otimes_{\mathcal{O}} \mathcal{O}(n H)=\mathcal{M} \otimes_{\mathcal{O}} j_{*} \mathcal{O} \approx j_{*} \mathcal{M}_{\mathbb{U}^{0}} \tag{6.7.16}
\end{equation*}
$$

The last isomorphism needs explanation:
6.7.17. Lemma. Let $\mathcal{M}$ be an $\mathcal{O}$-module on projective space, and let $\mathbb{U} \xrightarrow{j} \mathbb{P}^{d}$ be the inclusion of the standard affine open set $\mathbb{U}=\mathbb{U}^{0}$. Then $\mathcal{M} \otimes_{\mathcal{O}} j_{*} \mathcal{O}_{\mathbb{U}}$ and $j_{*} \mathcal{M}_{\mathbb{U}}$ are isomorphic.
proof. Let $\alpha \otimes f$ be a section of $\mathcal{M} \otimes j_{*} \mathcal{O}_{\mathbb{U}}$ on an open set $V$, where $\alpha$ is a section of $\mathcal{M}$ on $V$ and and $f$ is a regular function on $V \cap \mathbb{U}$. Then, denoting the restriction of $\alpha$ to $V \cap \mathbb{U}$ by the same symbol, $\alpha f$ will be a section of $\mathcal{M}$ on $V \cap \mathbb{U}$ and therefore a section of $j_{*} \mathcal{M}_{\mathbb{U}}$ on $V$. Sending $\alpha \otimes f$ to $\alpha f$ defines a homomorphism $\mathcal{M} \otimes_{\mathcal{O}} j_{*} \mathcal{O}_{b} b u \rightarrow j_{*} \mathcal{M}_{\mathbb{U}}$. To show that this map is an isomorphism, it suffices to verify that it defines a bijective map on each of the standard affine open sets $\mathbb{U}^{i}$. We leave the case $i=0$ as exercise and look at $\mathbb{U}^{1}$. On that open set, the coordinate algebra is $\mathbb{C}\left[v_{0}, \ldots, v_{d}\right]$ with $v_{=}=x_{i} / x_{1}$. and $v_{1}=1$. Let $M_{0}=\mathcal{M}\left(\mathbb{U}^{0}\right)$. Then $\left[j_{*} \mathcal{M}_{\mathbb{U}}\right]\left(\mathbb{U}^{1}\right)=M_{0}\left[v_{0}^{-1}\right]$. Also, $\left[j_{*} \mathcal{O}_{\mathbb{U}}\right]\left(\mathbb{U}^{1}\right)=\mathcal{O}_{\left(\mathbb{U}^{0} \cap \mathbb{U}^{1}\right) \text {, so by definition of the tensor product, }}$

$$
\left[\mathcal{M} \otimes_{\mathcal{O}} j_{*} \mathcal{O}_{\mathbb{U}^{0}}\right]\left(\mathbb{U}^{1}\right)=\mathcal{M}\left(\mathbb{U}^{1}\right) \otimes_{\mathcal{O}\left(\mathbb{U}^{1}\right)} \mathcal{O}_{\left(\mathbb{U}^{0} \cap \mathbb{U}^{1}\right)=\mathcal{M}\left(b b u^{1} \cap \mathbb{U}^{0}\right)=M_{0}\left[v_{0}^{-1}\right]}
$$

## (6.7.18) generating an $\mathcal{O}$-module

Le $\mathcal{M}$ be an $\mathcal{O}$-module on a variety $X$, and let $m=\left(m_{1}, \ldots, m_{k}\right)$ be a set of global sections of $\mathcal{M}$. This set defines a map

$$
\begin{equation*}
\mathcal{O}^{k} \xrightarrow{m} \mathcal{M} \tag{6.7.19}
\end{equation*}
$$

n that sends a section $\left(\alpha_{1}, \ldots, \alpha_{k}\right)$ of $\mathcal{O}^{k}$ on an open set to the combination $\sum \alpha_{i} m_{i}$. The global sections $m_{1}, \ldots, m_{k}$ generate $\mathcal{M}$ if this map is surjective. If the sections generate $\mathcal{M}$, then they (more precisely, their restrictions), generate the $\mathcal{O}(U)$-module $\mathcal{M}(U)$ for every affine open set $U$. They may fail to generate $\mathcal{M}(U)$ when $U$ isn't affine.
6.7.20. Example. Let $X=\mathbb{P}^{1}$. For $n \geq 0$, the global sections of the twisting module $\mathcal{O}(n)$ are the polynomials of degree $n$ in the coordinate variables $x_{0}, x_{1}$ 6.7.2. Consider the map $\mathcal{O}^{2} \xrightarrow{\left(x_{0}^{n}, x_{1}^{n}\right)} \mathcal{O}(n)$. On $\mathbb{U}^{0}, \mathcal{O}(n)$ has basis $x_{0}^{n}$. Therefore this map is surjective on $\mathbb{U}^{0}$. Similarly, it is surjective on $\mathbb{U}^{1}$. So it is a surjective map on all of $X$ 6.3.6. The global sections $x_{0}^{n}, x_{1}^{n}$ generate $\mathcal{O}(n)$. However, the global sections of $\mathcal{O}(n)$ are the homogeneous polynomials of degree $n$. When $n>1$, the two sections $x_{0}^{n}, x_{1}^{n}$ don't span the space of global sections, and the map $\mathcal{O}^{2} \xrightarrow{\left(x_{0}^{n}, x_{1}^{n}\right)} \mathcal{O}(n)$ isn't surjective.

The next theorem explains the importance of the twisting operation.
6.7.21. Theorem. Let $\mathcal{M}$ be a finite $\mathcal{O}$-module on a projective variety $X$. For large $k$, the twist $\mathcal{M}(k)$ is generated by global sections.
proof. We may assume that $X$ is projective space $\mathbb{P}^{n}$.
We are to show that if $\mathcal{M}$ is a finite $\mathcal{O}$-module, the global sections generate $\mathcal{M}(k)$ when $k$ is large, and it suffices to show that for each $i=0, \ldots, n$, the restrictions of those global sections to $\mathbb{U}^{i}$ generate the $\mathcal{O}\left(\mathbb{U}^{i}\right)$ module $[\mathcal{M}(k)]\left(\mathbb{U}^{i}\right)$ 6.3.6). We work with the index 0 as usual.

We replace $\mathcal{M}(k)$ by the isomorphic module $\mathcal{M}(k H)$. We recall that the restriction of $\mathcal{M}(k H)$ to $\mathbb{U}^{0}$ is equal to $\mathcal{M}_{\mathbb{U}^{0}}$. and that $\xrightarrow{\lim } \mathcal{M}(k H)=j_{*} \mathcal{M}_{\mathbb{U}^{0}}$.

Let $A_{0}=\mathcal{O}\left(\mathbb{U}^{0}\right)$, and $M_{0}=\mathcal{M}\left(\mathbb{U}^{0}\right)$. The global sections of $j_{*} \mathcal{M}_{\mathbb{U}^{0}}$ are the sections of $\mathcal{M}$ on $\mathbb{U}^{0}$, which we are denoting by $M_{0}$. We choose a finite set of generators $m_{1}, \ldots, m_{r}$ for the finite $A_{0}$-module $M_{0}$. These generators are global sections of $j_{*} \mathcal{M}_{\mathbb{U}^{0}}$, and since $\underset{\longrightarrow}{\lim } \mathcal{M}(k H)=j_{*} \mathcal{M}_{\mathbb{U}^{0}}$, the sections $m_{i}$ are represented by global sections $m_{1}^{\prime}, \ldots, m_{r}^{\prime}$ of $\mathcal{M}(k H)$ when $k$ is large. The restrictions of $m_{i}$ and $m_{i}^{\prime}$ to $\mathbb{U}^{0}$ are equal, so the restrictions of $m_{1}^{\prime}, \ldots, m_{r}^{\prime}$ generate $M_{0}$ too. Then $M_{0}$ is generated by global sections of $\mathcal{M}(k H)$, as was to be shown.

### 6.8 Proof of Theorem 6.3.2.

The statement to be proved is that an $\mathcal{O}$-module $\mathcal{M}$ on a variety $X$ has a unique extension to a functor

$$
\text { (opens) } \xrightarrow{\widetilde{\mathcal{M}}}(\text { modules })
$$

that has the sheaf property 6 6.3.5,$(6.3 .10)$ and that a homomorphism $\mathcal{M} \rightarrow \mathcal{N}$ of $\mathcal{O}$-modules has a unique extension to a homomorphism $\mathcal{M} \rightarrow \mathcal{N}$.

The proof has the following steps:

1. Verification of the sheaf property for a covering of an affine open set by localizations.
2. Extension of the functor $\mathcal{M}$ to all morphisms between affine open sets.
3. Definition of $\widetilde{\mathcal{M}}$.

## Step 1. (the sheaf property for a covering of an affine open set by localizations)

\#\#\#has this been done before??\#\#\#
Suppose that an affine open set $Y=\operatorname{Spec} A$ is covered by a family of localizations $\mathbb{U}_{0}=\left\{U^{s_{i}}\right\}$, et $\mathcal{M}$ be an $\mathcal{O}$-module, and let $M=\mathcal{M}(Y), M_{i}=\mathcal{M}\left(U^{s_{i}}\right)$, and $M_{i j}=\mathcal{M}\left(U^{s_{i} s_{j}}\right)$. The exact sequence that expresses the sheaf property for the covering diagram $Y \longleftarrow \mathbb{U}_{0} \leftleftarrows \mathbb{U}_{1}$ becomes

$$
\begin{equation*}
0 \rightarrow M \stackrel{\alpha}{\longrightarrow} \prod M_{i} \xrightarrow{\beta} \prod M_{i j} \tag{6.8.1}
\end{equation*}
$$

In this sequence, the map $\alpha$ sends an element $m$ of $M$ to the vector $(m, \ldots, m)$ of its images in $\prod_{i} M_{i}$, and $\beta$ sends a vector $\left(m_{1}, \ldots, m_{k}\right)$ in $\prod_{i} M_{i}$ to the matrix $\left(z_{i j}\right)$, with $z_{i j}=m_{j}-m_{i}$ in $M_{i j}$ To be precise, $M_{i}$ and $M_{j}$ map to $M_{i j}$, and $z_{i j}$ is the difference of their images. We must show that the sequence (6.8.1) is exact.
exactness at $M$ : Since the open sets $U^{i}$ cover $Y$, the elements $s_{1}, \ldots, s_{k}$ generate the unit ideal. Let $m$ be an element of $M$ that maps to zero in every $M_{i}$. Then there exists an $n$ such that $s_{i}^{n} m=0$, and we can use the same exponent $n$ for all $i$. The elements $s_{i}^{n}$ generate the unit ideal. Writing $\sum a_{i} s_{i}^{n}=1$, we have $m=\sum a_{i} s_{i}^{n} m=\sum a_{i} 0=0$.
exactness at $\prod M_{i}$ : Let $m_{i}$ be elements of $M_{i}$ such that $m_{j}=m_{i}$ in $M_{i j}$ for all $i, j$. We must find an element $w$ in $M$ that maps to $m_{j}$ in $M_{j}$ for every $j$.

We write $m_{i}$ as a fraction: $m_{i}=s_{i}^{-n} x_{i}$, or $x_{i}=s_{i}^{n} m_{i}$, with $x_{i}$ in $M$, using the same integer $n$ for all $i$. The equation $m_{j}=m_{i}$ in $M_{i j}$ tells us that $s_{i}^{n} x_{j}=s_{j}^{n} x_{i}$ in $M_{i j}$, and then $\left(s_{i} s_{j}\right)^{r} s_{i}^{n} x_{j}=\left(s_{i} s_{j}\right)^{r} s_{j}^{n} x_{i}$ will be true in $M$, if $r$ is large (see 2.6.11.

We adjust the notation. Let $\widetilde{x}_{i}=s_{i}^{r} x_{i}$, and $\widetilde{s}_{i}=s_{i}^{r+n}$. Then in $M, \widetilde{x}_{i}=\widetilde{s}_{i} m_{i}$ and $\widetilde{s}_{j} \widetilde{x}_{i}=\widetilde{s}_{i} \widetilde{x}_{j}$. The elements $\widetilde{s}_{i}$ generate the unit ideal So there is an equation in $A$, of the form $\sum a_{i} \widetilde{s}_{i}=1$. Let $w=\sum a_{i} \widetilde{x}_{i}$. This is an element of $M$, and

$$
\widetilde{x}_{j}=\left(\sum_{i} a_{i} \widetilde{s}_{i}\right) \widetilde{x}_{j}=\sum_{i} a_{i} \widetilde{s}_{j} \widetilde{x}_{i}=\widetilde{s}_{j} w
$$

Since $m_{j}=\widetilde{s}_{j}^{-1} \widetilde{x}_{j}, m_{j}=w$ is true in $M_{j}$. Since $j$ is arbitrary, $w$ is the required element of $M$.
Step 2. (extending an $\mathcal{O}$-module to all morphisms between affine open sets)
The $\mathcal{O}$-module $\mathcal{M}$ comes with localization maps $\mathcal{M}(U) \rightarrow \mathcal{M}\left(U_{s}\right)$. It doesn't come with homomorphisms $\mathcal{M}(U) \rightarrow \mathcal{M}(V)$ when $V \rightarrow U$ is an arbitrary inclusion of affine open sets. We define those maps here.

Let $\mathcal{M}$ be an $\mathcal{O}$-module and let $V \rightarrow U$ be an inclusion of affine open sets. To describe the canonical homomorphism $\mathcal{M}(U) \rightarrow \mathcal{M}(V)$, we cover $V$ by a family $\mathbb{V}_{0}=\left\{V^{i}\right\}$ of open sets that are localizations of $U$ and therefore also localizations of $V$, and we inspect the covering diagram $V \leftarrow \mathbb{V}_{0} \leftleftarrows \mathbb{V}_{1}$ and the corresponding exact sequence $0 \rightarrow \mathcal{M}(V) \xrightarrow{\alpha} \mathcal{M}\left(\mathbb{V}_{0}\right) \xrightarrow{\beta} \mathcal{M}\left(\mathbb{V}_{1}\right)$. We add the map $V \rightarrow U$ to the covering diagram:

$$
U \leftarrow V \leftarrow \mathbb{V}_{0} \leftleftarrows \mathbb{V}_{1}
$$

Since $V^{i}$ are localizations of $U$ and $V^{i j}$ are localizations of $V^{i}$ and of $V^{j}$, the $\mathcal{O}$-module $\mathcal{M}$ comes with maps $\mathcal{M}(U) \xrightarrow{\psi} \mathcal{M}\left(\mathbb{V}_{0}\right) \rightrightarrows \mathcal{M}\left(\mathbb{V}_{1}\right)$. Let $\psi$ be the map $\mathcal{M}(U) \rightarrow \mathcal{M}\left(\mathbb{V}_{0}\right)$. The two composed maps $U \leftarrow V \leftleftarrows \mathbb{V}_{1}$ are equal, and so are the maps $\mathcal{M}(U) \rightrightarrows \mathcal{M}\left(\mathbb{V}_{1}\right)$. Their difference, which is $\beta \psi$, is the zero map. Therefore $\psi$ maps $\mathcal{M}(U)$ to the kernel of $\beta$ which, according to Step 1, is $\mathcal{M}(V)$. This gives us a map $\mathcal{M}(U) \xrightarrow{\eta} \mathcal{M}(V)$ that makes a diagram


Both $\psi$ and $\alpha$ are compatible with multiplication by a regular function $f$ on $U$, and $\alpha$ is injective. So $\eta$ is also compatible with multiplication by $f$.

We must check that $\eta$ is independent of the covering $\mathbb{V}_{0}$. Let $\mathbb{V}_{0}^{\prime}=\left\{V^{\prime j}\right\}$ be another covering of $V$ by localizations of $U$. We cover each of the open sets $V^{i} \cap V^{\prime j}$ by localizations $W^{i j \nu}$ of $U$. Taken together, these open sets form a covering $\mathbb{W}_{0}$ of $V$. We have a map $\mathbb{W}_{0} \xrightarrow{\epsilon} \mathbb{V}_{0}$ that gives us a diagram

and therefore a diagram

whose rows are exact sequences. Here $\mathcal{M}(U)$ maps to the kernels of $\beta_{\mathbb{V}}$ and $\beta_{\mathbb{W}}$, both of which are equal to $\mathcal{M}(V)$. Looking at the diagram, one sees that the map $\mathcal{M}(U) \rightarrow \mathcal{M}\left(\mathbb{W}_{0}\right)$ is the composition of the maps $\mathcal{M}(U) \rightarrow \mathcal{M}\left(\mathbb{V}_{0}\right) \rightarrow \mathcal{M}\left(\mathbb{W}_{0}\right)$. Therefore the maps $\mathcal{M}(U) \rightarrow \mathcal{M}(V)$ defined using the two coverings $\mathbb{V}_{0}$ and $\mathbb{W}_{0}$ are equal.

We show that this extended functor has the sheaf property for an affine covering $\mathbb{U}_{0}=\left\{U^{i}\right\}$ of an affine variety $U$. Let $\mathbb{V}_{0}$ be a covering of $U$ that is obtained by covering each $U^{i}$ by localizations of $U$. This gives us a diagram


Because $\mathbb{V}_{0}$ covers $\mathbb{U}_{0}, \mathbb{V}_{1}$ covers $\mathbb{U}_{1}$ as well. So the maps $\beta$ and $\gamma$ are injective. Step 1 tells us that the bottom row is exact. Then Lemma 6.3.16(ii) shows that the top row is exact.

## Step 3. (definition of $\widetilde{\mathcal{M}}$ )

Let $Y$ be an open subset of $X$. We use the sheaf property to define $\widetilde{\mathcal{M}}(Y)$. We choose a (finite) covering $\mathbb{U}_{0}=\left\{U^{i}\right\}$ of $Y$ by affine open sets, and we define $\widetilde{\mathcal{M}}(Y)$ to be the kernel $K_{\mathbb{U}}$ of the map $\mathcal{M}\left(\mathbb{U}_{0}\right) \xrightarrow{\beta_{\mathbb{U}}}$
$\mathcal{M}\left(\mathbb{U}_{1}\right)$, where $\beta_{\mathbb{U}}$ is the map described in (??). When we show that this kernel is independent of the covering $\mathbb{U}_{0}$, it will follow that $\widetilde{\mathcal{M}}$ is well-defined, and that it has the sheaf property.

Let $\mathbb{W}_{0}=\left\{W^{\nu}\right\}$ be another covering of $Y$ by affine open sets. One can go from $\mathbb{U}_{0}$ to $\mathbb{V}_{0}$ and back in a finite number of steps, each of which changes a covering by adding or deleting a single affine open set.

We consider a family $\mathbb{W}_{0}=\left\{U^{i}, V\right\}$ obtained by adding one affine open subset $V$ of $Y$ to $\mathbb{U}_{0}$, and we let $\mathbb{W}_{1}$ be the family of intersections of pairs of elements of $\mathbb{W}_{0}$. Then with notation as above, we have a map $K_{\mathbb{W}} \rightarrow K_{\mathbb{U}}$. We show that, for any element $\left(u_{i}\right)$ in the kernel $K_{\mathbb{U}}$, there is a unique element $v$ in $\mathcal{M}(V)$ such that $\left(\left(u_{i}\right), v\right)$ is in the kernel $K_{\mathbb{W}}$. This will show that $K_{\mathbb{W}}=K_{\mathbb{U}}$.

To define the element $v$, we let $V^{i}=U^{i} \cap V$. Since $\mathbb{U}_{0}=\left\{U^{i}\right\}$ is a covering of $Y, \mathbb{V}_{0}=\left\{V^{i}\right\}$ is a covering of $V$ by affine open sets. Let $v_{i}$ be the restriction of the section $u_{i}$ to $V^{i}$. Since $\left(u_{i}\right)$ is in the kernel of $\beta_{\mathbb{U}}, u_{i}=u_{j}$ on $U^{i j}$. Then it is also true that $v_{i}=v_{j}$ on the smaller open set $V^{i j}$. So $\left(v_{i}\right)$ is in the kernel of the map $\mathcal{M}\left(\mathbb{V}_{0}\right) \xrightarrow{\beta_{\mathbb{V}}} \mathcal{M}\left(\mathbb{V}_{1}\right)$, and since $\mathbb{V}_{0}$ is a covering of the affine variety $V$ by affine open sets, Step 2 tells us that the kernel is $\mathcal{M}(V)$. So there is a unique element $v$ in $\mathcal{M}(V)$ that restricts to $v_{i}$ on $V^{i}$. With this element $V, \quad\left(u_{i}, v\right)$ is in the kernel of $\beta_{\mathbb{W}}$.

When the subsets in the family $\mathbb{W}_{1}$ are listed in the order

$$
\mathbb{W}_{1}=\left\{U^{i} \cap U^{j}\right\},\left\{V \cap U^{j}\right\},\left\{U^{i} \cap V\right\},\{V \cap V\}
$$

the map $\beta_{\mathbb{W}}$ sends $\left(\left(u_{i}\right), v\right)$ to $\left(\left(u_{j}-u_{i}\right),\left(u_{j}-v\right),\left(v-u_{i}\right), 0\right)$, restricted appropriately. Here $u_{i}=u_{j}$ on $U^{i} \cap U^{j}$ because $\left(u_{i}\right)$ is in the kernel of $\beta_{\mathbb{U}}$, and $u_{j}=v_{j}=v$ on $U^{j} \cap V=V^{i}$ by definition.

This completes the proof of Theorem 6.3.2

## Chapter 7 COHOMOLOGY

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7.1 Cohomology
7.2 Complexes
7.3 Characteristic Properties of Cohomology
7.4 Existence of Cohomology
7.5 Cohomology of the Twisting Modules
7.6 Cohomology of Hypersurfaces
7.7 Three Theorems about Cohomology
7.8 Bézout's Theorem

### 7.1 Cohomology

cohqcoh coefficient cohomology $H^{q}(X, \mathbb{Z})$ in the Zariski topology is zero for all $q>0$.

The functor
that carries an $\mathcal{O}$-module $\mathcal{M}$ to $H^{0}(X, \mathcal{M})$ is left exact: If
is an exact sequence of $\mathcal{O}$-modules, the associated sequence of global sections sequence of functors ( $\mathcal{O}$-modules) $\xrightarrow{H^{q}}$ (vector spaces),

In a classic 1956 paper "Faisceaux Algébriques Cohérents", Serre showed that the Zariski topology could be used to define cohomology of $\mathcal{O}$-modules on a variety. This cohomology is the topic of the chapter. By the way, the Zariski topology has limited use for cohomology with other coefficients. In particular, the constant

Let $\mathcal{M}$ be an $\mathcal{O}$-module on a variety $X$. The zero-dimensional cohomology of $\mathcal{M}$ is the space $\mathcal{M}(X)$ of its global sections. When speaking of cohomology, one denotes that space by $H^{0}(X, \mathcal{M})$.

$$
(\mathcal{O} \text {-modules }) \xrightarrow{H^{0}} \text { (vector spaces) }
$$

$$
0 \rightarrow \mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{P} \rightarrow 0
$$

$$
0 \rightarrow H^{0}(X, \mathcal{M}) \rightarrow H^{0}(X, \mathcal{N}) \rightarrow H^{0}(X, \mathcal{P})
$$

is exact. But unless $X$ is affine, the map $H^{0}(X, \mathcal{N}) \rightarrow H^{0}(X, \mathcal{P})$ needn't be surjective. The cohomology is a

$$
H^{0}, H^{1}, H^{2}, \ldots
$$

beginning with $H^{0}$, one for each dimension, that compensates for the lack of exactness in the following way:
(a) To every short exact sequence (7.1.1) of $\mathcal{O}$-modules, there is an associated long exact cohomology sequence

$$
\begin{aligned}
& 0 \rightarrow H^{0}(X, \mathcal{M}) \\
& \xrightarrow{\delta^{0}} H^{0}(X, \mathcal{N}) \rightarrow H^{0}(X, \mathcal{P}) \xrightarrow{\delta^{0}} \\
& \\
&\left.H^{1}\right) \rightarrow H^{1}(X, \mathcal{N}) \rightarrow H^{1}(X, \mathcal{P}) \xrightarrow{\delta^{1}}
\end{aligned}
$$

$$
\xrightarrow{\delta^{q-1}} H^{q}(X, \mathcal{M}) \rightarrow H^{q}(X, \mathcal{N}) \rightarrow H^{q}(X, \mathcal{P}) \xrightarrow{\delta^{q}} \cdots
$$

The maps $\delta^{q}$ in this sequence are the coboundary maps.
(b) Given a (commutative) diagram

whose rows are short exact sequences of $\mathcal{O}$-modules, there is a map of cohomology sequences


Most of this diagram arises from the fact that the $H^{q}$ are functors. The only property that doesn't follow is that the squares

deltadia-
gram
deltadia-
gram
that involve the coboundary maps $\delta$ commute.
Thus a map of exact sequences of $\mathcal{O}$-modules induces a map of cohomology sequences.
A sequence $H^{q}, q=0,1, \ldots$ of functors from $\mathcal{O}$-modules to vector spaces that comes with long cohomology sequences for every short exact sequence of $\mathcal{O}$-modules is called a cohomological functor .

Unfortunately, there is no natural construction of the cohomology. We present a construction in Section 7.4, but it isn't canonical. Though one needs to look at an explicit construction at times, it is usually best to work with the characteristic properties that are described below, in the Section 7.3 .

The cohomology $H^{1}$ in dimension one has an interesting interpretation that you can read about if you like. We won't use it. The higher cohomology $H^{q}$ has no useful direct interpretation.

### 7.2 Complexes

We need complexes because they are used in the construction of cohomology.
A complex of vector spaces is a sequence of homomorphisms of vector spaces

$$
\begin{equation*}
\cdots \rightarrow V^{n-1} \xrightarrow{d^{n-1}} V^{n} \xrightarrow{d^{n}} V^{n+1} \xrightarrow{d^{n+1}} \cdots \tag{7.2.1}
\end{equation*}
$$

indexed by the integers, such that the composition $d^{n} d^{n-1}$ of adjacent maps is zero: The image of $d^{n-1}$ is contained in the kernel of $d^{n}$. Such a complex may be denoted by $V^{\bullet}$.

The $q$-dimensional cohomology of a complex $V^{\bullet}$ is the quotient

$$
\begin{equation*}
\mathbf{C}^{q}\left(V^{\bullet}\right)=\left(\operatorname{ker} d^{q}\right) /\left(\operatorname{im} d^{q-1}\right) \tag{7.2.2}
\end{equation*}
$$

An exact sequence of vector spaces is a complex whose cohomology is zero.

A finite sequence of homomorphisms $7.2 .1,, V^{k} \xrightarrow{d^{k}} V^{k+1} \ldots \xrightarrow{d^{\ell-1}} V^{\ell}$, such that the compositions $d^{n} d^{n-1}$ are zero can be made into a complex by defining $V^{n}=0$ for all other integers $n$. In all of the complexes we consider, $V^{q}$ will be zero when $q<0$.

A homomorphism of vector spaces $V^{0} \xrightarrow{d^{0}} V^{1}$ can be made into the complex

$$
\cdots \rightarrow 0 \rightarrow V^{0} \xrightarrow{d^{0}} V^{1} \rightarrow 0 \rightarrow \cdots
$$

For this complex, $\mathbf{C}^{0}=\operatorname{ker} d^{0}, \mathbf{C}^{1}=$ coker $d^{0}$, and $\mathbf{C}^{q}=0$ for all other $q$.

A map $V^{\bullet} \xrightarrow{\varphi} V^{\prime \bullet}$ of complexes is a collection of homomorphisms $V^{n} \xrightarrow{\varphi^{n}} V^{\prime n}$ making a diagram


A map of complexes induces maps on the cohomology

$$
\mathbf{C}^{q}\left(V^{\bullet}\right) \rightarrow \mathbf{C}^{q}\left(V^{\prime \bullet}\right)
$$

because $\operatorname{ker} d^{q}$ maps to $\operatorname{ker} d^{\prime q}$ and $\operatorname{im} d^{q}$ maps to $\operatorname{im} d^{\prime q}$.
A sequence
exseqcplx
mapcplx
cohcplx
snakecohomology

$$
\begin{equation*}
\cdots \rightarrow V^{\bullet} \xrightarrow{\varphi} V^{\prime \bullet} \xrightarrow{\psi} V^{\prime \prime \bullet} \rightarrow \cdots \tag{7.2.3}
\end{equation*}
$$

of maps of complexes consists of maps

$$
\begin{equation*}
\cdots \rightarrow V^{q} \xrightarrow{\varphi^{q}} V^{\prime q} \xrightarrow{\psi^{q}} V^{\prime \prime q} \rightarrow \cdots \tag{7.2.4}
\end{equation*}
$$

for each $q$. A sequence of maps of complexes is exact if those sequences (7.2.4) are exact for every $q$.
7.2.5. Proposition.

Let $0 \rightarrow V^{\bullet} \rightarrow V^{\bullet \bullet} \rightarrow V^{\prime \prime \bullet} \rightarrow 0$ be a short exact sequence of complexes. For every $q$, there are maps $\mathbf{C}^{q}\left(V^{\prime \prime \bullet}\right) \xrightarrow{\delta^{q}} \mathbf{C}^{q+1}\left(V^{\bullet}\right)$ such that the sequence

$$
\rightarrow \mathbf{C}^{0}\left(V^{\bullet}\right) \rightarrow \mathbf{C}^{0}\left(V^{\prime \bullet}\right) \rightarrow \mathbf{C}^{0}\left(V^{\prime \prime \bullet}\right) \xrightarrow{\delta^{0}} \mathbf{C}^{1}\left(V^{\bullet}\right) \rightarrow \mathbf{C}^{1}\left(V^{\prime \bullet}\right) \rightarrow \mathbf{C}^{1}\left(V^{\prime \prime \bullet}\right) \xrightarrow{\delta^{1}} \mathbf{C}^{2}\left(V^{\bullet}\right) \rightarrow \cdots
$$

is exact.
The proof of the proposition is below.
The long exact sequence is the cohomology sequence associated to the short exact sequence of complexes. This proposition makes the set of functors $\left\{\mathbf{C}^{q}\right\}$ into a cohomological functor on the category of complexes.
7.2.6. Example. We make the Snake Lemma 5.1.2 into a cohomology sequence. Suppose given a diagram

with exact rows. We form the complex $0 \rightarrow V \xrightarrow{f} W \rightarrow 0$ with $V$ in degree zero, so that $\mathbf{C}^{0}\left(V^{\bullet}\right)=\operatorname{ker} f$ and $\mathbf{C}^{1}\left(V^{\bullet}\right)=$ coker $f$. When we do the analogous thing for the maps $f^{\prime}$ and $f^{\prime \prime}$, the Snake Lemma becomes an exact sequence

$$
\mathbf{C}^{0}\left(V^{\bullet}\right) \rightarrow \mathbf{C}^{0}\left(V^{\prime \bullet}\right) \rightarrow \mathbf{C}^{0}\left(V^{\prime \prime \bullet}\right) \rightarrow \mathbf{C}^{1}\left(V^{\bullet}\right) \rightarrow \mathbf{C}^{1}\left(V^{\prime \bullet}\right) \rightarrow \mathbf{C}^{1}\left(V^{\prime \prime \bullet}\right)
$$

proof of Proposition 7.2.5 Let

$$
V^{\bullet}=\left\{\cdots \rightarrow V^{q-1} \xrightarrow{d^{q-1}} V^{q} \xrightarrow{d^{q}} V^{q+1} \xrightarrow{d^{q+1}} \cdots\right\}
$$

be a complex, let $B^{q}$ be the image of $d^{q-1}$, and let $Z^{q}$ be the kernel of $d^{q}$. So $B^{q} \subset Z^{q} \subset V^{q}$. The cohomology of the complex is $\mathbf{C}^{q}\left(V^{\bullet}\right)=Z^{q} / B^{q}$. Let $D^{q}$ be the cokernel $V^{q} / B^{q}$ of $d^{q-1}$, o that we have an exact sequence

$$
0 \rightarrow B^{q} \rightarrow V^{q} \rightarrow D^{q} \rightarrow 0
$$

The map $V^{q} \xrightarrow{d^{q}} V^{q+1}$ factors therouh $D^{q}$ because $B^{q} \subset Z^{q}$, and its image is in $Z^{q+1}$. This gives us a map $D^{q} \xrightarrow{f^{q}} Z^{q+1}$ such that $d^{q}$ is the composition of three maps

$$
V^{q} \xrightarrow{\pi^{q}} D^{q} \xrightarrow{f^{q}} Z^{q+1} \xrightarrow{i^{q+1}} V^{q+1}
$$

where $\pi^{q}$ is the projection from $V^{q}$ to $D^{q}$ and $i^{q+1}$ is the inclusion of $Z^{q+1}$ into $V^{q+1}$. Studying these maps, one sees that

$$
\begin{equation*}
\mathbf{C}^{q}\left(V^{\bullet}\right)=\operatorname{ker} f^{q} \quad \text { and } \quad \mathbf{C}^{q+1}\left(V^{\bullet}\right)=\operatorname{coker} f^{q} . \tag{7.2.7}
\end{equation*}
$$

hiskerandcoker

Suppose given a short exact sequence of complexes $0 \rightarrow V^{\bullet} \rightarrow V^{\prime \bullet} \rightarrow V^{\prime \prime} \rightarrow 0$ as in the proposition. In the diagram below, the rows are exact because cokernel is a right exact functor and kernel is a left exact functor.


When we apply 7.2.7) and the Snake Lemma to this diagram, we obtain an exact sequence

$$
\mathbf{C}^{q}\left(V^{\bullet}\right) \rightarrow \mathbf{C}^{q}\left(V^{\prime \bullet}\right) \rightarrow \mathbf{C}^{q}\left(V^{\prime \prime \bullet}\right) \xrightarrow{\delta^{q}} \mathbf{C}^{q+1}\left(V^{\bullet}\right) \rightarrow \mathbf{C}^{q+1}\left(V^{\prime \bullet}\right) \rightarrow \mathbf{C}^{q+1}\left(V^{\prime \prime \bullet}\right)
$$

The cohomology sequence associated to the short exact sequence of complexes is obtained by splicing these sequences together.

The coboundary maps $\delta^{q}$ in cohomology sequences are related in a natural way. If

is a diagram of maps of complexes whose rows are short exact sequences, the diagrams

commute. It isn't difficult to check this. Thus a map of short exact sequences induces a map of cohomology sequences.

### 7.3 Characteristic Properties of Cohomology

The cohomology $H^{q}(X, \cdot)$ of $\mathcal{O}$-modules, the sequence of functors $H^{0}, H^{1}, H^{2}, \cdots$ from ( $\mathcal{O}$-modules) to (vector spaces), is characterized by the three properties below. The first two have already been mentioned.
charpro-
existco-
cohzeroaffine

1. $H^{0}(X, \mathcal{M})$ is the space $\mathcal{M}(X)$ of global sections of $\mathcal{M}$.
2. The sequence $H^{0}, H^{1}, H^{2}, \cdots$ is a cohomological functor on $\mathcal{O}$-modules: A short exact sequence of $\mathcal{O}$-modules produces a long exact cohomology sequence.
3. Let $Y \xrightarrow{f} X$ be the inclusion of an affine open subset $Y$ into $X$, let $\mathcal{N}$ be an $\mathcal{O}_{Y}$-module, and let $f_{*} \mathcal{N}$ be its direct image on $X$. The cohomology $H^{q}\left(X, f_{*} \mathcal{N}\right)$ is zero for all $q>0$.

When $Y$ is an affine variety, the global section functor is exact. If $0 \rightarrow \mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{P} \rightarrow 0$ is a short exact sequence of $\mathcal{O}$-modules on $Y$, the sequence

$$
0 \rightarrow H^{0}(Y, \mathcal{M}) \rightarrow H^{0}(Y, \mathcal{N}) \rightarrow H^{0}(Y, \mathcal{P}) \rightarrow 0
$$

is exact. There is no need for the higher cohomology $H^{q}$. One may as well define $H^{q}(Y, \cdot)=0$ to be zero when $Y$ is affine and $q>0$. The third characteristic property is based on this observation.
7.3.2. Theorem. There exists a cohomology theory with the properties (7.3.1), and it is unique up to unique isomorphism.

The proof is in the next section.
7.3.3. Corollary. If $X$ is an affine variety, $H^{q}(X, \mathcal{M})=0$ for all $\mathcal{O}$-modules $\mathcal{M}$ and all $q>0$.

This follows when one applies the third characteristic property to the identity map $X \rightarrow X$.
7.3.4. Example. Let $j$ be the inclusion of the standard affine open set $\mathbb{U}^{0}$ into projective space $X$. The third property tells us that the cohomology $H^{q}\left(X, j_{*} \mathcal{O}_{\mathbb{U}^{0}}\right)$ of the direct image $j_{*} \mathcal{O}_{\mathbb{U}^{0}}$ is zero when $q>0$. The direct image is isomorphic to the limit $\underset{\longrightarrow}{\lim } \mathcal{O}_{X}(n H)$ 6.7.13. We will see below 7.4 .26 that cohomology commutes with direct limits. Therefore $\underset{\longrightarrow}{\lim } H^{q}\left(X, \mathcal{O}_{X}(n H)\right.$ and $\underset{\longrightarrow}{\lim } H^{q}\left(X, \mathcal{O}_{X}(n)\right)$ are zero when $q>0$. This will be useful.

Intuitively, the third property tells us that allowing poles on the complement of an affine open set kills cohomology in positive dimension.

### 7.4 Existence of Cohomology

The proof of existence of cohomology and its uniqueness are based on the following facts:

- The intersection of two affine open subsets of a variety is an affine open set.
- A sequence $\cdots \rightarrow \mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{P} \rightarrow \cdots$ of $\mathcal{O}$-modules on a variety $X$ is exact if and only if, for every affine open subset $U$, the sequence of sections $\cdots \rightarrow \mathcal{M}(U) \rightarrow \mathcal{N}(U) \rightarrow \mathcal{P}(U) \rightarrow \cdots$ is exact. This is the definition of exactness.

We begin by choosing an arbitrary affine covering $\mathbb{U}=\left\{U^{\nu}\right\}$ of our variety $X$ by finitely many affine open sets $U^{\nu}$, and we use this covering to describe the cohomology. When we have shown that the cohomology is unique, we will know that it doesn't depend on our choice of covering.

Let $\mathbb{U} \xrightarrow{j} X$ denote the inclusions $U^{\nu} \xrightarrow{j^{\nu}} X$ of our chosen affine open sets into $X$. If $\mathcal{M}$ is an $\mathcal{O}$ module, $\mathcal{R}_{\mathcal{M}}$ will denote the $\mathcal{O}$-module $j_{*} \mathcal{M}_{\mathbb{U}}=\prod j_{*}^{\nu} \mathcal{M}_{U^{\nu}}$, where $\mathcal{M}_{U^{\nu}}$ is the restriction of $\mathcal{M}$ to $U^{\nu}$. As has been noted (6.5.9), there is a canonical map $\mathcal{M} \rightarrow j_{*}^{\nu} \mathcal{M}_{U^{\nu}}$, and therefore a canonical map $\mathcal{M} \rightarrow \mathcal{R}_{\mathcal{M}}$.
7.4.1. Lemma. (i) Let $X^{\prime}$ be an open subset of $X$. The module of sections $\mathcal{R}_{\mathcal{M}}\left(X^{\prime}\right)$ of $\mathcal{R}_{\mathcal{M}}$ on $X^{\prime}$ is is the product $\prod_{\nu} \mathcal{M}\left(X^{\prime} \cap U^{\nu}\right)$. In particular, the space of global sections $\mathcal{R}_{\mathcal{M}}(X)$, which is $H^{0}\left(X, \mathcal{R}_{\mathcal{M}}\right)$, is the product $\prod_{\nu} \mathcal{M}\left(U^{\nu}\right)$.
(ii) The canonical map $\mathcal{M} \rightarrow \mathcal{R}_{\mathcal{M}}$ is injective. Thus, if $\mathcal{S}_{\mathcal{M}}$ denotes the cokernel of that map, there is a short exact sequence of $\mathcal{O}$-modules

$$
\begin{equation*}
0 \rightarrow \mathcal{M} \rightarrow \mathcal{R}_{\mathcal{M}} \rightarrow \mathcal{S}_{\mathcal{M}} \rightarrow 0 \tag{7.4.2}
\end{equation*}
$$

(iii) For any cohomology theory with the characteristic properties and for any $q>0, H^{q}\left(X, \mathcal{R}_{\mathcal{M}}\right)=0$.
proof. (i) This is seen by going through the definitions:

$$
\mathcal{R}\left(X^{\prime}\right)=\prod_{\nu}\left[j_{*}^{\nu} \mathcal{M}_{U^{\nu}}\right]\left(X^{\prime}\right)=\prod_{\nu} \mathcal{M}_{U^{\nu}}\left(X^{\prime} \cap U^{\nu}\right)=\prod_{\nu} \mathcal{M}\left(X^{\prime} \cap U^{\nu}\right)
$$

(ii) Let $X^{\prime}$ be an open subset of $X$. The map $\mathcal{M}\left(X^{\prime}\right) \rightarrow \mathcal{R}_{\mathcal{M}}\left(X^{\prime}\right)$ is the product of the restriction maps $\mathcal{M}\left(X^{\prime}\right) \rightarrow \mathcal{M}\left(X^{\prime} \cap U^{\nu}\right)$. Because the open sets $U^{\nu}$ cover $X$, the intersections $X^{\prime} \cap U^{\nu}$ cover $X^{\prime}$. The sheaf property of $\mathcal{M}$ tells us that the map $\mathcal{M}\left(X^{\prime}\right) \rightarrow \prod_{\nu} \mathcal{M}\left(X^{\prime} \cap U^{\nu}\right)$ is injective.
(iii) This follows from the third characteristic property.
7.4.3. Lemma. (i) A short exact sequence $0 \rightarrow \mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{P} \rightarrow 0$ of $\mathcal{O}$-modules embeds into a diagram

ttdiagr

Rsequence
uniquecoh hone

$$
0 \rightarrow H^{0}(X, \mathcal{M}) \rightarrow H^{0}\left(X, \mathcal{R}_{\mathcal{M}}\right) \rightarrow H^{0}\left(X, \mathcal{S}_{\mathcal{M}}\right) \xrightarrow{\delta^{0}} H^{1}(X, \mathcal{M}) \rightarrow 0
$$

and isomorphisms

$$
\begin{equation*}
0 \rightarrow H^{q}\left(X, \mathcal{S}_{\mathcal{M}}\right) \xrightarrow{\delta^{q}} H^{q+1}(X, \mathcal{M}) \rightarrow 0 \tag{7.4.8}
\end{equation*}
$$

for every $q>0$. The first three terms of the sequence (7.4.7), and the arrows connecting them, depend on our choice of covering of $X$, but the important point is that they don't depend on the cohomology. So that sequence determines $H^{1}(X, \mathcal{M})$ up to unique isomorphism as the cokernel of a map that is independent of the cohomology, and this is true for every $\mathcal{O}$-module $\mathcal{M}$, including for the module $\mathcal{S}_{\mathcal{M}}$. Therefore it is also true that $H^{1}\left(X, \mathcal{S}_{\mathcal{M}}\right)$ is determined uniquely. This being so, $H^{2}(X, \mathcal{M})$ is determined uniquely for every $\mathcal{M}$, by the isomorphism (7.4.8), with $q=1$. The isomorphisms (7.4.8) determine the rest of the cohomology up to unique isomorphism by induction on $q$.
existcoh

MtoRtoMone cmonesequence RMone RM

Rcmexact

Rsequencetwo

RMX

## (7.4.9) construction of cohomology

One can use the sequence $\sqrt{7.4 .2}$ and induction to construct cohomology, but it seems clearer to proceed by iterating the construction of $\mathcal{R}_{\mathcal{M}}$.

Let $\mathcal{M}$ be an $\mathcal{O}$-module. We rewrite the exact sequence 7.4.2, labeling $\mathcal{R}_{\mathcal{M}}$ as $\mathcal{R}_{\mathcal{M}}^{0}$, and $\mathcal{S}_{\mathcal{M}}$ as $\mathcal{M}^{1}$ :

$$
\begin{equation*}
0 \rightarrow \mathcal{M} \rightarrow \mathcal{R}_{\mathcal{M}}^{0} \rightarrow \mathcal{M}^{1} \rightarrow 0 \tag{7.4.10}
\end{equation*}
$$

and we repeat the construction with $\mathcal{M}^{1}$. Let $\mathcal{R}_{\mathcal{M}}^{1}=\mathcal{R}_{\mathcal{M}^{1}}^{0}\left(=j_{*} \mathcal{M}_{\mathbb{U}}^{1}\right)$, so that there is an exact sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{M}^{1} \rightarrow \mathcal{R}_{\mathcal{M}}^{1} \rightarrow \mathcal{M}^{2} \rightarrow 0 \tag{7.4.11}
\end{equation*}
$$

analogous to the sequence 7.4.10, with $\mathcal{M}^{2}=\mathcal{R}_{\mathcal{M}}^{1} / \mathcal{M}^{1}$. We combine the sequences 7.4.10 and 7.4.11, into an exact sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{M} \rightarrow \mathcal{R}_{\mathcal{M}}^{0} \rightarrow \mathcal{R}_{\mathcal{M}}^{1} \rightarrow \mathcal{M}^{2} \rightarrow 0 \tag{7.4.12}
\end{equation*}
$$

Then we let $\mathcal{R}_{\mathcal{M}}^{2}=\mathcal{R}_{\mathcal{M}^{2}}^{0}$. We continue in this way, to construct modules $\mathcal{R}_{\mathcal{M}}^{k}$ that form an exact sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{M} \rightarrow \mathcal{R}_{\mathcal{M}}^{0} \rightarrow \mathcal{R}_{\mathcal{M}}^{1} \rightarrow \mathcal{R}_{\mathcal{M}}^{2} \rightarrow \cdots \tag{7.4.13}
\end{equation*}
$$

The next lemma follows by induction from Lemmas 7.4.1 and 7.4.3.

### 7.4.14. Lemma.

(i) Let $0 \rightarrow \mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{P} \rightarrow 0$ be a short exact sequence of $\mathcal{O}$-modules. For every $n$, the sequences

$$
0 \rightarrow \mathcal{R}_{\mathcal{M}}^{n} \rightarrow \mathcal{R}_{\mathcal{N}}^{n} \rightarrow \mathcal{R}_{\mathcal{P}}^{n} \rightarrow 0
$$

are exact, and so are the sequences of global sections

$$
0 \rightarrow \mathcal{R}_{\mathcal{M}}^{n}(X) \rightarrow \mathcal{R}_{\mathcal{N}}^{n}(X) \rightarrow \mathcal{R}_{\mathcal{P}}^{n}(X) \rightarrow 0
$$

(ii) If $H^{0}, H^{1}$, .. is a cohomology theory, then $H^{q}\left(X, \mathcal{R}_{\mathcal{M}}^{n}\right)=0$ for all $n$ and all $q>0$.

An exact sequence such as 7.4 .13 is called a resolution of $\mathcal{M}$, and because $H^{q}\left(X, \mathcal{R}_{\mathcal{M}}^{n}\right)=0$ when $q>0$, it is called an acyclic resolution.

Continuing with the proof of existence, we consider the complex of $\mathcal{O}$-modules $\mathcal{R}_{\mathcal{M}}^{\boldsymbol{\mathcal { M }}}$ that is obtained by omitting the first term from 7.4.13:

$$
\begin{equation*}
0 \rightarrow \mathcal{R}_{\mathcal{M}}^{0} \rightarrow \mathcal{R}_{\mathcal{M}}^{1} \rightarrow \mathcal{R}_{\mathcal{M}}^{2} \rightarrow \cdots \tag{7.4.15}
\end{equation*}
$$

and the complex $\mathcal{R}_{\mathcal{M}}^{\bullet}(X)$ of its global sections:

$$
\begin{equation*}
0 \rightarrow \mathcal{R}_{\mathcal{M}}^{0}(X) \rightarrow \mathcal{R}_{\mathcal{M}}^{1}(X) \rightarrow \mathcal{R}_{\mathcal{M}}^{2}(X) \rightarrow \cdots \tag{7.4.16}
\end{equation*}
$$

which we can also write as

$$
0 \rightarrow H^{0}\left(X, \mathcal{R}_{\mathcal{M}}^{0}\right) \rightarrow H^{0}\left(X, \mathcal{R}_{\mathcal{M}}^{1}\right) \rightarrow H^{0}\left(X, \mathcal{R}_{\mathcal{M}}^{2}\right) \rightarrow \cdots
$$

The sequence $\mathcal{R}_{\mathcal{M}}^{\bullet}$ becomes the resolution (7.4.13) when the module $\mathcal{M}$ is inserted. So the complex (7.4.15) is exact except at $\mathcal{R}_{\mathcal{M}}^{0}$. But the global section functor is only left exact, and the sequence 7.4.16, of global sections $\mathcal{R}_{\mathcal{M}}^{\bullet}(X)$ needn't be exact anywhere. However, it is a complex because $\mathcal{R}_{\mathcal{M}}$ is a complex. The composition of adjacent maps is zero.

Recall that the cohomology of a complex $0 \rightarrow V^{0} \xrightarrow{d^{0}} V^{1} \xrightarrow{d^{1}} \cdots$ of vector spaces is defined to be $\mathbf{C}^{q}\left(V^{\bullet}\right)=\left(\operatorname{ker} d^{q}\right) /\left(\operatorname{im} d^{q-1}\right)$, and that $\left\{\mathbf{C}^{q}\right\}$ is a cohomological functor on complexes 7.2.5.
7.4.17. Definition. The cohomology of an $\mathcal{O}$-module $\mathcal{M}$ is the cohomology of the complex $\mathcal{R}_{\mathcal{M}}^{\bullet}(X)$ :

$$
H^{q}(X, \mathcal{M})=\mathbf{C}^{q}\left(\mathcal{R}_{\mathcal{M}}^{\bullet}(X)\right)
$$

Thus if we denote the maps in the complex (7.4.16) by $d^{q}$,

$$
0 \rightarrow \mathcal{R}_{\mathcal{M}}^{0}(X) \xrightarrow{d^{0}} \mathcal{R}_{\mathcal{M}}^{1}(X) \xrightarrow{d^{1}} \mathcal{R}_{\mathcal{M}}^{2}(X) \rightarrow \cdots
$$

then $H^{q}(X, \mathcal{M})=\left(\operatorname{ker} d^{q}\right) /\left(\operatorname{im} d^{q-1}\right)$.
7.4.18. Lemma. Let $X$ be an affine variety. With cohomology defined as above, $H^{q}(X, \mathcal{M})=0$ for all $\mathcal{O}$-modules $\mathcal{M}$ and all $q>0$.
proof. When $X$ is affine, the sequence of global sections of the exact sequence 7.4 .13 is exact.
To show that our definition gives the (unique) cohomology, we verify the three characteristic properties. Since the sequence 7.4 .13 is exact and since the global section functor is left exact, $\mathcal{M}(X)$ is the kernel of the map $\mathcal{R}_{\mathcal{M}}^{0}(X) \rightarrow \mathcal{R}_{\mathcal{M}}^{1}(X)$, and this kernel is also equal to $\mathbf{C}^{0}\left(\mathcal{R}_{\mathcal{M}}^{\bullet}(X)\right)$. So our cohomology has the first property: $H^{0}(X, \mathcal{M})=\mathcal{M}(X)$.

To show that we obtain a cohomological functor, we apply Lemma 7.4.14 to conclude that, for a short exact sequence $0 \rightarrow \mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{P} \rightarrow 0$, the global sections

$$
\begin{equation*}
0 \rightarrow \mathcal{R}_{\mathcal{M}}^{\bullet}(X) \rightarrow \mathcal{R}_{\mathcal{N}}^{\bullet}(X) \rightarrow \mathcal{R}_{\mathcal{P}}^{\bullet}(X) \rightarrow 0 \tag{7.4.19}
\end{equation*}
$$

form an exact sequence of complexes. The cohomology $H^{q}(X, \cdot)$ is a cohomological functor because cohomology of complexes is a cohomological functor.

We make a digression before verifying the third characteristic property.

## (7.4.20) affine morphisms

Let $Y \xrightarrow{f} X$ be a morphism of varieties. Let $U \xrightarrow{j} X$ be the inclusion of an open subvariety into $X$ and let $V$ be the inverse image $f^{-1} U$, which is an open subvariety of $Y$. These varieties and maps form a diagram

affinemorph
jstarfstar
jstarfstartwo
proof. Let $U^{\prime}$ be an open subset of $U$, and let $V^{\prime}=g^{-1} U^{\prime}$. Then

$$
\left[f_{*} \mathcal{N}\right]_{U}\left(U^{\prime}\right)=\left[f_{*} \mathcal{N}\right]\left(U^{\prime}\right)=\mathcal{N}\left(V^{\prime}\right)=\mathcal{N}_{V}\left(V^{\prime}\right)=\left[g_{*}\left[\mathcal{N}_{V}\right]\right]\left(U^{\prime}\right)
$$

defaffmorph
7.4.23. Definition. An affine morphism is a morphism $Y \xrightarrow{f} X$ of varieties with the property that the inverse image $f^{-1}(U)$ of every affine open subset $U$ of $X$ is an affine open subset of $Y$.

The following are examples of affine morphisms:

- the inclusion of an affine open subset $Y$ into $X$,
- the inclusion of a closed subvariety $Y$ into $X$,
- a finite morphism, or an integral morphism.

But, the inclusion of a nonaffine open set may not be an affine morphism, and if $Y$ is a closed subset of $\mathbb{P}^{n} \times X$, the projection $Y \rightarrow X$ will not be an affine morphism unless its fibres are finite, in which case Chevalley's Finitenss Theorem tells us that it is a finite morphism.
7.4.24. Lemma. Lert $Y \xrightarrow{f} C$ be an affine morphism and let $\mathcal{N} \rightarrow \mathcal{N}^{\prime} \rightarrow \mathcal{N}^{\prime \prime}$ be an exact sequence of $\mathcal{O}_{Y \text {-modules }}$ The sequence of direct images $f_{*} \mathcal{N} \rightarrow f_{*} \mathcal{N}^{\prime} \rightarrow f_{*} \mathcal{N}^{\prime \prime}$ is exact.

Let $Y \xrightarrow{f} X$ be an affine morphism, let $j$ be the map from our chosen affine covering $\mathbb{U}=\left\{U^{\nu}\right\}$ to $X$, and let $\mathbb{V}$ denote the family $\left\{V^{\nu}\right\}=\left\{f^{-1} U^{\nu}\right\}$ of inverse images. Then $\mathbb{V}$ is an affine covering of $Y$, and there is a morphism $\mathbb{V} \xrightarrow{g} \mathbb{U}$. We form a diagram analogous to 7.4.21, in which $\mathbb{V}$ and $\mathbb{U}$ replace $V$ and $U$, respectively:

7.4.25. Proposition. Let $Y \xrightarrow{f} X$ be an affine morphism, and let $\mathcal{N}$ be an $\mathcal{O}_{Y}$-module. Let $H^{q}(X, \cdot)$ be cohomology defined in (7.4.17), and let $H^{q}(Y, \cdot)$ be cohomology defined in the analogous way, using the covering $\mathbb{V}$ of $Y$. Then $H^{q}\left(X, f_{*} \mathcal{N}\right)$ is isomorphic to $H^{q}(Y, \mathcal{N})$.
proof. This proof is easy, except that one has to untangle the notation.
To compute the cohomology of $f_{*} \mathcal{N}$ on $X$, we substitute $\mathcal{M}=f_{*} \mathcal{N}$ into 7.4.17):

$$
H^{q}\left(X, f_{*} \mathcal{N}\right)=\mathbf{C}^{q}\left(\mathcal{R}_{f_{*} \mathcal{N}}(X)\right)
$$

To compute the cohomology of $\mathcal{N}$ on $Y$, we let

$$
\mathcal{R}_{\mathcal{N}}^{\prime 0}=i_{*}\left[\mathcal{N}_{V}\right]
$$

and we continue, to construct a resolution $\mathcal{R}^{\prime \bullet}=0 \rightarrow \mathcal{N} \rightarrow \mathcal{R}_{\mathcal{N}}^{\prime 0} \rightarrow \mathcal{R}^{\prime \prime}{ }_{\mathcal{N}} \rightarrow \cdots$ and the complex of its global sections $\mathcal{R}_{\mathcal{N}}^{\prime \bullet}(Y)$. The prime is there to remind us that $\mathcal{R}^{\prime}$ is defined using the covering $\mathbb{V}$ of $Y$. Then

$$
H^{q}(Y, \mathcal{N})=\mathbf{C}^{q}\left(\mathcal{R}_{\mathcal{N}}^{\prime \bullet}(Y)\right)
$$

It suffices to show that the complexes $\mathcal{R}_{f_{* \mathcal{N}}}^{\bullet}(X)$ and $\mathcal{R}_{\mathcal{N}}^{\prime \bullet}(Y)$ are isomorphic. If so, we will have

$$
H^{q}\left(X, f_{*} \mathcal{N}\right)=\mathbf{C}^{q}\left(\mathcal{R}_{f_{*} \mathcal{N}}^{\bullet}(X)\right) \approx \mathbf{C}^{q}\left(\mathcal{R}_{\mathcal{N}}^{\prime \bullet}(Y)\right)=H^{q}(Y, \mathcal{N})
$$

as required.
By definition of the direct image, $\left[f_{*} \mathcal{R}_{\mathcal{N}}^{\prime q}\right](X)=\mathcal{R}_{\mathcal{N}}^{\prime q}(Y)$. So it suffices to show that $\mathcal{R}_{f_{*} \mathcal{N}}^{q} \approx f_{*} \mathcal{R}_{\mathcal{N}}^{\prime q}$. We look back at the definition 7.4.11 of the modules $\mathcal{R}^{0}$ in its rewritten form 7.4.10. On $Y$, the analogous sequence for $\mathcal{N}$ analogous to

$$
0 \rightarrow \mathcal{N} \rightarrow{\mathcal{R}_{\mathcal{N}}^{\prime}}^{0} \rightarrow \mathcal{N}^{1} \rightarrow 0
$$

where $\mathcal{R}^{\prime}{ }_{\mathcal{N}}=i_{*} \mathcal{N}_{\mathbb{V}}$. When $f$ is an affine orphism, the direct image of this sequence is exact:

$$
0 \rightarrow f_{*} \mathcal{N} \rightarrow f_{*} \mathcal{R}_{\mathcal{N}}^{\prime 0} \rightarrow f_{*} \mathcal{N}^{1} \rightarrow 0
$$

Here

So $f_{*} \mathcal{R}^{\prime 0}{ }_{\mathcal{N}}=\mathcal{R}_{f_{c} n}^{0}$. Now induction on $q$ applies.
We go back to the proof of existence of cohomology to verify the third characteristic property, which is that when $Y \xrightarrow{f} X$ is the inclusion of an affine open subset, $H^{q}\left(X, f_{*} \mathcal{N}\right)=0$ for all $\mathcal{O}_{Y}$-modules $\mathcal{N}$ and all $q>0$. The inclusion of an affine open set is an affine morphism, so $H^{q}(Y, \mathcal{N})=H^{q}\left(X, f_{*} \mathcal{N}\right)$ 7.4.25, and since $Y$ is affine, $H^{q}(Y, \mathcal{N})=0$ for all $q>0$ 7.4.18.

Proposition 7.4.25 is one of the places where a specific construction of cohomology is used. The characteristic properties don't apply directly. The next proposition is another such place.
7.4.26. Lemma. Cohomology is compatible with limits of directed sets of $\mathcal{O}$-modules: $H^{q}\left(X, \underset{\longrightarrow}{\lim } \mathcal{M}_{\bullet}\right) \approx$ $\xrightarrow{\lim } H^{q}\left(X, \mathcal{M}_{\bullet}\right)$ for all $q$.
proof. The direct and inverse image functors and the global section functor are all compatible with $\underset{\rightarrow}{\lim }$, and $\underset{\longrightarrow}{\lim }$ is exact (??). So the module $\mathcal{R}_{\underset{l}{q}}^{q} \mathcal{M}_{\bullet}$ that is used to compute the cohomology of $\underset{\longrightarrow}{\lim } \mathcal{M}_{\bullet}$ is isomorphic to $\underset{\longrightarrow}{\lim }\left[\mathcal{R}_{\mathcal{M}_{\bullet}}^{q}\right]$ and $\mathcal{R}_{\underline{\lim } \mathcal{M}_{\bullet}}^{q}(X)$ is isomorphic to $\xrightarrow[\longrightarrow]{\lim }\left[\mathcal{R}_{\mathcal{M}_{\bullet}}^{q}\right](X)$.

## (7.4.27) uniqueness of the coboundary maps

We have constructed a cohomology $\left\{H^{q}\right\}$ that has the characteristic properties, and we have shown that the functors $H^{q}$ are unique. We haven't shown that the coboundary maps $\delta^{q}$ that appear in the cohomology sequences are unique. To make it clear that there is something to show, we note that the cohomology sequence 7.1.3 remains exact when some of the coboundary maps $\delta^{q}$ are multiplied by -1 . Why can't we define a new collection of coboundary maps by changing some signs? The reason we can't do this is that we used the coboundary maps $\delta^{q}$ in 7.4.7) and 7.4.8, to identify $H^{q}(X, \mathcal{M})$. Having done that, we aren't allowed to change $\delta^{q}$ for the particular short exact sequences 7.4 .2 . We show that the coboundary maps for those sequences determine the coboundary maps for every short exact sequence of $\mathcal{O}$-modules

$$
\begin{equation*}
0 \rightarrow \mathcal{M} \longrightarrow \mathcal{N} \longrightarrow \mathcal{P} \rightarrow 0 \tag{A}
\end{equation*}
$$

The sequences $\overline{7.4 .2}$ were rewritten as 7.4 .10 . We will use that form.
To show that the coboundaries for the sequence $(A)$ are determined uniquely, we relate it to a sequence for which the coboundary maps are fixed:

$$
\begin{equation*}
0 \rightarrow \mathcal{M} \longrightarrow \mathcal{R}_{\mathcal{M}}^{0} \longrightarrow \mathcal{M}^{1} \rightarrow 0 \tag{B}
\end{equation*}
$$

We map $(A)$ and $(B)$ to a third short exact sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{M} \xrightarrow{\psi} \mathcal{R}_{\mathcal{N}}^{0} \longrightarrow \mathcal{Q} \rightarrow 0 \tag{C}
\end{equation*}
$$

where $\psi$ is the composition of the injective maps $\mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{R}_{\mathcal{N}}^{0}$ and $\mathcal{Q}$ is the cokernel of $\psi$.
First, we inspect the diagram

and its diagram of coboundary maps


This diagram shows that the coboundary map $\delta_{A}^{q}$ for the sequence $(A)$ is determined by the coboundary map $\delta_{C}^{q}$ for $(C)$.

Next, we inspect the diagram

BtoC

and its diagram of coboundary maps


When $q>0, \delta_{C}^{q}$ and $\delta_{B}^{q}$ are bijective because the cohomology of $\mathcal{R}_{\mathcal{M}}^{0}$ and $\mathcal{R}_{\mathcal{N}}^{0}$ is zero in positive dimension. Then $\delta_{C}^{q}$ is uniquely determined by $\delta_{B}^{q}$, and so is $\delta_{A}^{q}$.

We have to look more closely to settle the case $q=0$. The map labeled $u$ in 7.4.28 is injective, and the Snake Lemma shows that $v$ is injective. The cokernels of $u$ and $v$ are isomorphic. We write both of the cokernels as $\mathcal{R}_{\mathcal{P}}^{0}$. When we add the cokernels to the diagram, and pass to cohomology, we obtain a diagram whose relevant part is


The rows and columns in the diagram are exact. We want to show that the map $\delta_{C}^{0}$ is determined uniquely by $\delta_{B}^{0}$. It is determined by $\delta_{B}^{0}$ on the image of $v$ and it is zero on the image of $\beta$. To show that $\delta_{C}^{0}$ is determined by $\delta_{B}^{0}$, it suffices to show that the images of $v$ and $\beta$ together span $H^{0}(X, \mathcal{Q})$. This follows from the fact that $\gamma$ is surjective. Thus $\delta_{C}^{0}$ is determined uniqely by $\delta_{B}^{0}$, and so is $\delta_{A}^{0}$.

### 7.5 Cohomology of the Twisting Modules

We determine the cohomology of the twisting modules $\mathcal{O}(d)$ on $\mathbb{P}^{n}$ here. As we will see, $H^{q}\left(\mathbb{P}^{n}, \mathcal{O}(d)\right)$ is zero for most values of $q$. This will help us to determine the cohomology of other modules.

Lemma 7.4.18 about vanishing of cohomology on an affine variety, and Lemma 7.4.25 about the direct image via an affine morphism, were stated using a particular affine covering. Since we know that cohomology is unique, that particular covering is irrelevant. Though it isn't necessary, we restate those lemmas here as a corollary:
7.5.1. Corollary. (i) On an affine variety $X, H^{q}(X, \mathcal{M})=0$ for all $\mathcal{O}$-modules $\mathcal{M}$ and all $q>0$.
 isomorphic. If $Y$ is an affine variety, $H^{q}\left(X, f_{*} \mathcal{N}\right)=0$ for all $q>0$.

One case in which (ii) applies is that $f$ is the inclusion of a closed subvariety $Y$ into $X$.
Let $\mathcal{M}$ be a finite $\mathcal{O}$-module on projective space $\mathbb{P}^{n}$. The twisting modules $\mathcal{O}(d)$ and the twists $\mathcal{M}(d)=$ $\mathcal{M} \otimes_{\mathcal{O}} \mathcal{O}(d)$ are isomorphic to the modules $\mathcal{O}(d H)$ and $\mathcal{M}(d H)$, respectively 6.7.11. They form maps of directed sets

(See (??)). The second diagram is obtained from the first one by tensoring with $\mathcal{M}$. Let $\mathbb{U}$ denote the standard affine open subset $\mathbb{U}^{0}$ of $\mathbb{P}^{n}$, and let $j$ be the inclusion of $\mathbb{U}$ into $\mathbb{P}^{n}$. Then $\underline{\lim } \mathcal{O}(d H) \approx j_{*} \mathcal{O}_{\mathbb{U}}(?$ ? $)$ and because $\xrightarrow{\lim }$ is compatible with tensor products, $\xrightarrow{\lim } \mathcal{M}(d H) \approx j_{*} \mathcal{M}_{\mathbb{U}}$ 6.7.16. Since $j$ is an affine morphism and $\mathbb{U}^{0}$ is an affine open set, $H^{q}\left(\mathbb{P}^{n}, j_{*} \mathcal{O}_{\mathbb{U}}\right)=0$ and $H^{q}\left(\mathbb{P}^{n}, j_{*} \mathcal{M}_{\mathbb{U}}\right)=0$ for all $q>0$.

The next corollary follows from the facts that $\mathcal{M}(d)$ is isomorphic to $\mathcal{M}(d H)$, and that cohomology is compatible with direct limits 7.4.26.
7.5.2. Corollary. For all projective varieties $X$, all $\mathcal{O}$-modules $\mathcal{M}$, and all $q>0 \lim _{\longrightarrow} H^{q}(X, \mathcal{O}(d))=0$ and $\underset{\longrightarrow}{\lim } H^{q}(X, \mathcal{M}(d))=0$.
7.5.3. Notation. If $\mathcal{M}$ is an $\mathcal{O}$-module, we denote the dimension of $H^{q}(X, \mathcal{M})$ by $\mathbf{h}^{q}(X, \mathcal{M})$ or by $\mathbf{h}^{q} \mathcal{M}$. We can write $\mathbf{h}^{q} \mathcal{M}=\infty$ if the dimension is infinite. However, in Section 7.7 we will see that when $\mathcal{M}$ is a finite $\mathcal{O}$-module on a projective variety $X, H^{q}(X, \mathcal{M})$ has finite dimension for every $q$.

### 7.5.4. Theorem.

(i) For $d \geq 0, \quad \mathbf{h}^{0}\left(\mathbb{P}^{n}, \mathcal{O}(d)\right)=\binom{d+n}{n}$ and $\mathbf{h}^{q}\left(\mathbb{P}^{n}, \mathcal{O}(d)\right)=0$ if $q \neq 0$.
(ii) For $r>0, \quad \mathbf{h}^{n}\left(\mathbb{P}^{n}, \mathcal{O}(-r)\right)=\binom{r-1}{n}$ and $\mathbf{h}^{q}\left(\mathbb{P}^{n}, \mathcal{O}(-r)\right)=0$ if $q \neq n$.

In particular, the case $r=1$ in part (ii) asserts that $\mathbf{h}^{q}\left(\mathbb{P}^{n}, \mathcal{O}(-1)\right)=0$ for all $q$.
proof. We have described the global sections of $\mathcal{O}(d)$ before: If $d \geq 0, H^{0}(X, \mathcal{O}(d))$ is the space of homogeneous polynomials of degree $d$ in the coordinate variables, and if $d<0, H^{0}(X, \mathcal{O}(d))=0$ (see 6.7.2).
(i) (the case $d \geq 0$ )

Let $X=\mathbb{P}^{n}$, and let $Y \xrightarrow{i} X$ be the inclusion of the hyperplane at infinity $x_{0}=0$ into $X$. By induction on $n$, we may assume that the theorem has been proved for $Y$, which is a projective space of dimension $n-1$. We consider the exact sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{O}_{X}(-1) \xrightarrow{x_{0}} \mathcal{O}_{X} \rightarrow i_{*} \mathcal{O}_{Y} \rightarrow 0 \tag{7.5.5}
\end{equation*}
$$

and its twists

$$
\begin{equation*}
0 \rightarrow \mathcal{O}_{X}(d-1) \xrightarrow{x_{0}} \mathcal{O}_{X}(d) \rightarrow i_{*} \mathcal{O}_{Y}(d) \rightarrow 0 \tag{7.5.6}
\end{equation*}
$$

The twisted sequences are exact because they are obtained by tensoring 7.5 .5 with the invertible $\mathcal{O}$-modules $\mathcal{O}(d)$. Because the inclusion $i$ is an affine morphism, $H^{q}\left(X, i_{*} \mathcal{O}_{Y}(d)\right) \approx H^{q}\left(Y, \mathcal{O}_{Y}(d)\right)$.

The monomials of degree $d$ in $n+1$ variables form a basis of the space of global sections of $\mathcal{O}_{X}(d)$. Setting $x_{0}=0$ and deleting terms that become zero gives us a basis of $\mathcal{O}_{Y}(d)$. Every global section of $\mathcal{O}_{Y}(d)$ is the restriction of a global section of $\mathcal{O}_{X}(d)$. The sequence of global sections

$$
0 \rightarrow H^{0}\left(X, \mathcal{O}_{X}(d-1)\right) \xrightarrow{x_{0}} H^{0}\left(X, \mathcal{O}_{X}(d)\right) \rightarrow H^{0}\left(Y, \mathcal{O}_{Y}(d)\right) \rightarrow 0
$$

is exact. This tells us that the map $H^{1}\left(X, \mathcal{O}_{X}(d-1)\right) \longrightarrow H^{1}\left(X, \mathcal{O}_{X}(d)\right)$ is injective.
By induction on the dimension $n, H^{q}\left(Y, \mathcal{O}_{Y}(d)\right)=0$ for $d \geq 0$ and $q>0$. When combined with the injectivity noted above, the cohomology sequence of 7.5 .6 shows that the maps $H^{q}\left(X, \mathcal{O}_{X}(d-1)\right) \rightarrow$ $H^{q}\left(X, \mathcal{O}_{X}(d)\right)$ are bijective for every $q>0$. Since the limits are zero 7.5.2, $H^{q}\left(X, \mathcal{O}_{X}(d)\right)=0$ for all $d \geq 0$ and all $q>0$.
(ii) (the case $d<0$, or $r>0$.)

We use induction on the integers $r$ and $n$. We suppose the theorem proved for $r$, and we substitute $d=-r$ into the sequence (7.5.6:

The base case $r=0$ is the exact sequence 7.5.5). In the cohomology sequence associated to that sequence, the terms $H^{q}\left(X, \mathcal{O}_{X}\right)$ and $H^{q}\left(Y, \mathcal{O}_{Y}\right)$ are zero when $q>0$, and $H^{0}\left(X, \mathcal{O}_{X}\right)=H^{0}\left(Y, \mathcal{O}_{Y}\right)=\mathbb{C}$. Therefore $H^{q}\left(X, \mathcal{O}_{X}(-1)\right)=0$ for every $q$. This proves (ii) for $r=1$.

Our induction hypothesis is that, $\mathbf{h}^{n}(X, \mathcal{O}(-r))=\binom{r-1}{n}$ and $\mathbf{h}^{q}=0$ if $q \neq n$. By induction on $n$, we may suppose that $\mathbf{h}^{n-1}(Y, \mathcal{O}(-r))=\binom{r-1}{n-1}$ and $\mathbf{h}^{q}(Y, \mathcal{O}(-r))=0$ if $q \neq n-1$. Instead of displaying the cohomology sequence associated to 7.5.7, we assemble the dimensions of cohomology into a table in which the asterisks stand for entries that are to be determined:

|  | $\mathcal{O}_{X}(-(r+1))$ | $\mathcal{O}_{X}(-r)$ | $i_{*} \mathcal{O}_{Y}(-r)$ |
| :---: | :---: | :---: | :---: |
| $\mathbf{h}^{0} \quad \vdots$ | $*$ | 0 | 0 |
| $\vdots$ | $\vdots$ | $\vdots$ | $\vdots$ |
| $\mathbf{h}^{n-2}:$ | $*$ | 0 | 0 |
| $\mathbf{h}^{n-1}:$ | $*$ | 0 | $\binom{r-1}{n-1}$ |
| $\mathbf{h}^{n}$ | $:$ | $*$ | $\binom{r-1}{n}$ |

The second column is determined by induction on $r$ and the third by induction on $n$. The cohomology sequence shows that that

$$
\mathbf{h}^{n}(X, \mathcal{O}(-(r+1)))=\binom{r-1}{n-1}+\binom{r-1}{n}
$$

and that the other entries labeled with an asterisk are zero. The right side of this equation is equal to $\binom{r}{n}$.

### 7.6 Cohomology of Hypersurfaces

We determine the cohomology of a plane projective curve first. Let $X$ be the projective plane $\mathbb{P}^{2}$ and let $C \xrightarrow{i} X$ denote the inclusion of a plane curve $C$ of degree $k$. The ideal $\mathcal{I}$ of functions that vanish on $C$ is isomorphic to the twisting module $\mathcal{O}_{X}(-k)$ 6.7.8. So one has an exact sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{O}_{X}(-k) \rightarrow \mathcal{O}_{X} \rightarrow i_{*} \mathcal{O}_{C} \rightarrow 0 \tag{7.6.1}
\end{equation*}
$$

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We form a table showing dimensions of the cohomology. Theorem7.5.4 determines the first two columns, and the cohomology sequence determines the last column.

$$
\begin{array}{lccc} 
& \mathcal{O}_{X}(-k) & \mathcal{O}_{X} & i_{*} \mathcal{O}_{C} \\
\hline \mathbf{h}^{0}: & 0 & 1 & 1  \tag{7.6.2}\\
\mathbf{h}^{1}: & 0 & 0 & \binom{k-1}{2} \\
\mathbf{h}^{2}: & \binom{k-1}{2} & 0 & 0
\end{array}
$$

Since the inclusion of the curve $C$ into the projective plane $X$ is an affine morphism, $\mathbf{h}^{q}\left(X, i_{*} \mathcal{O}_{C}\right)=$ $\mathbf{h}^{q}\left(C, \mathcal{O}_{C}\right)$. Therefore

$$
\mathbf{h}^{0}\left(C, \mathcal{O}_{C}\right)=1, \mathbf{h}^{1}\left(C, \mathcal{O}_{C}\right)=\binom{k-1}{2}, \text { and } \mathbf{h}^{q}=0 \text { when } q>1
$$

The dimension $\mathbf{h}^{1}\left(C, \mathcal{O}_{C}\right)$, which is $\binom{k-1}{2}$, is called the arithmetic genus of $C$. It is denoted by $p_{a}\left(=p_{a}(C)\right)$. We will see later (8.8.2 that when $C$ is a smooth curve, its arithmetic genus is equal to its topological genus: $p_{a}=g$. But the arithmetic genus of a plane curve of degree $k$ is $\binom{k-1}{2}$ also when the curve $C$ is singular.

We restate the results as a corollary.
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7.6.3. Corollary. Let $C$ be a plane curve of degree $k$. Then $\mathbf{h}^{0} \mathcal{O}_{C}=1, \mathbf{h}^{1} \mathcal{O}_{C}=\binom{k-1}{2}=p_{a}$, and $\mathbf{h}^{q}=0$ if $q \neq 0,1$.

The fact that $\mathbf{h}^{0} \mathcal{O}_{C}=1$ tells us that the cibstabts are the only rational functions that are regular everywhere on $C$. This reflects a fact that will be proved later: A plane curve is compact and connected in the classical topology. However, it isn't a proof of that fact.

We will need more technique in order to compute cohomology of a curve that is embedded in a higher dimensional projective space. Cohomology of projective curves is the topic of Chapter 8 In the next section we will see that the cohomology on any projective curve is zero except in dimensions 0 and 1 .

One can make a similar computation for the hypersurface $Y$ in $X=\mathbb{P}^{n}$ defined by an irreducible homogeneous polynomial $f$ of degree $k$. The ideal of $Y$ is isomorphic to $\mathcal{O}_{X}(-k)$, and there is an exact sequence

$$
0 \rightarrow \mathcal{O}_{X}(-k) \xrightarrow{f} \mathcal{O}_{X} \rightarrow i_{*} \mathcal{O}_{Y} \rightarrow 0
$$

Since we know the cohomology of $\mathcal{O}_{X}(-k)$ and of $\mathcal{O}_{X}$, and since $H^{q}\left(X, i_{*} \mathcal{O}_{Y}\right) \approx H^{q}\left(Y, \mathcal{O}_{Y}\right)$, we can use this sequence to compute the dimensions of the cohomology of $\mathcal{O}_{Y}$.
7.6.4. Corollary. Let $Y$ be a hypersurface of dimension $d$ and degree $k$ in a projective space of dimension $d+1$. Then $\mathbf{h}^{0}\left(Y, \mathcal{O}_{Y}\right)=1, \mathbf{h}^{d}\left(Y, \mathcal{O}_{Y}\right)=\binom{k-1}{d+1}$, and $\mathbf{h}^{q}\left(Y, \mathcal{O}_{Y}\right)=0$ for all other $q$.

In particular, when $S$ is a surface in $\mathbb{P}^{3}$ defined by an irreducible polynomial of degree $k, \mathbf{h}^{0}\left(S, \mathcal{O}_{S}\right)=1$, $\mathbf{h}^{1}\left(S, \mathcal{O}_{S}\right)=0, \mathbf{h}^{2}\left(S, \mathcal{O}_{S}\right)=\binom{k-1}{3}$, and $\mathbf{h}^{q}=0$ if $q>2$. For a projective surface $S$ that isn't embedded into $\mathbb{P}^{3}$, it is still true that $\mathbf{h}^{q}=0$ when $q>2$, but $\mathbf{h}^{1}\left(S, \mathcal{O}_{S}\right)$ may be nonzero. The dimensions $\mathbf{h}^{1}\left(S, \mathcal{O}_{S}\right)$ and $\mathbf{h}^{2}\left(S, \mathcal{O}_{S}\right)$ are invariants of the surface somewhat analogous to the genus of a curve. In classical terminology, $\mathbf{h}^{2}\left(S, \mathcal{O}_{S}\right)$ is the geometric genus $p_{g}$ and $\mathbf{h}^{1}\left(S, \mathcal{O}_{S}\right)$ is the irregularity q . The arithmetic genus $p_{a}$ is defined to be

$$
\begin{equation*}
p_{a}=\mathbf{h}^{2}\left(S, \mathcal{O}_{S}\right)-\mathbf{h}^{1}\left(S, \mathcal{O}_{S}\right)=p_{g}-q \tag{7.6.5}
\end{equation*}
$$

Therefore the irregularity is $q=p_{g}-p_{a}$. When $S$ is a surface in $\mathbb{P}^{3}, q=0$ and $p_{g}=p_{a}$.
In modern terminology, it would be more natural to replace the arithmetic genus by the Euler characteristic of the structure sheaf $\chi\left(\mathcal{O}_{S}\right)$, which is defined to be $\sum_{q}(-1)^{q} \mathbf{h}^{q}\left(\mathcal{O}_{S}\right)$ (see 7.7.7 below. The Euler characteristic of the structure sheaf on a curve is

$$
\chi\left(\mathcal{O}_{C}\right)=\mathbf{h}^{0}\left(C, \mathcal{O}_{C}\right)-\mathbf{h}^{1}\left(C, \mathcal{O}_{C}\right)=1-p_{a}
$$

and on a surface $S$ it is

$$
\chi\left(\mathcal{O}_{S}\right)=\mathbf{h}^{0}\left(S, \mathcal{O}_{S}\right)-\mathbf{h}^{1}\left(S, \mathcal{O}_{S}\right)+\mathbf{h}^{2}\left(S, \mathcal{O}_{S}\right)=1+p_{a}
$$

But because of tradition, the arithmetic genus is still used quite often.

### 7.7 Three Theorems about Cohomology

These theorems are taken from Serre's 1956 paper.
7.7.1. Theorem. Let $X$ be a projective variety, and let $\mathcal{M}$ be a finite $\mathcal{O}_{X}$-module. If the support of $\mathcal{M}$ has dimension $k$, then $H^{q}(X, \mathcal{M})=0$ for all $q>k$. In particular, if $X$ has dimension $n$, then $H^{q}(X, \mathcal{M})=0$ for all $q>n$.

See Sectdon 6.6 for the definition of support.
7.7.2. Theorem. Let $\mathcal{M}(d)$ be the twist of a finite $\mathcal{O}_{X}$-module $\mathcal{M}$ on a projective variety $X$. For sufficiently large $d, H^{q}(X, \mathcal{M}(d))=0$ for all $q>0$.
7.7.3. Theorem. Let $\mathcal{M}$ be a finite $\mathcal{O}$-module on a projective variety $X$. The cohomology $H^{q}(X, \mathcal{M})$ is a finite-dimensional vector space for every $q$.
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7.7.4. Notes. (a) As the first theorem asserts, the highest dimension in which cohomology of an $\mathcal{O}_{X}$-module on a variety $X$ can be nonzero is the dimension of $X$. It is also true that, on any projective variety $X$ of dimension $n$, there will be $\mathcal{O}$-modules $\mathcal{M}$ such that $H^{n}(X, \mathcal{M}) \neq 0$. In contrast, in the classical topology on a projective variety $X$ of dimnsion $n$, the constant coefficient cohomology $H^{2 n}\left(X_{\text {class }}, \mathbb{Z}\right)$ isn't zero. As we have mentioned, the constant coefficient cohomology is zero for every $q>0$ in the Zariski topology. And when $X$ is affine, the cohomology of any $\mathcal{O}$-module is for all $q>0$.
(b) The third theorem tells us that the space of global sections of a finite $\mathcal{O}$-module $\mathcal{M}$ is finite-dimensional. This is one of the most important consequences of the theorem, and it isn't easy to prove directly. Cohomology needn't be finite-dimensional when a variety isn't projective. For example, on an affine variety $X=\operatorname{Spec} A$, $H^{0}(X, \mathcal{O})=A$ isn't finite-dimensional unless $X$ is a point. When $X$ is the complement of a point in $\mathbb{P}^{2}$, $H^{1}(X, \mathcal{O})$ isn't finite-dimensional.
(c) The structure of the proofs is interesting. The first theorem allows us to use descending induction to prove the second and third theorems, beginning with the fact that $\mathrm{s} H^{k}(X, \mathcal{M})=0$ when $k$ is greater than the dimension of $X$.

In these theorems, we are given that $X$ is a closed subvariety of a projective space $\mathbb{P}^{n}$. We can replace an $\mathcal{O}_{X}$-module by its extension by zero to $\mathbb{P}^{n}$ 7.5.1. This doesn't change the cohomology or the dimension of support. So we may assume that $X$ is a projective space. In fact, the twist $\mathcal{M}(d)$ of an $\mathcal{O}_{X}$-module that is referred to in the second theorem is defined in terms of the extension by zero (??).

The proofs are based on the cohomology of the twisting modules 7.5.4 and the vanishing of the limit $\xrightarrow{l i m} H^{q}(X, \mathcal{M}(d))$ for $q>0$ (7.5.2).
proof of Theorem 7.7.1 (vanishing in large dimension)
Here $\mathcal{M}$ is a finite $\mathcal{O}$-module whose support $S$ has dimension $k$ or less. We are to show that $H^{q}(X, \mathcal{M})=0$ when $q>k$. We choose coordinates so that the hyperplane $H: x_{0}=0$ doesn't contain any component of the support $S$. Then $H \cap S$ has dimension at most $k-1$. We inspect the multiplication map $\mathcal{M}(-1) \xrightarrow{x_{0}} \mathcal{M}$. The kernel $\mathcal{K}$ and cokernel $\mathcal{Q}$ are annihilated by $x_{0}$, so the supports of $\mathcal{K}$ and $\mathcal{Q}$ are contained in $H$. Since they are also in $S$, the supports have dimension at most $k-1$. We can apply induction on $k$ to them. In the base case $k=0$, the supports of $\mathcal{K}$ and $\mathcal{Q}$ will be empty, and the cohomology will be zero.

We break the exact sequence $0 \rightarrow \mathcal{K} \rightarrow \mathcal{M}(-1) \rightarrow \mathcal{M} \rightarrow \mathcal{Q} \rightarrow 0$ into two short exact sequences by introducing the $\mathcal{O}$-module $\mathcal{N}=\operatorname{ker}(\mathcal{M} \rightarrow \mathcal{Q})$ :

$$
\begin{equation*}
0 \rightarrow \mathcal{K} \rightarrow \mathcal{M}(-1) \rightarrow \mathcal{N} \rightarrow 0 \quad \text { and } \quad 0 \rightarrow \mathcal{N} \rightarrow \mathcal{M} \rightarrow \mathcal{Q} \rightarrow 0 \tag{7.7.5}
\end{equation*}
$$

The induction hypothesis applies to $\mathcal{K}$ and to $\mathcal{Q}$. When $q>k$, $H^{q-1}(X, \mathcal{K})=H^{q-1}(X, \mathcal{Q})=0$. The relevant parts of the cohomology sequences associated to these exact sequences for $q>k$ become

$$
0 \rightarrow H^{q}(X, \mathcal{M}(-1)) \rightarrow H^{q}(X, \mathcal{N}) \rightarrow 0 \quad \text { and } \quad 0 \rightarrow H^{q}(X, \mathcal{N}) \rightarrow H^{q}(X, \mathcal{M}) \rightarrow 0
$$

respectively. Therefore the maps $H^{q}(X, \mathcal{M}(-1)) \rightarrow H^{q}(X, \mathcal{N}) \rightarrow H^{q}(X, \mathcal{M})$ are bijective, and this is true for every $\mathcal{O}$-module whose upport has dimension at most $k$, including the $\mathcal{O}$-module $\mathcal{M}(d)$. For every $d$, $H^{q}(X, \mathcal{M}(d-1)) \approx H^{q}(X, \mathcal{M}(d))$.

According to $7.5 .2,, \underset{\longrightarrow}{\lim } H^{q}(X, \mathcal{M}(d))=0$. It follows that $H^{q}(X, \mathcal{M}(d))=0$ for all $d$, and in particular, $H^{q}(X, \mathcal{M})=0$, when $q>k$.
proof of Theorem 7.7.2 (vanishing for a large twist)
Let $\mathcal{M}$ be a finite $\mathcal{O}$-module. Then $\mathcal{M}(r)$ is generated by global sections when $r$ is sufficiently large 6.7.21). Choosing generators gives us a surjective map $\mathcal{O}^{r} \rightarrow \mathcal{M}(r)$. Let $\mathcal{N}$ be the kernel of this map. When we twist the sequence $0 \rightarrow \mathcal{N} \rightarrow \mathcal{O}^{r} \rightarrow \mathcal{M}(r) \rightarrow 0$, we obtain short exact sequences

$$
\begin{equation*}
0 \rightarrow \mathcal{N}(d) \rightarrow \mathcal{O}(d)^{r} \rightarrow \mathcal{M}(d+r) \rightarrow 0 \tag{7.7.6}
\end{equation*}
$$

for every $d \geq 0$. These sequences are useful because $H^{q}(X, \mathcal{O}(d))=0$ when $q>0$.
We go to the proof of Theorem 7.7.2 now. We must show this:

Let $n$ be the dimnsion of $X$. By Theorem 7.7.1. $H^{q}(X, \mathcal{M})=0$ for any $\mathcal{O}$-module $\mathcal{M}$, when $q>n$, in particular, for the twists $\mathcal{M}(d)$ of $\mathcal{M}$. This leaves a finite set of integers $q=1, \ldots, n$ to consider, and it suffices to consider them one at a time. If $\left({ }^{*}\right)$ is true for each individual $q$ it will be true for the finite set of integers $q=1, \ldots, n$ at the same time, and therefore for all positive integers $q$, as the theorem asserts.

We use descending induction on $q$, the base case being $q=n+1$, for which (*) is true. We suppose that ${ }^{( }$) is true for every finite $\mathcal{O}$-module $\mathcal{M}$ when $q=p+1$, and that $p>0$, and we show that $\left(^{*}\right.$ ) is true for every finite $\mathcal{O}$-module $\mathcal{M}$ when $q=p$.

We substitute $q=p$ into the cohomology sequence associated to the sequence 7.7.6. The relevant part of that sequence is

$$
H^{p}\left(X, \mathcal{O}(d)^{m}\right) \rightarrow H^{p}(X, \mathcal{M}(d+r)) \xrightarrow{\delta^{p}} H^{p+1}(X, \mathcal{N}(d))
$$

Since $p$ is positive, $H^{p}(X, \mathcal{O}(d))=0$ for all $d \geq 0$, and therefore the map $\delta^{p}$ is injective. Our induction hypothesis, applied to the $\mathcal{O}$-module $\mathcal{N}$, shows that $H^{p+1}(X, \mathcal{N}(d))=0$ if $d$ is large, and then

$$
H^{p}(X, \mathcal{M}(d+r))=0
$$

The particular integer $d+r$ isn't useful. Our conclusion is that, for every finite $\mathcal{O}$-module $\mathcal{M}, H^{p}(X, \mathcal{M}(k))=$ 0 when $k$ is large enough.

## proof of Theorem 7.7 .3 (finiteness of cohomology)

This proof uses ascending induction on the dimension of support as well as descending induction on the degree $d$ of a twist. As was mentioned above, it isn't easy to prove directly that the space $H^{0}(X, \mathcal{M})$ of global sections is finite-dimensional.

We go back to the sequences (7.7.5) and their cohomology sequences. Ascending induction on the dimension of the support of $\mathcal{M}$ allows us to assume that $H^{r}(X, \mathcal{K})$ and $H^{r}(X, \mathcal{Q})$ are finite-dimensional for all $r$. Denoting finite-dimensional spaces ambiguously by $F$, the two cohomology sequencs become

$$
\cdots \rightarrow F \rightarrow H^{q}(X, \mathcal{M}(-1)) \rightarrow H^{q}(X, \mathcal{N}) \rightarrow F \rightarrow \cdots
$$

and

$$
\cdots \rightarrow F \rightarrow H^{q}(X, \mathcal{N}) \rightarrow H^{q}(X, \mathcal{M}) \rightarrow F \rightarrow \cdots
$$

The first sequence shows that if $H^{q}(X, \mathcal{M}(-1))$ has infinite dimension, then $H^{q}(X, \mathcal{N})$ also has infinite dimension, and the second sequence shows that if $H^{q}(X, \mathcal{N})$ has infinite dimension, then $H^{q}(X, \mathcal{M})$ has infinite dimenson too. This applies to the twisted module $\mathcal{M}(d)$ as well as to $\mathcal{M}$. Therefore $H^{q}(X, \mathcal{M}(d-1))$ and $H^{q}(X, \mathcal{M}(d))$ are either both finite-dimensional, or else they are both infinite-dimensional, and this is true for every $d$.

Suppose that $q>0$. Then $H^{q}(X, \mathcal{M}(d))=0$ when $d$ is large enough (Theorem 7.7.2. Since the zero space is finite-dimensional, we can use the sequence together with descending induction on $d$, to conclude that $H^{q}(X, \mathcal{M}(d))$ is finite-dimensional for every finite module $\mathcal{M}$ and every $d$. In particular, $H^{q}(X, \mathcal{M})$ is finite-dimensional.

This leaves the case that $q=0$. To prove that $H^{0}(X, \mathcal{M})$ is finite-dimensional, we set $d=-r$ in the sequence 7.7.6:

$$
0 \rightarrow \mathcal{N}(-r) \rightarrow \mathcal{O}(-r)^{m} \rightarrow \mathcal{M} \rightarrow 0
$$

The corresponding cohomology sequence is

$$
0 \rightarrow H^{0}(X, \mathcal{N}(-r)) \rightarrow H^{0}(X, \mathcal{O}(-r))^{m} \rightarrow H^{0}(X, \mathcal{M}) \xrightarrow{\delta^{0}} H^{1}(X, \mathcal{N}(-r)) \rightarrow \cdots .
$$

Here $H^{0}(X, \mathcal{O}(-r))^{m}=0$, and we've shown that $H^{1}(X, \mathcal{N}(-r))$ is finite-dimensional. It follows that $H^{0}(X, \mathcal{M})$ is finite-dimensional, and this completes the proof.

Notice that the finiteness of $H^{0}$ comes out only at the end. The higher cohomology is essential for the proof.
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Theorems 7.7.1 and 7.7.3 allow us to define the Euler characteristic of a finite module on projective variety.
7.7.8. Definition. Let $X$ be a projective variety. The Euler characteristic of a finite $\mathcal{O}$-module $\mathcal{M}$ is the alternating sum of the dimensions of the cohomology:

$$
\begin{equation*}
\chi(\mathcal{M})=\sum(-1)^{q} \mathbf{h}^{q}(X, \mathcal{M}) \tag{7.7.9}
\end{equation*}
$$

This makes sense because $\mathbf{h}^{q}(X, \mathcal{M})$ is finite for every $q$, and is zero when $q$ is large.
Try not to confuse the Euler characterstic of an $\mathcal{O}$-module with the topological Euler characteristic of the variety $X$.
7.7.10. Proposition. (i) If $0 \rightarrow \mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{P} \rightarrow 0$ is a short exact sequence of finite $\mathcal{O}$-modules on a projective variety $X$, then $\chi(\mathcal{M})-\chi(\mathcal{N})+\chi(\mathcal{P})=0$.
(ii) If $0 \rightarrow \mathcal{M}_{0} \rightarrow \mathcal{M}_{1} \rightarrow \cdots \rightarrow \mathcal{M}_{n} \rightarrow 0$ is an exact sequence of finite $\mathcal{O}$-modules on $X$, the alternating sum $\sum(-1)^{i} \chi\left(\mathcal{M}_{i}\right)$ is zero.
7.7.11. Lemma. Let $0 \rightarrow V^{0} \rightarrow V^{1} \rightarrow \cdots \rightarrow V^{n} \rightarrow 0$ be an exact sequence of finite dimensional vector spaces. The alternating sum $\sum(-1)^{q} \operatorname{dim} V^{q}$ is zero.
proof of Proposition 7.7.10 (i) Let $n$ be the dimension of $X$. The cohomology sequence associated to the given sequence is

$$
0 \rightarrow H^{0}(\mathcal{M}) \rightarrow H^{0}(\mathcal{N}) \rightarrow H^{0}(\mathcal{P}) \rightarrow H^{1}(\mathcal{M}) \rightarrow H^{1}(\mathcal{N}) \rightarrow H^{1}(\mathcal{P}) \rightarrow \cdots \rightarrow H^{n}(\mathcal{P}) \rightarrow 0
$$

and the lemma tells us that the alternating sum of its dimensions is zero. That alternating sum is also equal to $\chi(\mathcal{M})-\chi(\mathcal{N})+\chi(\mathcal{P})$.
(ii) Let 's denote the given sequence by $\mathbb{S}_{0}$ and the alternating sum $\sum_{i}(-1)^{i} \chi\left(\mathcal{M}_{i}\right)$ by $\chi\left(\mathbb{S}_{0}\right)$.

Let $\mathcal{N}=\mathcal{M}_{1} / \mathcal{M}_{0}$. The sequence $\mathbb{S}_{0}$ decomposes into the two exact sequences

$$
\mathbb{S}_{1}: \quad 0 \rightarrow \mathcal{M}_{0} \rightarrow \mathcal{M}_{1} \rightarrow \mathcal{N} \rightarrow 0 \quad \text { and } \quad \mathbb{S}_{2}: \quad 0 \rightarrow \mathcal{N} \rightarrow \mathcal{M}_{2} \rightarrow \cdots \rightarrow \mathcal{M}_{k} \rightarrow 0
$$

One sees directly that $\chi\left(\mathbb{S}_{0}\right)=\chi\left(\mathbb{S}_{1}\right)-\chi\left(\mathbb{S}_{2}\right)$, and the assertion follows from (i) by induction on $n$.

### 7.8 Bézout's Theorem

As an application of cohomology, we use it to prove Bézout's Theorem.
If $f(x)=p_{1}(x)^{e_{1}} \cdots p_{k}(x)^{e_{k}}$ is a factorization of a homogeneous polyomial $f\left(x_{0}, x_{1}, x_{2}\right)$ into irreducible polynomials, the divisor of $f$ is defined to be the integer combination $e_{1} C_{1}+\cdots+e_{k} C_{k}$, where $C_{i}$ is the curve of zeros of $p_{i}$.

We restate the theorem to be proved.
7.8.1. Bézout's Theorem. Let $Y$ and $Z$ be the divisors in the projective plane $X$ defined by relatively prime homogeneous polynomials $f$ and $g$ of degrees $m$ and $n$, respectively. The number of intersection points $Y \cap Z$, counted with an appropriate multiplicity, is equal to $m n$. Moreover, the multiplicity is 1 at a point at which $Y$ and $Z$ intersect transversally.

The definition of the multiplicity will emerge during the proof.
7.8.2. Example. Suppose that $f$ and $g$ are products of linear polynomials, so that $Y$ is the union of $m$ lines and $Z$ is the union of $n$ lines, and suppose that those lines are distinct. Since distinct lines intersect transversally in a single point, there are $m n$ intersection points of multiplicity 1.
proof of Bézout's Theorem. Multiplication by $f$ defines a short exact sequence

$$
0 \rightarrow \mathcal{O}_{X}(-m) \xrightarrow{f} \mathcal{O}_{X} \rightarrow \mathcal{A} \rightarrow 0
$$

where $\mathcal{A}$ stands for the quotient $\mathcal{O} / f \mathcal{O}(-m)$. If $f$ is a irreducible, $\mathcal{A}$ will be the extension by zero of the structure sheaf on $Y$. When $f$ is reducible, $Y$ won't be a variety. Let's use that imprecise notation anyway. This sequence describes $\mathcal{O}_{X}(-m)$ as the ideal $\mathcal{I}$ of $Y$, and there is a similar sequence describing the module $\mathcal{O}_{X}(-n)$ as the ideal $\mathcal{J}$ of $Z$. The zero locus of the ideal $\mathcal{I}+\mathcal{J}$ is the intersection $Y \cap Z$.

Let $\overline{\mathcal{O}}$ denote the quotient $\mathcal{O}_{X} /(\mathcal{I}+\mathcal{J})$. Since $f$ and $g$ have no common factor, $Y \cap Z$ is a finite set of points $\left\{p_{1}, \ldots, p_{k}\right\}$, and $\overline{\mathcal{O}}$ is isomorphic to a direct sum $\bigoplus_{\overline{\mathcal{O}}}^{i}$, where $\overline{\mathcal{O}}_{i}$ is a finite-dimensional algebra whose support is $p_{i}$ 6.6.4. We define the intersection multiplicity of $Y$ and $Z$ at $p_{i}$ to be the dimension of $\overline{\mathcal{O}}_{i}$, which is also equal to the dimension $\mathbf{h}^{0}\left(X, \overline{\mathcal{O}}_{i}\right)$ of its space of its global sections. Let's denote that multiplicity by $\mu_{i}$. The dimension of $H^{0}(X, \overline{\mathcal{O}})$ is the sum $\mu_{1}+\cdots+\mu_{k}$, and $H^{q}(X, \overline{\mathcal{O}})=0$ for all $q>0$ (Theorem 7.7.1 . So the Euler characteristic $\chi(\overline{\mathcal{O}})$ is equal to $\mathbf{h}^{0}(X, \overline{\mathcal{O}})$. We'll show that $\chi(\overline{\mathcal{O}})=m n$, and therefore that $\mu_{1}+\cdots+\mu_{k}=m n$. This will prove Bézout's Theorem.

We form an exact sequence of $\mathcal{O}$-modules, in which $\mathcal{O}=\mathcal{O}_{X}$ :

$$
\begin{equation*}
0 \rightarrow \mathcal{O}(-m-n) \xrightarrow{(g, f)^{t}} \mathcal{O}(-m) \times \mathcal{O}(-n) \xrightarrow{(f,-g)} \mathcal{O} \xrightarrow{\pi} \overline{\mathcal{O}} \rightarrow 0 \tag{7.8.3}
\end{equation*}
$$

In order to interpret the maps in this sequence as matrix multiplication with homomorphisms acting on the left, sections of $\mathcal{O}(-m) \times \mathcal{O}(-n)$ should be represented as column vectors $(u, v)^{t}, u$ and $v$ being sections of $\mathcal{O}(-m)$ and $\mathcal{O}(-n)$, respectively.
7.8.4. Lemma. The sequence 7.8.3 is exact.
proof. To prove exactness, it suffices to show that the sequence of sections on each of the standard affine open sets is exact. Let's suppose that coordinates have been chosen so that none of the points making up $Y \cap Z$ lie on the coordinate axes. We look at $\mathbb{U}^{0}$, as usual. Let $A$ be the algebra of regular functions on $\mathbb{U}^{0}$, the polynomial algebra $\mathbb{C}\left[u_{1}, u_{2}\right]$, with $u_{i}=x_{i} / x_{0}$. We identify $\mathcal{O}(k)$ with $\mathcal{O}(k H), H$ being the hyperplane at infinity. The restriction of the module $\mathcal{O}(k H)$ to $\mathbb{U}^{0}$ is isomorphic to $\mathcal{O}_{\mathbb{U}^{0}}$. Its sections on $\mathbb{U}^{0}$ are the elements of $A$. Let $\bar{A}$ be the algebra of sections of $\overline{\mathcal{O}}$ on $\mathbb{U}^{0}$. Since $f$ and $g$ are relatively prime, so are their dehomogenizations $F=f\left(1, u_{1}, u_{2}\right)$ and $G=g\left(1, u_{1}, u_{2}\right)$. The sequence of sections of 7.8 .3 on $\mathbb{U}^{0}$ is

$$
0 \rightarrow A \xrightarrow{(G, F)^{t}} A \times A \xrightarrow{(F,-G)} A \rightarrow \bar{A} \rightarrow 0
$$

and the only place at which exactness of this sequence isn't obvious is at $A \times A$. Suppose that $(u, v)^{t}$ is in the kernel of the map $(F,-G)$, i.e., that $F u=G v$ Since $F$ and $G$ are relatively prime, $F$ divides $v, G$ divides $u$, and $v / F=u / G$. Let $w=v / F=u / G$. Then $(u, v)^{t}=(G, F)^{t} w$.

We go back to the proof of Bézout's Theorem. Proposition 7.7.10(ii), applied to the exact sequence (7.8.3), tells us that the alternating sum

$$
\begin{equation*}
\chi(\mathcal{O}(-m-n))-[\chi(\mathcal{O}(-m))+\chi(\mathcal{O}(-n) \times \mathcal{O}(-n))-\chi(\overline{\mathcal{O}}) \tag{7.8.5}
\end{equation*}
$$

is zero. Since cohomology is compatible with products, $\chi(\mathcal{M} \times \mathcal{N})=\chi(\mathcal{M})+\chi(\mathcal{N})$ for any $\mathcal{O}$-modules $\mathcal{M}$ and $\mathcal{N}$. Solving for $\chi(\overline{\mathcal{O}})$ and applying Theorem 7.5.4

$$
\chi(\overline{\mathcal{O}})=\binom{n+m-1}{2}-\binom{m-1}{2}-\binom{n-1}{2}+1
$$

This equation shows that the term $\chi(\overline{\mathcal{O}})$ depends only on the integers $m$ and $n$. Since we know that the answer is $m n$ when $Y$ and $Z$ are unions of distinct lines, it is $m n$ in every case. This completes the proof.

If you are suspicious of this reasoning, you can evaluate the right side of the equation.
We still need to explain the assertion that the mutiplicity at a transversal intersection $p$ is equal to 1 . This will be true if and only if $\mathcal{I}+\mathcal{J}$ generates the maximal ideal $\mathfrak{m}$ of $A=\mathbb{C}[y, z]$ at $p$ locally, and it is obvious when $Y$ and $Z$ are lines. In that case we may choose affine coordinates so that $p$ is the origin in $\mathbb{A}^{2}=\operatorname{Spec} A$ and the curves are the coordinate axes $\{z=0\}$ and $\{y=0\}$. The variables $y$, $z$ generate the maximal ideal at the origin, so the quotient algebra $k=A / \mathfrak{m}$ has dimension 1 .
bezoutresolution
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chiforbezout

Suppose that $Y$ and $Z$ intersect transverally at $p$, but that they aren't lines. We choose affine coordinates so that $p$ is the origin and that the tangent directions of $Y$ and $Z$ at $p$ are the coordinate axes. The affine equations of $Y$ and $Z$ will have the form $y^{\prime}=0$ and $z^{\prime}=0$, where $y^{\prime}=y+g(y, z)$ and $z^{\prime}=z+h(y, z), g$ and $h$ being polynomials all of whose terms have degree at least 2 . Because $Y$ and $Z$ may intersect at points other than $p$, the elements $y^{\prime}$ and $z^{\prime}$ may not generate the maximal ideal $\mathfrak{m}$ at $p$. However, it suffices to show that they generate the maximal ideal locally.

The locus $Y \cap Z$ in the plane $X$ is the intersection of two distinct curves in $X$, so it is a finite set. We choose an element $s$ of $A$ that is not zero at $p$, but is zero at the other points of $Y \cap Z$. Then in $X_{s}$ the locus $y^{\prime}=z^{\prime}=0$ is $p$.

In $A_{s}$, let $I$ be the ideal generated by $x^{\prime}, y^{\prime}$. We map a free $A_{s}$-module $V$ with basis $x^{\prime}, y^{\prime}$ to $I$. Let $C$ be the cokernel of that map. Tensoring the exact sequence $\bar{V} \rightarrow I \rightarrow C \rightarrow 0$ with the residue field $k=A / \mathfrak{m}$ gives us an exact sequence $\bar{V} \rightarrow \bar{I} \rightarrow \bar{C} \rightarrow 0$. The map $\bar{V} \rightarrow \bar{I}$ is surjective, so $\bar{C}=0$. The Local Nakayama Lemma tells us that the localization of $C$ at $p$ is zero. Then, since $C$ is supported at $p, C=0$.

# Chapter 8 THE RIEMANN-ROCH THEOREM FOR CURVES 

july 11
8.1 Branched CoveringsDivisors
The Riemann-Roch Theorem
8.4 The Birkhoff-Grothendieck Theorem
8.5 Differentials
8.6 Trace
8.7 The Riemann-Roch Theorem II
8.8 Using Riemann-Roch

The topic of this chapter is a classical problem of algebraic geometry, to determine the rational functions on a smooth projective curve with given poles. This is can be difficult. The functions whose poles have orders at most $r_{i}$ at $p_{i}, \mathrm{i}=1, \ldots, \mathrm{k}$, form a vector space, and one is usually happy if one can determine the dimension of that space. The most important tool for determining the dimension is the Riemann-Roch Theorem.

### 8.1 Branched Coverings

covercurve

Smooth affine curves were discussed in Chapter ??. An affine curve is smooth if its local rings are valuation rings, or if its coordinate ring is a normal domain. An arbitrary curve is smooth if it has an open covering by smooth affine curves.

We take a brief look at modules on a smooth curve.
Recall that a module over a domain $A$ is torsion-free if its only torsion elementis zero 2.6.11. This definition is extended to $\mathcal{O}$-modules by applying it to affine open sets.
8.1.1. Lemma. Let $Y$ be a smooth curve.
(i) A finite $\mathcal{O}$-module $\mathcal{M}$ is locally free if and only if it is torsion-free.
(ii) An $\mathcal{O}$-module $\mathcal{M}$ that isn't torsion-free has a nonzero global section.
proof. (i) We may assume that $Y$ is affine, $Y=\operatorname{Spec} B$, and that $\mathcal{M}$ is the $\mathcal{O}$-module associated to a $B$-module $M$. Let $\widetilde{B}$ and $\widetilde{M}$ be the localizations of $B$ and $M$ at a point $q$, respectively. Then $\widetilde{M}$ is a finite, torsion-free module over the local ring $\widetilde{B}$. It suffices to show that, for every point $q$ of $Y, \widetilde{M}$ is a free $\widetilde{B}$-module 5.1.15, The local ring $\widetilde{B}$ is a valuation ring. It is a principal ideal domain its nonzero ideals are powers of the maximal ideal $\widetilde{\mathfrak{m}}$, which is a principal ideal. Every finite, torsion-free module over a principal ideal domain is free.
(ii) If the torsion submodule of $\mathcal{M}$ isn't zero, there will be an affine open set $U$, with nonzero elements $m$ in $\mathcal{M}(U)$ and $a$ in $\mathcal{O}(U)$, such that $a m=0$. Let $C$ be the finite set of zeros of $a$ in $U$, and let $V=Y-C$ be the complement of $C$ in $Y$. Then $a$ is invertible on the intersection $W=U \cap V$, and since $a m=0$, the restriction of $m$ to $W$ is zero.

The open sets $U$ and $V$ cover $Y$, and the sheaf property for this covering can be written as an exact sequence

$$
0 \rightarrow \mathcal{M}(Y) \rightarrow \mathcal{M}(U) \times \mathcal{M}(V) \xrightarrow{+,-} \mathcal{M}(W)
$$

(Lemma 6.3.11). In this sequence, the section $(m, 0)$ of $\mathcal{M}(U) \times \mathcal{M}(V)$ maps to zero in $\mathcal{M}(W)$. Therefore it is the image of a nonzero global section of $\mathcal{M}$.
8.1.2. Lemma. Let $Y$ be a smooth curve. Every nonzero ideal $\mathcal{I}$ of $\mathcal{O}_{Y}$ is a product of powers of maximal ideals of $\mathcal{O}_{Y}: \quad \mathcal{I}=\mathfrak{m}_{1}^{e_{1}} \cdots \mathfrak{m}_{k}^{e_{k}}$.
proof. This follows for any smooth curve from the case that $Y$ is affine, which is Proposition 5.3.7

We come now the main topic of this section: branched coverings of smooth curves.
An integral morphism $Y \xrightarrow{\pi} X$ of smooth curves will be called a branched covering. It follows from Chevalley's Finiteness Theorem that every morphism of smooth projective curves is a branched covering, unless it maps $Y$ to a point.

Let $Y \rightarrow X$ be a branched covering. The function field $K$ of $Y$ will be a finite extension of the function field $F$ of $X$. The degree of the covering is the degree $[K: F]$ of that field extension. It will be denoted by $[Y: X]$. If $X^{\prime} \operatorname{Spec} A$ is an affine open subset of $X$, its inverse image $Y^{\prime}$ will be an affine open subset $Y^{\prime}=\operatorname{Spec} B$ of $Y$, and $B$ will be a locally free $A$-module whose rank is the degee $[Y: X]$ of the covering.

To describe the fibre of a branched covering $Y \xrightarrow{\pi} X$ over a point $p$ of $X$, we may localize. So we may assume that $X$ and $Y$ are affine, $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$, and that the maximal ideal $\mathfrak{m}_{p}$ of $A$ at a point $p$ is a principal ideal, generated by an element $x$ of $A$. If a point $q$ of $Y$ lies over $p$, the ramification index at $q$, which we denote by $e$, is defined to be $\mathrm{v}_{q}(x)$, where $\mathrm{v}_{q}$ is the valuation of the function field $K$ corresponding to $q$. WE usually denote theramification index by $e$. Then, if $y$ is a local generator for the maximal ideal $\mathfrak{m}_{q}$ of $B$ at $q$, we will have

$$
x=u y^{e}
$$

where $u$ is a local unit. (A rational function is a local unit at a point $p$ if it is a unit in some open neighborhood of $p$.)

The next lemma follows from Lemma 8.1.2 and the Chinese Remainder Theorem.
8.1.3. Lemma. Let $Y \xrightarrow{\pi} X$ be a branched covering, with $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$. Let $q_{1}, \ldots, q_{k}$ be the points of $Y$ that lie over a point $p$ of $X$, let $x$ be a generator for the the maximal ideal $\mathfrak{m}_{p}$ at $p$, and let $\mathfrak{m}_{i}$ and $e_{i}$ be the maximal ideal and ramification index at $q_{i}$, respectively.
(i) The extended ideal $\mathfrak{m}_{p} B=x B$ is the product ideal $\mathfrak{m}_{1}^{e_{1}} \cdots \mathfrak{m}_{k}^{e_{k}}$.
(ii) Let $\bar{B}_{i}=B / \mathfrak{m}_{i}^{e_{i}}$. The quotient $\bar{B}=B / x B$ is isomorphic to the product $\bar{B}_{1} \times \cdots \times \bar{B}_{k}$.
(iii) The degree $[Y: X]$ of the covering is the sum $e_{1}+\cdots+e_{k}$ of the ramification indices at the points $q_{i}$.

Points of $Y$ whose ramification indices are greater than one are called branch points. We also call a point $p$ of $X$ a branch point of the covering if there is a branch point of $Y$ whose image is $p$.
8.1.4. Lemma. A branched covering $Y \rightarrow X$ has finitely many branch points. If a point $p$ of $X$ is not a branch point, the fibre over $p$ consists of $n=[Y: X]$ points with ramification indices equal to 1 .
proof. This is very simple. We can delete finite sets of points, so we may suppose that $X$ and $Y$ are affine, $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$. Then $B$ is a finite $A$-module of rank $n$. Let $F$ and $K$ be the fraction fields of $A$ and $B$, respectively, and let $\beta$ be an element of $B$ that generates the field extension $K / F$. Then $A[\beta] \subset B$, and since these two rings have the same fraction field, there will be a nonzero element $s \in A$ such that $A_{s}[\beta]=B_{s}$. We may suppose that $B=A[\beta]$. Let $g$ be the monic irreducible polynomial for $\beta$ over $A$. The discriminant of $g$ is not the zero ideal 1.7.19, so for all but finitely many points $p$ of $X$, there will be $n$ points of $Y$ over $p$ with ramification indices equal to 1 .

### 8.1.5. Corollary. A branched covering $Y \xrightarrow{\pi} X$ of degree one is an isomorphism.

proof. When $[Y: X]=1$, the function fields of $Y$ and $X$ are equal. Then, because $Y \rightarrow X$ is an integral morphism and $X$ is normal, $Y=X$.

> figure: a branched covering

## (8.1.6) local analytic structure

The local analytic structure of a branched covering $Y \xrightarrow{\pi} X$ in the classical topology is very simple. We explain it here because it is helpful for intuition. It is also useful.
extende-didealisprod
8.1.7. Proposition. In the classical topology, $Y$ is locally isomorphic to the $e$-th root covering $y^{e}=x$.
proof. Let $q$ be a point of $Y$, let $p$ be its image in $X$, let $x$ and $z$ be local generators for the maximal ideals $\mathfrak{m}_{p}$ of $\mathcal{O}_{X}$ at $p$, and $\mathfrak{m}_{q}$ of $\mathcal{O}_{Y}$ at $q$, respectively. Let $e=\mathrm{v}_{q}(x)$ be the ramification index at $q$. Then $x$ has the form $u z^{e}$, where $u$ is a local unit at $q$. In a neighborhood of $q$ in the classical topology, $u$ will have an analytic $e$-th root $w$. The element $y=w z$ also generates $\mathfrak{m}_{q}$ locally, and $x=y^{e}$. It follows from the implicit function theorem that $x$ and $y$ are local analytic coordinate functions on $X$ and $Y$ (see (??)).
8.1.8. Corollary. Let $Y \xrightarrow{\pi} X$ be a branched covering, let $\left\{q_{1}, \ldots, q_{k}\right\}$ be the fibre over a point $p$ of $X$, and let $e_{i}$ be the ramification index at $q_{i}$. As a point $p^{\prime}$ of $X$ approaches $p$, $e_{i}$ points of the fibre over $p^{\prime}$ approach $q_{i}$.
8.1.9. Notation. When considering a branched covering $Y \xrightarrow{\pi} X$ of smooth curves, we will often pass between an $\mathcal{O}_{Y}$-module $\mathcal{M}$ and its direct image $\pi_{*} \mathcal{M}$, and it will be convenient to work primarily on $X$. Recall that if $Y^{\prime}$ is the inverse image of an open subset $X^{\prime}$ of $X$, then

$$
\left[\pi_{*} \mathcal{M}\right]\left(X^{\prime}\right)=\mathcal{M}\left(Y^{\prime}\right)
$$

One can think of the direct image $\pi_{*} \mathcal{M}$ as working with $\mathcal{M}$, but looking only at open subsets $Y^{\prime}$ that are inverse images of open subsets $X^{\prime}$ of $X$. If we look only at such subsets, the only significant difference between $\mathcal{M}$ and its direct image will be that the $\mathcal{O}_{Y}\left(Y^{\prime}\right)$-module $\mathcal{M}\left(Y^{\prime}\right)$ is made into an $\mathcal{O}_{X}\left(X^{\prime}\right)$-module by restriction of scalars.

To simplify notation, we will often drop the symbol $\pi_{*}$, and write $\mathcal{M}$ instead of $\pi_{*} \mathcal{M}$. If $X^{\prime}$ is an open subset of $X, \mathcal{M}\left(X^{\prime}\right)$ will stand for $\left[\pi_{*} \mathcal{M}\right]\left(X^{\prime}\right)=\mathcal{M}\left(\pi^{-1} X^{\prime}\right)$. When denoting the direct image of an $\mathcal{O}_{Y^{-}}$ module $\mathcal{M}$ by the same symbol $\mathcal{M}$, we may refer to it as an $\mathcal{O}_{X}$-module. In accordance with this convention, we may also write $\mathcal{O}_{Y}$ for $\pi_{*} \mathcal{O}_{Y}$, but we must be careful to include the subscript $Y$.

This abbreviation is analogous to the one used for restriction of scalars in a module. When $A \rightarrow B$ is an algebra homomorphism, the $A$-module obtained from a $B$-module $M$ by restriction of scalars is usually denoted by the same letter $M$.
8.1.10. Lemma. Let $Y \rightarrow X$ be a branched covering of smooth curves.
(i) A finite $\mathcal{O}_{Y}$-module $\mathcal{M}$ is a torsion $\mathcal{O}_{Y}$-module if and only if it is a torsion $\mathcal{O}_{X}$-module .
(ii) A finite $\mathcal{O}_{Y}$-module $\mathcal{M}$ is a locally free $\mathcal{O}_{Y}$-module if and only if it is a locally free $\mathcal{O}_{X}$-module. If $\mathcal{M}$ is a locally free $\mathcal{O}_{Y}$-module of rank $r$, then its rank as $\mathcal{O}_{X}$-module will be $n r$, where $n=[Y: X]$ is the degree of the covering.

## (8.1.11) the module of homomorphisms

We begin by discussing homomorphisms of modules over a ring.
Let $M$ and $N$ be modules over a noetherian ring $A$. The set of homomorphisms $M \rightarrow N$ is usually denoted by $\operatorname{Hom}_{A}(M, N)$. It becomes an $A$-module with some fairly obvious laws of composition: If $\varphi$ and $\psi$ are homomorphisms and $a$ is an element of $A$, then $\varphi+\psi$ and $a \varphi$ are defined by

$$
\begin{equation*}
[\varphi+\psi](m)=\varphi(m)+\psi(m) \quad \text { and } \quad[a \varphi](m)=a \varphi(m) \tag{8.1.12}
\end{equation*}
$$

Because $\varphi$ is a module homomorphism, we also have $\varphi\left(m_{1}+m_{2}\right)=\varphi\left(m_{1}\right)+\varphi\left(m_{2}\right)$ and $a \varphi(m)=\varphi(a m)$.
8.1.13. Lemma. (i) An A-module $N$ is canonically isomorphic to $\operatorname{Hom}_{A}(A, N)$. The homomorphism $A \xrightarrow{\varphi}$ $N$ that corresponds to an element $n$ of $N$ is multiplication by $n: ~ \varphi(a)=a n$. The element of $N$ that corresponds to a homomorphism $A \xrightarrow{\varphi} N$ is $n=\varphi(1)$.
(ii) $\operatorname{Hom}_{A}\left(A^{k}, N\right)$ is isomorphic to $N^{k}$, and $\operatorname{Hom}_{A}\left(A^{k}, A^{\ell}\right)$ is isomorphic to the module $A^{\ell \times k}$ of $k \times \ell A$ matrices.

A presentation of an $A$-module $M$ is an exact sequence of the form $A^{\ell} \rightarrow A^{k} \rightarrow M \rightarrow 0$. Every finite module over a noetherian ring $A$ has such a presentation.
8.1.14. Lemma.
(i) The functor $\mathrm{Hom}_{A}$ is a left exact and contravariant in the first variable. For any $A$-module $N$, an exact sequence $M_{1} \rightarrow M_{2} \rightarrow M_{3} \rightarrow 0$ of $A$-modules induces an exact sequence

$$
0 \rightarrow \operatorname{Hom}_{A}\left(M_{3}, N\right) \rightarrow \operatorname{Hom}_{A}\left(M_{2}, N\right) \rightarrow \operatorname{Hom}_{A}\left(M_{1}, N\right)
$$

(ii) The functor $\operatorname{Hom}_{A}$ is a left exact and covariant in the second variable. For any A-module $M$, an exact sequence $0 \rightarrow N_{1} \rightarrow N_{1} \rightarrow N_{3}$ of A-modules induces an exact sequence

$$
0 \rightarrow \operatorname{Hom}_{A}\left(M, N_{1}\right) \rightarrow \operatorname{Hom}_{A}\left(M, N_{2}\right) \rightarrow \operatorname{Hom}_{A}\left(M, N_{3}\right)
$$

8.1.15. Corollary. If $M$ and $N$ are finite $A$-modules over a notherian ring $A$, then $\operatorname{Hom}_{A}(M, N)$ is a finite A-module.

This follows when part (i) of the lemma is applied to a presentation of $M$.
The module Hom is compatible with localization:
8.1.16. Lemma. Let $M$ and $N$ be modules over a noetherian domain $A$, and suppose that $M$ is a finite module. Let $S$ be a multiplicative system in $A$. The localization $S^{-1} \operatorname{Hom}_{A}(M, N)$ is canonically isomorphic to $\operatorname{Hom}_{S^{-1} A}\left(S^{-1} M, S^{-1} N\right)$.
proof. We choose a presentation $A^{\ell} \rightarrow A^{k} \rightarrow M \rightarrow 0$ of the $A$-module $M$. Its localization, which is $\left(S^{-1} A\right)^{\ell} \rightarrow\left(S^{-1} A\right)^{k} \rightarrow S^{-1} M \rightarrow 0$ is a presentation of the $S^{-1} A$-module $S^{-1} M$. The sequence

$$
0 \rightarrow \operatorname{Hom}_{A}(M, N) \rightarrow \operatorname{Hom}_{A}\left(A^{k}, N\right) \rightarrow \operatorname{Hom}_{A}\left(A^{\ell}, N\right)
$$

is exact, as is its localization. So it suffices to prove the lemma in the case that $M=A$. It is true in that case.
This lemma shows that, when $\mathcal{M}$ and $\mathcal{N}$ are finite $\mathcal{O}$-modules on a variety $X$, there is an $\mathcal{O}$-module of homomorphisms $\mathcal{M} \rightarrow \mathcal{N}$. This $\mathcal{O}$-module is usually denoted by $\operatorname{Hom}_{\mathcal{O}}(\mathcal{M}, \mathcal{N})$. If $U=\operatorname{Spec} A$ is an affine open set, $M=\mathcal{M}(U)$ and $N=\mathcal{N}(U)$, the module of sections of $\operatorname{Hom}_{\mathcal{O}}(\mathcal{M}, \mathcal{N})$ on $U$ is the $A$ module $\operatorname{Hom}_{A}(M, N)$. We are using the symbol Hom here because the vector space of homomorphisms $\mathcal{M} \rightarrow \mathcal{N}$ defined on all of $X$, which is the space of global sections of $\underline{\operatorname{Hom}}_{\mathcal{O}}(\mathcal{M}, \mathcal{N})$, is customarily denoted by $\operatorname{Hom}_{\mathcal{O}}(\mathcal{M}, \mathcal{N})$.

The analogues of Lemma 8.1.13 and lemma 8.1.14 are true for Hom:
8.1.17. Corollary. (i) An $\mathcal{O}$-module $\mathcal{M}$ on a smooth curve $Y$ is isomorphic to $\underline{\operatorname{Hom}}_{\mathcal{O}}(\mathcal{O}, \mathcal{M})$.
(ii) The functor Hom is left exact and contravariant in the first variable, and it is left exact and covariant in the first variable.

Notation. The notation $\operatorname{Hom}_{A}(M, N)$ is cumbersome. It seems permissible to drop the symbol Hom, and to write ${ }_{A}(M, N)$ for $\operatorname{Hom}_{A}(M, N)$. Similarly, if $\mathcal{M}$ and $\mathcal{M}$ are $\mathcal{O}$-modules on a variety $X$, we will write $\mathcal{O}(\mathcal{M}, \mathcal{N})$ or $x_{x}(\mathcal{M}, \mathcal{N})$ for $\underline{\operatorname{Hom}}_{\mathcal{O}}(\mathcal{M}, \mathcal{N})$.

## (8.1.18) the dual module

The dual module $\mathcal{M}^{*}$ of a locally free $\mathcal{O}$-module $\mathcal{M}$ is the $\mathcal{O}$-module $\mathcal{O}(\mathcal{M}, \mathcal{O})$ of $\mathcal{O}$-module homomorphisms $\mathcal{M} \rightarrow \mathcal{O}$. A section of $\mathcal{M}^{*}$ on an open set $U$ is a $\mathcal{O}(U)$-module homomorphism $\mathcal{M}(U) \rightarrow \mathcal{O}(U)$.

The dualizing operation is contravariant. A homomorphism $\mathcal{M} \rightarrow \mathcal{N}$ of locally free $\mathcal{O}$-modules induces a homomorphism $\mathcal{M}^{*} \leftarrow \mathcal{N}^{*}$.

If $\mathcal{M}$ is a free module with basis $v_{1}, \ldots, v_{k}$, then $\mathcal{M}^{*}$ will also be free, with the dual basis $v_{i}^{*}$ defined by $v_{i}^{*}\left(v_{i}\right)=1$ and $v_{i}^{*}\left(v_{j}\right)=0$ if $i \neq j$. Therefore, when $\mathcal{M}$ is locally free, $\mathcal{M}^{*}$ is also locally free. The dual $\mathcal{O}^{*}$ of the structure sheaf $\mathcal{O}$ is the module $\mathcal{O}$ itself. If $\mathcal{M}$ and $\mathcal{N}$ are locally free $\mathcal{O}$-modules, the dual $\left(\mathcal{M} \otimes_{\mathcal{O}} \mathcal{N}\right)^{*}$ is isomorphic to the tensor product $\mathcal{M}^{*} \otimes_{\mathcal{O}} \mathcal{N}^{*}$.

There is a canonical $\mathcal{O}$-bilinear map $\mathcal{M}^{*} \times \mathcal{M} \rightarrow \mathcal{O}$. If $\mu$ and $m$ are sections of $\mathcal{M}^{*}$ and $\mathcal{M}$, respectively, the bilinear map evaluates $\mu$ at $m:\langle\mu, m\rangle=\mu(m)$.
bidualm 8.1.19. Corollary. (i) A locally free $\mathcal{O}$-module $\mathcal{M}$ is canonically isomorphic to its bidual: $\left(\mathcal{M}^{*}\right)^{*} \approx \mathcal{M}$.
(ii) If $\mathcal{M}$ and $\mathcal{N}$ are locally free $\mathcal{O}$-modules, $\mathcal{M}^{*} \otimes_{\mathcal{O}} \mathcal{N}^{*}$ is isomorphic to $\left(\mathcal{N} \otimes_{\mathcal{O}} \mathcal{M}\right)^{*}$.
\#\#any functor commutes with $\oplus$ \#\#
dualseq 8.1.20. Proposition. Let $\mathcal{O}$ be the strucdture sheaf on a variety $X$.
(i) Let $\mathcal{P} \xrightarrow{f} \mathcal{N} \xrightarrow{g} \mathcal{P}$ be homomorphisms of $\mathcal{O}$-modules whose composition gf is the identity map on $\mathcal{P}$. Then $\mathcal{N}$ is the direct sum of the image of $f$, which is isomorphic to $\mathcal{P}$, and the kernel $\mathcal{K}$ of $g: \mathcal{N} \approx \mathcal{P} \oplus \mathcal{K}$.
(ii) Let $g: \mathcal{N} \rightarrow \mathcal{P}$ be a surjective homomorphism $\mathcal{O}$-modules, and suppose that $\mathcal{P}$ is a free module. If the map $H^{0}(\mathcal{N}) \rightarrow H^{0}(\mathcal{P})$ of global sections is surjective, the sequence splits: $\mathcal{N}$ is isomorphic to the direct sum $\mathcal{M} \oplus \mathcal{P}$.
(iii) Let $0 \rightarrow \mathcal{M} \rightarrow \mathcal{N} \rightarrow \mathcal{P} \rightarrow 0$ be an exact sequence of $\mathcal{O}$-modules. If $\mathcal{P}$ is locally free, the dual modules form an exact sequence $0 \rightarrow \mathcal{P}^{*} \rightarrow \mathcal{N}^{*} \rightarrow \mathcal{M}^{*} \rightarrow 0$.
proof. (i) This follows from the analogous statement about modules over a ring.
(ii) Let $\left\{p_{i}\right\}$ be a basis of global sections of $\mathcal{P}$, let $v_{i}$ be global sections of $\mathcal{N}$ such that $g\left(v_{i}\right)=p_{i}$, and let $f$ be the map $f: \mathcal{P} \rightarrow \mathcal{N}$ defined by $f\left(p_{i}\right)=v_{i}$. Then $f g$ is the identity map on $\mathcal{P}$.
(iii) The sequence $0 \rightarrow \mathcal{P}^{*} \rightarrow \mathcal{N}^{*} \rightarrow \mathcal{M}^{*}$ is exact whether or not the modules are locally free 8.1.14(ii)). The zero on the right comes from the fact that, when $\mathcal{P}$ is locally free, it is free on some affine covering, so the given sequence splits locally.

### 8.2 Divisors

## (8.2.1) the divisor of a function

Let $f$ be a rational function on a smooth curve $Y$. The divisor of $f$ is

$$
\operatorname{div}(f)=\sum_{q \in Y} \mathrm{v}_{q}(f) q
$$

where, as usual, $\mathrm{v}_{q}$ denotes the valuation of $K$ that corresponds to the point $q$ of $Y$.
The divisor of $f$ is written here as a sum over all points $q$, but it becomes a finite sum when we disregard terms with coefficient zero, because $f$ has finitely many zeros and poles. The coefficients will be zero at all other points.

The map

$$
K^{\times} \rightarrow(\text { divisors })^{+}
$$

that sends a rational function to its divisor is a homomorphism from the multiplicative group $K^{\times}$of nonzero elements of $K$ to the additive group of divisors:

$$
\operatorname{div}(f g)=\operatorname{div}(f)+\operatorname{div}(g)
$$

As before, a rational function $f$ has a zero of order $r>0$ at $q$ if $\mathrm{v}_{q}(f)=r$, and it has a pole of order $r$ at $q$ if $\mathrm{v}_{q}(f)=-r$. Thus the divisor of $f$ is the difference of two effective divisors:

$$
\operatorname{div}(f)=\operatorname{zeros}(f)-\operatorname{poles}(f)
$$

A rational function $f$ is regular on $Y$ if and only if $\operatorname{div}(f)$ is effective - if and only if poles $(f)=0$.
The divisor of a rational function is a principal divisor, and two divisors $D$ and $E$ are linearly equivalent if their difference $D-E$ is a principal divisor. For instance, the divisors $z \operatorname{eros}(f)$ and poles $(f)$ of a rational function $f$ are linearly equivalent.
8.2.2. Lemma. Let $f$ be a rational function on a smooth curve $Y$. For all complex numbers $c$, the divisors of zeros of $f-c$, which are the level sets of $f$, are linearly equivalent.
proof. The functions $f-c$ have the same poles as $f$.

## (8.2.3) invertible modules

An invertible $\mathcal{O}$-module is a locally free module of rank one - a module that is isomorphic to the free module $\mathcal{O}$ in a neighborhood of any point. The tensor product $\mathcal{L} \otimes_{\mathcal{O}} \mathcal{M}$ of invertible modules is invertible. The dual $\mathcal{L}^{*}$ of an invertible module $\mathcal{L}$ is invertible.

Part (i) of the next lemma explains the adjective 'invertible'.
8.2.4. Lemma. Let $\mathcal{L}$ be an invertible $\mathcal{O}$-module.
(i) Let $\mathcal{L}^{*}$ be the dual module. The canonical map $\mathcal{L}^{*} \otimes_{\mathcal{O}} \mathcal{L} \rightarrow \mathcal{O}$ defined by $\gamma \otimes \alpha \mapsto \gamma(\alpha)$ is an isomorphism.
(ii) The map $\mathcal{O} \rightarrow \mathcal{O}(\mathcal{L}, \mathcal{L})$ that sends a regular function $\alpha$ to multiplication by $\alpha$ is an isomorphism.
(iii) Every nonzero homomorphism $\mathcal{L} \xrightarrow{\varphi} \mathcal{M}$ to a locally free module $\mathcal{M}$ is injective.
proof. (i),(ii) It is enough to verify these assertions in the case that $\mathcal{L}$ is free, isomorphic to $\mathcal{O}$, in which case they are clear.
(iii) The problem is local, so we may assume that the variety is affine, say $Y=\operatorname{Spec} A$, and that $\mathcal{L}$ and $\mathcal{M}$ are free. Then $\varphi$ becomes a nonzero homomorphism $A^{1} \rightarrow A^{k}$. Such a homomorphism is injective because $A$ is a domain.

## (8.2.5) the module $\mathcal{O}(D)$

To analyze functions with given poles on a smooth curve $Y$, we associate an invertible $\mathcal{O}$-module $\mathcal{O}(D)$ to a divisor $D$. The nonzero sections of $\mathcal{O}(D)$ on an open subset $V$ of $Y$ are the rational functions $f$ such that the the divisor $\operatorname{div}(f)+D$ is effective on $V$ - such that its restriction to $V$ is effective.

Thus, when $D$ is effetcive, the global sections of $\mathcal{O}(D)$ are the rational functions with poles bounded by $D$. They are the solutions of the classical problem that was mentioned at the beginning of hte chapter.

$$
\begin{equation*}
[\mathcal{O}(D)](V)=\{f \mid \operatorname{div}(f)+D \text { is effective on } V\} \cup\{0\} \tag{8.2.6}
\end{equation*}
$$

moduleOD

Points that aren't in $V$ impose no conditions on the sections of $\mathcal{O}(D)$ on $V$. A section on $V$ can have arbitrary zeros or poles at points not in $V$.

When $D$ is an effective divisor, a rational function $f$ is a global section of $\mathcal{O}(D)$ if $\operatorname{poles}(f) \leq D$.
Say that $D=\sum r_{i} q_{i}$. If $q_{i}$ is a point of an open set $V$ and if $r_{i}>0$, a section of $\mathcal{O}(D)$ on $V$ may have a pole of order at most $r_{i}$ at $q_{i}$, and if $r_{i}<0$ a section must have a zero of order at least $-r_{i}$ at $q_{i}$. For example, the module $\mathcal{O}(-q)$ is the maximal ideal $\mathfrak{m}_{q}$. The sections of $\mathcal{O}(-q)$ on an open set $V$ that contains $q$ are the regular functions on $V$ that are zero at $q$. Similarly, the sections of $\mathcal{O}(q)$ on an open set $V$ that contains $q$ are the rational functions that have a pole of order at most 1 at $q$ and are regular at every other point of $V$. The sections of $\mathcal{O}(-q)$ and of $\mathcal{O}(q)$ on an open set $V$ that doesn't contain $p$ are the regular functions on $V$.

The fact that a section of $\mathcal{O}(D)$ is allowed to have a pole at $q_{i}$ if $r_{i}>0$ contrasts with the divisor of a function. If $\operatorname{div}(f)=\sum r_{i} q_{i}$, then $r_{i}>0$ means that $f$ has a zero at $q_{i}$. Thus, if $\operatorname{div}(f)=D$, then $f$ will be a global section of $\mathcal{O}(-D)$.

## LisOD <br> 8.2.7. Proposition. Let $D$ and $E$ be divisors on a smooth curve $Y$.

(i) The map $\mathcal{O}(D) \otimes_{\mathcal{O}} \mathcal{O}(E) \rightarrow \mathcal{O}(D+E)$ that sends $f \otimes g$ to the product $f g$ is an isomorphism.
(ii) The dual module $\mathcal{O}(D)^{*}$ is $\mathcal{O}(-D)$.
(iii) $\mathcal{O}(D) \subset \mathcal{O}(E)$ if and only if $E-D$ is effective.
proof. We may assume that $Y$ is affine and that the supports of $D$ and $E$ contain at most one point: $D=r p$ and $E=s p$. We may also assume that the maximal ideal at $p$ is a principal ideal, generated by an element $x$. Then $\mathcal{O}(D), \mathcal{O}(-D), \mathcal{O}(E)$, and $\mathcal{O}(D+E)$ have bases $x^{r}, x^{-r}, x^{s}$ and $x^{r+s}$, respectively.
limODK 8.2.8. Corollary The union of the modules $\mathcal{O}(D)$ is the function field module $\mathcal{K}$.
idealOD 8.2.9. Corollary. Let $Y$ be a smooth curve, and let $\mathfrak{m}_{1}, \ldots, \mathfrak{m}_{k}$ be the maximal ideals at points $q_{1}, \ldots, q_{k}$ of $Y$. The product ideal $\mathcal{I}=\mathfrak{m}_{1}^{r_{1}} \cdots \mathfrak{m}_{k}^{r_{k}}$ of $\mathcal{O}_{Y}$ is equal to $\mathcal{O}_{Y}(-D)$, where $D$ is the effective divisor $\sum r_{i} q_{i}$. Thus the nonzero ideals of $\mathcal{O}_{Y}$ correspond bijectively to divisors $-D$, where $D$ is effective.

A module of the form $\mathcal{O}(D)$ is equal to $\mathcal{O}$ on the open complement of the support of $D$. Therefore, when we localize by forming the tensor product $\mathcal{O}(D) \otimes_{\mathcal{O}} \mathcal{K}$ with the function field module $\mathcal{K}$, we get $\mathcal{K}$. Since an invertible module $\mathcal{L}$ is isomorphic to $\mathcal{O}$ on an open set, the localization $\mathcal{L}_{\mathcal{K}}=\mathcal{L} \otimes_{\mathcal{O}} \mathcal{K}$ is a one dimensional $K$ vector space, but without a chosen basis (see (??)). As the next proposition shows, this is is the only difference between an invertible module $\mathcal{L}$ and a module of the form $\mathcal{O}(D)$.

LD Notation. Let $\mathcal{L}$ be an invertible module and let $D$ be a divisor. We denote by $\mathcal{L}(D)$ the invertible module $\mathcal{L} \otimes_{\mathcal{O}} \mathcal{O}(D)$. If $\mathcal{L} \approx \mathcal{O}(D)$, the degree of $\mathcal{L}$ is defined to be the degree of $D$.
proof of proposition 8.2.10 Since the function field module $\mathcal{K}$ of $Y$ is the union $\mathcal{K}=\bigcup \mathcal{O}(D)$, we also have $\mathcal{L}_{K}=\bigcup \mathcal{L}(D)$, where $\mathcal{L}(D)$ denotes the tensor product $\mathcal{L} \otimes_{\mathcal{O}} \mathcal{O}(D)$. A nonzero global section $\alpha$ of $\mathcal{L}_{K}$ will be a global section of $\mathcal{L}(D)$ for some $D$. It will define a map $\mathcal{O} \xrightarrow{\alpha} \mathcal{L}(D)$. Passing to duals, $\mathcal{L}(D)^{*}=\mathcal{L}^{*} \otimes_{\mathcal{O}} \mathcal{O}(D)^{*} \approx \mathcal{L}^{*} \otimes_{\mathcal{O}} \mathcal{O}(-D) \approx \mathcal{L}^{*}(-D)$. The dual of the map $\alpha$ is a nonzero and therefore injective map $\mathcal{L}^{*}(-D) \rightarrow \mathcal{O}$ whose image is an ideal of $\mathcal{O}$. So $\mathcal{L}^{*}(-D)$ is isomorphic to $\mathcal{O}(-E)$ for some effective divisor $E$ 8.2.9. Therefore $\mathcal{L}^{*}$ is isomorphic to $\mathcal{O}(D-E)$. Dualizing once more, $\mathcal{L}$ is isomorphic to $\mathcal{O}(E-D)$.
8.2.11. Proposition. Let $\mathcal{L} \subset \mathcal{M}$ be an inclusion of invertible $\mathcal{O}$-modules. Then $\mathcal{M}=\mathcal{L}(E)$ for some effective divisor $E$.
proof. Since $\mathcal{L} \subset \mathcal{M}, \mathcal{L} \otimes_{\mathcal{O}} \mathcal{M}^{*} \subset \mathcal{M} \otimes_{\mathcal{O}} \mathcal{M}^{*}=\mathcal{O}$. Therefore $\mathcal{L} \otimes_{\mathcal{O}} \mathcal{M}^{*}=\mathcal{O}(-E)$ for some effective divisor $E$ (refidealOD), and $\mathcal{L}=\mathcal{M}(-E)$.

It is important to note that, though every invertible module $\mathcal{L}$ is isomorphic to one of the form $\mathcal{O}(D)$, the divisor $D$ isn't uniquely determined by $\mathcal{M}$. If $D$ and $E$ are divisors, $\mathcal{O}(D)$ is a submodule of $\mathcal{O}(E)$ only when $E-D$ is effective. But as the next proposition explains, there may be injective homomorphisms from $\mathcal{O}(D)$ to $\mathcal{O}(E)$ that aren't inclusions.
8.2.12. Proposition. Let $D$ and $E$ be divisors on a smooth curve $Y$. Multiplication by a rational function $f$ such that $\operatorname{div}(f)+E-D \geq 0$ defines a homomorphism of $\mathcal{O}$-modules $\mathcal{O}(D) \rightarrow \mathcal{O}(E)$, and every homomorphism $\mathcal{O}(D) \rightarrow \mathcal{O}(E)$ is multiplication by such a function.
proof. For any $\mathcal{O}$-module $\mathcal{M}$, a homomorphism $\mathcal{O} \rightarrow \mathcal{M}$ is multiplication by a global section of $\mathcal{M}$ 6.4.5. Then a homomorphism $\mathcal{O} \rightarrow \mathcal{O}(E-D)$ will be multiplication by a rational function $f$ such that $\operatorname{div}(f)+E-D \geq$ 0 . If $f$ is such a function, one obtains a homomorphism $\mathcal{O}(D) \longrightarrow \mathcal{O}(E)$ by tensoring with $\mathcal{O}(D)$.

### 8.2.13. Corollary.

(i) The modules $\mathcal{O}(D)$ and $\mathcal{O}(E)$ are isomorphic if and only if the divisors $D$ and $E$ are linearly equivalent.
(ii) Let $f$ be a rational function on $Y$, and let $D=\operatorname{div}(f)$. Multiplication by $f$ defines an isomorphism $\mathcal{O}(D) \rightarrow \mathcal{O}$.

### 8.3 The Riemann-Roch Theorem

Let $Y$ be a smooth projective curve. In Chapter 7, we learned that, when $\mathcal{M}$ is a finite $\mathcal{O}_{Y}$-module, the cohomology $H^{q}(Y, \mathcal{M})$ is a finite-dimensional vector space for all $q$, and is zero if $q \neq 0,1$. As before, we denote the dimension of the space $H^{q}(Y, \mathcal{M})$ by $\mathbf{h}^{q} \mathcal{M}$ or, if there is ambiguity about the variety, by $\mathbf{h}^{q}(Y, \mathcal{M})$.

The Euler characteristic 7.6.5 of a finite $\mathcal{O}$-module $\mathcal{M}$ is

$$
\begin{equation*}
\chi(\mathcal{M})=\mathbf{h}^{0} \mathcal{M}-\mathbf{h}^{1} \mathcal{M} \tag{8.3.1}
\end{equation*}
$$

rrone

In particular,

$$
\chi\left(\mathcal{O}_{Y}\right)=\mathbf{h}^{0} \mathcal{O}_{Y}-\mathbf{h}^{1} \mathcal{O}_{Y}
$$

The dimension $\mathbf{h}^{1} \mathcal{O}_{Y}$ is the arithmetic genus of $Y$. The arithmetic genus is denoted by $p_{a}$. We will see below, in 8.3.10 (iv), that $\mathbf{h}^{0} \mathcal{O}_{Y}=1$. So

$$
\begin{equation*}
\chi\left(\mathcal{O}_{Y}\right)=1-p_{a} \tag{8.3.2}
\end{equation*}
$$

8.3.3. Riemann-Roch Theorem (version 1). Let $D=\sum r_{i} p_{i}$ be a divisor on a smooth projective curve $Y$. Then

$$
\chi(\mathcal{O}(D))=\chi(\mathcal{O})+\operatorname{deg} D \quad\left(=\operatorname{deg} D+1-p_{a}\right)
$$

proof. We analyze the effect on cohomology when a divisor is changed by adding or subtracting a point. We do this by inspecting the inclusion $\mathcal{O}(D-p) \subset \mathcal{O}(D)$. Let $\epsilon$ be the cokernel of the inclusion map, so that there is a short exact sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{O}(D-p) \rightarrow \mathcal{O}(D) \rightarrow \epsilon \rightarrow 0 \tag{8.3.4}
\end{equation*}
$$

in which $\epsilon$ is a one-dimensional vector space supported at $p$. This sequence can be obtained by tensoring the sequence

$$
\begin{equation*}
0 \rightarrow \mathfrak{m}_{p} \rightarrow \mathcal{O} \rightarrow \kappa_{p} \rightarrow 0 \tag{8.3.5}
\end{equation*}
$$

chicurvetwo

RRcurve
with the invertible module $\mathcal{O}(D)$, because $\mathfrak{m}_{p}$ is isomorphic to $\mathcal{O}(-p)$, t
Since $\epsilon$ is a one-dimensional module supported at $p, \mathbf{h}^{0} \epsilon=1$, and $\mathbf{h}^{1} \epsilon=0$. Let's denote the onedimensional vector space $H^{0}(Y, \epsilon)$ (which isn't very different from $\epsilon$ itself) by [1]. The cohomology sequence associated to 8.3 .4 is

$$
\begin{equation*}
0 \rightarrow H^{0}(Y, \mathcal{O}(D-p)) \rightarrow H^{0}(Y, \mathcal{O}(D)) \xrightarrow{\gamma}[1] \xrightarrow{\delta} H^{1}(Y, \mathcal{O}(D-p)) \rightarrow H^{1}(Y, \mathcal{O}(D)) \rightarrow 0 \tag{8.3.6}
\end{equation*}
$$

In this exact sequence, one of the two maps, $\gamma$ or $\delta$, must be zero. Either
(1) $\gamma$ is zero and $\delta$ is injective. In this case

$$
\mathbf{h}^{0} \mathcal{O}(D-p)=\mathbf{h}^{0} \mathcal{O}(D) \quad \text { and } \quad \mathbf{h}^{1} \mathcal{O}(D-p)=\mathbf{h}^{1} \mathcal{O}(D)+1, \quad \text { or }
$$

(2) $\delta$ is zero and $\gamma$ is surjective, in which case

$$
\left.\mathbf{h}^{0} \mathcal{O}(D)-p\right)=\mathbf{h}^{0} \mathcal{O}(D)-1 \quad \text { and } \quad \mathbf{h}^{1} \mathcal{O}(D-p)=\mathbf{h}^{1} \mathcal{O}(D)
$$

In either case,

$$
\begin{equation*}
\chi(\mathcal{O}(D))=\chi(\mathcal{O}(D-p))+1 \tag{8.3.7}
\end{equation*}
$$

Also, deg $D=\operatorname{deg}(D-p)+1$. The Riemann-Roch theorem follows, because wecan get from $\mathcal{O}$ to $\mathcal{O}(D)$ by a finite number of operations, each of which changes the divisor by adding or subtracting a point.

The Riemann-Roch Theorem can be written in terms of an invertible module. Recall that the degree of an invertible module $\mathcal{L}$ is the degree of a divisor $D$ such that $\mathcal{L} \approx \mathcal{O}(D)$.

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onepole

RRcor
8.3.8. Corollary. Let $\mathcal{L}$ be an invertible $\mathcal{O}$-module on a smooth projective curve $Y$. Then $\chi(\mathcal{L})=\operatorname{deg} \mathcal{L}+$ $1-p_{a}$.

Because $\mathbf{h}^{0} \geq \mathbf{h}^{0}-\mathbf{h}^{1}=\chi$, this version of the Riemann-Roch Theorem gives reasonably good control of $H^{0}$. It is less useful for controlling $H^{1}$. To do that, one wants the full Riemann-Roch Theorem, which is in Section 8.7 It requires some preparation. However, the version which is above has important consequences:
8.3.9. Corollary. Let p be a point of a smooth projective curve $Y$. The dimension $\mathbf{h}^{0}(Y, \mathcal{O}(n p))$ tends to infinity with $n$. Therefore there exist rational functions with a pole of large order at $p$ and no other poles.
proof. When we change the module $\mathcal{O}(n p)$ to $\mathcal{O}((n+1) p)$, then, as we saw above, either $\mathbf{h}^{0}$ increases or $\mathbf{h}^{1}$ decreases. Since $H^{1}(Y, \mathcal{O}(n p))$ is finite-dimensional, the second possibility can occur only finitely many times.
8.3.10. Corollary. Let $Y$ be a smooth projective curve.
(i) The divisor of a rational function on $Y$ has degree zero: The number of zeros is equal to the number of poles.
(ii) Linearly equivalent divisors on $Y$ have equal degrees.
(iii) A nonconstant rational function on $Y$ takes every value, including infinity, the same number of times.
(iv) A rational function on $Y$ that is regular at every point of $Y$ is a constant: $H^{0}(Y, \mathcal{O})=\mathbb{C}$.
(v) Let $D$ be a divisor on $Y$. If $\operatorname{deg} D \geq p_{a}$, then $\mathbf{h}^{0} \mathcal{O}(D)>0$. If $\mathbf{h}^{0} \mathcal{O}(D)>0$, then $\operatorname{deg} D \geq 0$.
proof. (i) Let $f$ be a nonzero rational function and let $D=\operatorname{div}(f)$. Multiplication by $f$ defines an isomorphism $\mathcal{O}(D) \rightarrow \mathcal{O}$, so $\chi(\mathcal{O}(D))=\chi(\mathcal{O})$. On the other hand, by Riemann-Roch, $\chi(\mathcal{O}(D))=\chi(\mathcal{O})+\operatorname{deg} D$. Therefore $\operatorname{deg} D=0$.
(ii) If two divisors $D$ and $E$ are linearly equivalent, say $D-E=\operatorname{div}(f)$, then $D-E$ has degree zero, and $\operatorname{deg} D=\operatorname{deg} E$.
(iii) The zeros of the functions $f-c$ are linearly equivalent to the poles of $f$.
(iv) According to (iii), a nonconstant function must have a pole.
(v) If $\operatorname{deg} D \geq p_{a}$, then $\chi=\operatorname{deg} D+1-p_{a} \geq 1$, and $\mathbf{h}^{0} \geq \mathbf{h}^{0}-\mathbf{h}^{1}=\chi$. If $\mathcal{O}(D)$ has a nonzero global section $f$, a rational function such that $\operatorname{div}(f)+D=E$ is effective, then $\operatorname{deg} E \geq 0$. Since the degree of $\operatorname{div}(f)$ is zero, $\operatorname{deg} D \geq 0$.
8.3.11. Theorem. With its classical topology, a smooth projective curve $Y$ is a connected, compact, orientable two-dimensional manifold.
proof. We prove connectedness here. All other points have been discussed before (Theorem 1.7.24). A nonempty topological space is connected if it isn't the union of two disjoint, nonempty, closed subsets. Suppose that, in the classical topology, $Y$ is the union of disjoint, nonempty closed subsets $Y_{1}$ and $Y_{2}$. Both $Y_{1}$ and $Y_{2}$ will be compact manifolds. Let $q$ be a point of of $Y_{1}$. Corollary 8.3.9 shows that there is a nonconstant rational function $f$ whose only pole is at $q$. Then $f$ will be a regular function on the complement of $q$. It will be an analytic function on the compact manifold $Y_{2}$.

For review: For any point $q$ of $Y, Y$ can be represented as a branched covering of the projective line $\mathbb{P}^{1}$ that is unramified at $q$. Then a small neighborhood $V$ of $q$ in $Y$ maps bijectively to an neighborhood $U$ of its image $p$ in $\mathbb{P}^{1}$. To say that $f$ is analytic means that the function on $V$ that corresponds to $f$ is an analytic function of one variable on $U$.

The maximum principle for analytic functions asserts that a nonconstant analytic function on an open region of the complex plane has no maximal absolute value in the region. This applies to the neighborhood $U$ of $p$, and also to the neighborhood $V$ of $q$. Consequently, if $f$ isn't constant, $|f(q)|$ cannot be a maximum value of $f$. Since $q$ is arbitrary, $f$ cannot have a maximum on $Y_{2}$. On the other hand, since $Y_{2}$ is compact, a continuous function does have a maximum. So $f$ must be constant on $Y_{2}$. When we subtract that constant from $f$, we obtain a nonconstant rational function that is zero on $Y_{2}$. But since $Y$ has dimension 1, the zero locus of a rational function is finite. This is a contradiction.

### 8.4 The Birkhoff-Grothendieck Theorem

This theorem describes finite, torsion-free modules on the projective line.
8.4.1. Birkhoff-Grothendieck Theorem. A finite, torsion-free $\mathcal{O}$-module on the projective line $\mathbb{P}^{1}$ is isomorphic to a direct sum of twisting modules: $\mathcal{M} \approx \bigoplus \mathcal{O}\left(n_{i}\right)$.

This theorem was proved by Grothendieck in 1957 using cohomology. It had been proved by Birkhoff in 1909, much earlier, in the following equivalent form:

Birkhoff Factorization Theorem. Let $A_{0}=\mathbb{C}[u], A_{1}=\mathbb{C}\left[u^{-1}\right]$, and $A_{01}=\mathbb{C}\left[u, u^{-1}\right]$. Let $P$ be an invertible $A_{01}$-matrix. There exist an invertible $A_{0}$-matrix $Q_{0}$ and an invertible $A_{1}$-matrix $Q_{1}$ such that $Q^{-1} P Q_{1}$ is diagonal, and the diagonal entries are integer powers of $u$.

We recall the cohomology of the twisting modules on $\mathbb{P}^{1}$ : Let $r$ be a positive integer. According to Theorem 7.5.4.
$\mathbf{h}^{0} \mathcal{O}(r)=r+1, \quad \mathbf{h}^{1} \mathcal{O}(r)=0, \quad \mathbf{h}^{0} \mathcal{O}(-r)=0, \quad$ and $\quad \mathbf{h}^{1} \mathcal{O}(-r)=r-1$.
8.4.2. Lemma. Let $X$ denote the projective line, and let $\mathcal{M}$ be a finite, torsion-free $\mathcal{O}$-module on $X$. For sufficiently large $r$,
(i) the only module homomorphism $\mathcal{O}(r) \rightarrow \mathcal{M}$ is the zero map, and
(ii) $\mathbf{h}^{0}(X, \mathcal{M}(-r))=0$.
proof. (i) A nonzero map from the invertible module $\mathcal{O}(r)$ to the locally free module $\mathcal{M}$ will be injective 8.2.4, and the associated map $H^{0}(X, \mathcal{O}(r)) \rightarrow H^{0}(X, \mathcal{M})$ will also be injective. So $\mathbf{h}^{0}(X, \mathcal{O}(r)) \leq$ $\mathbf{h}^{0}(X, \mathcal{M})$. Since $\mathbf{h}^{0}(X, \mathcal{O}(r))=r+1$ and since $\mathbf{h}^{0}(X, \mathcal{M})$ is finite, $r$ is bounded.
(ii) A global section of $\mathcal{M}(-r)$ defines a map $\mathcal{O} \rightarrow \mathcal{M}(-r)$. Its twist by $r$ will be a map $\mathcal{O}(r) \rightarrow \mathcal{M}$.
proof of the Birkhoff-Grothendieck Theorem. This is Grothendieck's proof.
Lemma 8.1.1 tells us that $\mathcal{M}$ is locally free. We use induction on the rank of $\mathcal{M}$. We suppose that $\mathcal{M}$ has rank $r>0$, and that the theorem has been proved for locally free $\mathcal{O}$-modules of rank less than $r$. The plan is to show that $\mathcal{M}$ has a twisting module as a direct summand, so that $\mathcal{M}=\mathcal{W} \oplus \mathcal{O}(n)$ for some $\mathcal{W}$. Then induction on the rank can be applied to $\mathcal{W}$.

Since twisting is compatible with direct sums, we may replace $\mathcal{M}$ by a twist $\mathcal{M}(n)$. Instead of showing that $\mathcal{M}$ has a twisting module $\mathcal{O}(n)$ as a direct summand, we show that, after we replace $\mathcal{M}$ by a suitable twist, the structure sheaf $\mathcal{O}$ will be a direct summand.

As we know 6.7.21, the twist $\mathcal{M}(n)$ will have a nonzero global section when $n$ is sufficiently large, and it will have no nonzero global section when $n$ is sufficiently negative (Lemma8.4.2(ii)). Therefore, when we replace $\mathcal{M}$ by a suitable twist, we will have $H^{0}(X, \mathcal{M}) \neq 0$ but $H^{0}(X, \mathcal{M}(-1))=0$. We assume that this is true for $\mathcal{M}$.

We choose a nonzero global section $s$ of $\mathcal{M}$ and consider the injective multiplication map $\mathcal{O} \xrightarrow{s} \mathcal{M}$. Let $\mathcal{W}$ be its cokernel, so that we have a short exact sequence

$$
\begin{equation*}
0 \rightarrow \mathcal{O} \xrightarrow{s} \mathcal{M} \rightarrow \mathcal{W} \rightarrow 0 \tag{8.4.3}
\end{equation*}
$$

8.4.4. Lemma. Let $\mathcal{W}$ be the $\mathcal{O}$-module that appears in the sequence 8.4.3.
(i) $H^{0}(X, \mathcal{W}(-1))=0$.
(ii) $\mathcal{W}$ is torsion-free, and therefore locally free.
(iii) $\mathcal{W}$ is a direct sum $\bigoplus_{i=1}^{r-1} \mathcal{O}\left(n_{i}\right)$ of twisting modules on $\mathbb{P}^{1}$, with $n_{i} \leq 0$.
proof. (i) This follows from the cohomology sequence associated to the twisted sequence

$$
0 \rightarrow \mathcal{O}(-1) \rightarrow \mathcal{M}(-1) \rightarrow \mathcal{W}(-1) \rightarrow 0
$$

because $H^{0}(X, \mathcal{M}(-1))=0$ and $H^{1}(X, \mathcal{O}(-1))=0$.
(ii) If $\mathcal{W}$ had a nonzero torsion submodule, so would $\mathcal{W}(-1)$, and then $\mathcal{W}(-1)$ would have a nonzero global section 8.1.1.
(iii) The fact that $\mathcal{W}$ is a direct sum of twisting modules follows by induction on the rank: $\mathcal{W} \approx \bigoplus \mathcal{O}\left(n_{i}\right)$. Since $H^{0}(X, \mathcal{W}(-1))=0$, we must have $H^{0}\left(X, \mathcal{O}\left(n_{i}-1\right)\right)=0$ too. Therefore $n_{i}-1<0$, and $n_{i} \leq 0$.

We go back to the proof of Theorem 8.4.1. Lemma 8.1 .20 tells us that the dual of the sequence 8.4.3) is an exact sequence

$$
0 \rightarrow \mathcal{W}^{*} \longrightarrow \mathcal{M}^{*} \longrightarrow \mathcal{O}^{*} \rightarrow 0
$$

and $\mathcal{W}^{*} \approx \bigoplus \mathcal{O}\left(-n_{i}\right)$ with $-n_{i} \geq 0$. Therefore $\mathbf{h}^{1} \mathcal{W}^{*}=0$. The map $H^{0}(\mathcal{M}) \rightarrow H^{0}\left(\mathcal{O}^{*}\right)$ is surjective. Lemma 8.1.20 tells us that $\mathcal{M}^{*}$ is isomorphic to $\mathcal{W}^{*} \oplus \mathcal{O}^{*}$. Then $\mathcal{M}$ is isomorphic to $\mathcal{W} \oplus \mathcal{O}$.

### 8.5 Differentials

Why differentials enter into the Riemann-Roch Theorem is a mystery, but they do, so we introduce them here.
Let $A$ be an algebra and let $M$ be an $A$-module. A derivation $A \xrightarrow{\delta} M$ is a $\mathbb{C}$-linear map that satisfies the product rule for differentiation - a map that has these properties:

$$
\begin{equation*}
\delta(a b)=a \delta b+b \delta a, \quad \delta(a+b)=\delta a+\delta b, \quad \text { and } \quad \delta c=0 \tag{8.5.1}
\end{equation*}
$$

for all $a, b$ in $A$ and all $c$ in $\mathbb{C}$. The fact that $\delta$ is $\mathbb{C}$-linear, i.e., that it is a homomorphism of vector spaces, follows. Since $d c=0, \delta(c b)=c \delta b$. For example, differentiation $\frac{d}{d t}$ is a derivation $\mathbb{C}[t] \rightarrow \mathbb{C}[t]$.
8.5.2. Lemma. Let $A \xrightarrow{\varphi} B$ be an algbra homomorphism, and let $M \xrightarrow{f} N$ be a homomorphism of $B$-modules.
(i) Let $B \xrightarrow{\delta} M$ be a derivation. The composed maps $A \xrightarrow{\delta \varphi} M$ and $B \xrightarrow{f \delta} N$ are derivations.
(ii) Suppose that the homomorphism $\varphi$ is surjective. Let $B \xrightarrow{g} M$ be a map, and let $d=g \circ \varphi$. If $A \xrightarrow{d} M$ is a derivation, then $g$ is a derivation.

The module of differentials $\Omega_{A}$ of an algebra $A$ is an $A$-module that is generated by elements denoted by $d a$, one for each element $a$ of $A$. Its elements are (finite) combinations $\sum b_{i} d a_{i}$, with $a_{i}$ and $b_{i}$ in $A$. The defining relations among the generators $d a$ are the ones that make the map $A \xrightarrow{d} \Omega_{A}$ that sends $a$ to $d a$ a derivation. For all $a, b$ in $A$ and all $c$ in $\mathbb{C}$,

$$
\begin{equation*}
d(a b)=a d b+b d a, \quad d(a+b)=d a+d b, \quad \text { and } \quad d c=0 \tag{8.5.3}
\end{equation*}
$$

The elements of $\Omega_{A}$ are called differentials.

### 8.5.4. Lemma.

(i) When we compose a homomorphism $\Omega_{A} \xrightarrow{\varphi} M$ of $\mathcal{O}$-modules with the derivation $A \xrightarrow{d} \Omega_{A}$, we obtain a derivation $A \xrightarrow{\varphi \circ d} M$. Composition with $d$ defines a bijection between homomorphisms $\Omega_{A} \rightarrow M$ and derivations $A \xrightarrow{\delta} M$.
(ii) $\Omega$ is a functor: An algebra homomorphism $A \xrightarrow{u} B$ induces a homomorphism $\Omega_{A} \xrightarrow{v} \Omega_{B}$ that is compatible with the ring homomorphism $u$, and that makes a diagram


By compatibility of $v$ with $u$ we mean that, if $\omega$ is an element of $\Omega_{A}$ and $\alpha$ is in $A$, then $v(\alpha \omega)=u(\alpha) v(\omega)$. proof. (i) The composition $\delta=\varphi \circ d$ is a derivation $A \rightarrow M$. In the other direction, given a derivation $A \xrightarrow{\delta} M$, we define a map $\Omega_{A} \xrightarrow{\varphi} M$ by $\varphi(d a)=\delta(a)$. It follows from the defining relations for $\Omega_{A}$ that $\varphi$ is a well-defined homomorphism of $A$-modules.
(ii) When $\Omega_{B}$ is made into an $A$-module by restriction of scalars, the composed map $A \xrightarrow{u} B \xrightarrow{d} \Omega_{B}$ will be a derivation to which (i) applies.
8.5.5. Lemma. Let $R$ be the polynomial ring $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$. The $R$-module of differentials $\Omega_{R}$ is a free module with basis $d x_{1}, \ldots, d x_{n}$.
proof. The formula $d f=\sum \frac{d f}{d x_{i}} d x_{i}$ follows from the defining relations. It shows that the elements $d x_{1}, \ldots, d x_{n}$ generate the $R$-module $\Omega_{R}$. Let $v_{1}, \ldots, v_{n}$ be a basis of a free $R$-module $V$. The product rule for derivatives shows that the map $\delta: R \rightarrow V$ by $\delta(f)=\frac{\partial f}{\partial x_{i}} v_{i}$ is a derivation. It induces a module homomorphism $\Omega_{A} \rightarrow V$ that sends $d x_{i}$ to $v_{i}$. Since $d x_{1}, \ldots, d x_{n}$ generate $\Omega_{R}$ and since $v_{1}, \ldots, v_{n}$ is a basis of $V, \varphi$ is an isomorphism.
8.5.6. Proposition. Let I be an ideal of an algebra $R$, let $A$ be the quotient algebra $R / I$, and let dI denote the set of differentials df with $f$ in $I$. The subset $N=d I+I \Omega_{R}$ is a submodule of $\Omega_{R}$, and $\Omega_{A}$ is isomorphic to the quotient module $\Omega_{R} / N$.

The proposition can be interpreted this way: Suppose that the ideal $I$ is generated by elements $f_{1}, \ldots, f_{r}$ of $R$. Then $\Omega_{A}$ is the quotient of $\Omega_{R}$ that is obtained from $\Omega_{R}$ by introducing these two rules:

- $d f_{i}=0$, and
- multiplication by $f_{i}$ is zero.

For example, let $A$ be the quotient $\mathbb{C}[y] /\left(y^{2}\right)$ of a polynomial ring $R$ in one variable. Here $I$ is the principal ideal $\left(y^{2}\right), 2 y d y$ generates $d I$, and $y^{2} d y$ generates $I \Omega_{A}$. So $y d y$ generates the $R$-module $N$. Then, if $\bar{y}$ denotes the residue of $y$ in $A, \Omega_{A}$ is generated by an element $d \bar{y}$, with the relation $2 \bar{y} d \bar{y}=0$. In particular, it isn't the zero module.
proof of Proposition 8.5.6. First, $I \Omega_{R}$ is a submodule of $\Omega_{R}$, and $d I$ is an additive subgroup of $\Omega_{R}$. To show that $N$ is a submodule, we must show that scalar multiplication by an element of $R$ maps $d I$ to $N$, i.e., that if $g$ is in $R$ and $f$ is in $I$, then $g d f$ is in $N$. By the product rule, $g d f=d(f g)-f d g$. Since $I$ is an ideal, $f g$ is in $I$. Then $d(f g)$ is in $d I$, and $f d g$ is in $I \Omega_{R}$. So $g d f$ is in $N$.

The two rules shown above hold in $\Omega_{A}$ because the generators $f_{i}$ of $I$ are zero in $A$. Therefore $N$ is in the kernel of the surjective map $\Omega_{R} \xrightarrow{v} \Omega_{A}$ defined by the homomorphism $R \rightarrow A$. Let $\bar{\Omega}$ denote the quotient module $\Omega_{R} / N$. This is an $A$-module, and because $N \subset$ ker $v, v$ defines a surjective map of $A$-modules $\bar{\Omega} \xrightarrow{\bar{v}} \Omega_{A}$. We show that $\bar{v}$ is bijective. Let $r$ be an element of $R$, let $a$ be its image in $A$, and let $\overline{d r}$ be its image in $\bar{\Omega}$. The composed map $R \xrightarrow{d} \Omega_{R} \xrightarrow{\pi} \bar{\Omega}$ is a derivation that sends $r$ to $\overline{d r}$, and because $I$ is in its kernel, it defines a derivation $R / I=A \xrightarrow{\delta} \bar{\Omega}$ that sends $a$ to $\overline{d r}$. This derivation corresponds to a homomorphism of $A$-modules $\Omega_{A} \rightarrow \bar{\Omega}$ that sends $d a$ to $\overline{d r}$, and that inverts $\bar{v}$ 8.5.4.
8.5.7. Corollary. If $A$ is a finite-type algebra, then $\Omega_{A}$ is a finite $A$-module.

This follows from Proposition 8.5 .6 because the module of differentials on the polynomial ring $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ is a finite module.
8.5.8. Lemma. Let $S$ be a multiplicative system in a domain $A$, and let $S^{-1} \Omega_{A}$ be the module of fractions of $\Omega_{A}$. The modules $S^{-1} \Omega_{A}$ and $\Omega_{S^{-1} A}$ are canonically isomorphic. In particular, if $K$ is the field of fractions of $A$, then $K \otimes_{A} \Omega_{A} \approx \Omega_{K}$.
We have moved the symbol $S^{-1}$ to the left for clarity.
proof of Lemma 8.5.8 The composition $A \rightarrow S^{-1} A \xrightarrow{d} \Omega_{S^{-1} A}$ is a derivation that defines an $A$-module homomorphism $\Omega_{A} \rightarrow \Omega_{S^{-1} A}$. This homomorphism extends to an $S^{-1} A$-homomorphism $S^{-1} \Omega_{A} \xrightarrow{\varphi} \Omega_{S^{-1} A}$ because scalar multiplication by the elements of $S$ is invertible in $\Omega_{S^{-1} A}$. The relation $d s^{-k}=-k s^{k-1} d s$ follows from the definition of a differential, and it shows that $\varphi$ is surjective. The quotient rule

$$
\delta\left(s^{-k} a\right)=-k s^{-k-1} a d s+s^{-k} d a
$$

can be used to define a derivation $S^{-1} A \xrightarrow{\delta} S^{-1} \Omega_{A}$, which corresponds to a homomorphism $\Omega_{S^{-1} A} \rightarrow$ $S^{-1} \Omega_{A}$ that inverts $\varphi$. Here, one must show that $\delta$ is well-defined, that $\delta\left(s_{1}^{-k} a_{1}\right)=\delta\left(s_{2}^{-\ell} a_{2}\right)$ if $s_{1}^{-\ell} a_{1}=$ $s_{2}^{-k} a_{2}$, and that $\delta$ is a derivation. You will be able to do this.

Lemma 8.5 .8 shows that a finite $\mathcal{O}$-module $\Omega_{Y}$ of differentials on a variety $Y$ is defined. When $U=$ $\operatorname{Spec} A$ is an affine open subset of $Y, \Omega_{Y}(U)=\Omega_{A}$.
8.5.9. Proposition. The module $\Omega_{Y}$ of differentials on a smooth curve $Y$ is invertible. If y is a local generator for the maximal ideal at a point $q$, then in a suitable neighborhood of $q, \Omega_{Y}$ is a free $\mathcal{O}$-module with basis $d y$.
proof. We may assume that $Y$ is affine, say $Y=\operatorname{Spec} B$. Let $q$ be a point of $Y$, and let $y$ be an element of $B$ with $\mathrm{v}_{q}(y)=1$. To show that $d y$ generates $\Omega_{B}$ locally, we may localize, so we may suppose that $y$ generates the maximal ideal $\mathfrak{m}$ at $q$. We must show that after we localize $B$ once more, every differential $d f$ with $f$ in $B$ will be a multiple of $d y$. Let $c$ be the value of the function $f$ at $q$ : Then $f=c+y g$ for some $g$ in $B$, and because $d c=0, d f=g d y+y d g$. Here $g d y$ is in $B d y$ and $y d y$ is in $\mathfrak{m} \Omega_{B}$, so

$$
\Omega_{B}=B d y+\mathfrak{m} \Omega_{B}
$$

If $W$ denotes the quotient module $\Omega_{B} /(B d y)$, then $W=\mathfrak{m} W$. The Nakayama Lemma tells us that there is an element $z$ in $\mathfrak{m}$ such that $s=1-z$ annihilates $W$. When we replace $B$ by its localization $B_{s}$, we will have $W=0$ and $\Omega_{B}=B d y$, as required.

We must still verify that $d y$ isn't a torsion element. If it were, say $b d y=0$, then because $d y$ is a local generator, $\Omega_{B}$ would be the zero module except at the finite set of zeros of $b$. Since the chosen point $q$ of $Y$ was arbitrary, it suffices to show that the local generator $d y$ for $\Omega_{B}$ isn't equal to zero. Let $R=\mathbb{C}[y]$ and $A=\mathbb{C}[y] /\left(y^{2}\right)$. The module $\Omega_{R}$ is free, with basis $d y$, and as noted above, if $\bar{y}$ is the residue of $y$ in $A$, the $A$-module $\Omega_{A}$ is generated by $d \bar{y}$, with the relation $2 \bar{y} d \bar{y}=0$. It isn't the zero module. Proposition 5.3.7tells us that, at our point $q$, the algebra $B / \mathfrak{m}_{q}^{2}$ is isomorphic to $A$, and Proposition 8.5.6 tells us that $\Omega_{A}$ is a quotient of $\Omega_{B}$. Since $\Omega_{A}$ isn't zero, neither is $\Omega_{B}$.

### 8.6 Trace

## (8.6.1) trace of a function

Let $Y \xrightarrow{\pi} X$ be a branched covering of smooth curves, and let $F$ and $K$ be the function fields of $X$ and $Y$, respectively.

The trace map $K \xrightarrow{\operatorname{tr}} F$ for a field extension of finite degree has been defined before 4.3.10. If $\alpha$ is an element of $K$, multiplication by $\alpha$ on the $F$-vector space $K$ is an $F$-linear operator, and $\operatorname{tr}(k)$ is the trace of that operator. The trace is $F$-linear: If $f_{i}$ are in $F$ and $\alpha_{i}$ are in $K$, then $\operatorname{tr}\left(\sum f_{i} \alpha_{i}\right)=\sum f_{i} \operatorname{tr}\left(\alpha_{i}\right)$. Moreover, the trace carries regular functions to regular functions: If $X^{\prime}=\operatorname{Spec} A^{\prime}$ is an affine open subset of $X$, with inverse image $Y^{\prime}=\operatorname{Spec} B^{\prime}$, then because $A^{\prime}$ is a normal algebra, the trace of an element of $B^{\prime}$ will be in $A^{\prime}$ 4.3.6. Using our abbreviated notation $\mathcal{O}_{Y}$ for $\pi_{*} \mathcal{O}_{Y}$, the trace defines a homomorphism of $\mathcal{O}_{X}$-modules

$$
\begin{equation*}
\mathcal{O}_{Y} \xrightarrow{\operatorname{tr}} \mathcal{O}_{X} \tag{8.6.2}
\end{equation*}
$$

Analytically, the trace can be described as a sum over the sheets of the covering. Let $n=[Y: X]$. Over a point $p$ of $X$ that isn't a branch point, there will be $n$ points $q_{1}, \ldots, q_{n}$ of $Y$. If $U$ is a small neighborhood of $p$ in $X$ in the classical topology, its inverse image $V$ will consist of disjoint neighborhoods $V_{i}$ of $q_{i}$, each of which maps bijectively to $U$. On $V_{i}$, the ring $\mathcal{B}$ of analytic functions will be isomorphic to the ring $\mathcal{A}$ of analytic functions on $U$. So $\mathcal{B}$ is isomorphic to the direct sum $\mathcal{A}_{1} \oplus \cdots \oplus \mathcal{A}_{n}$ of $n$ copies of $\mathcal{A}$. If a rational function $g$ on $Y$ is regular on $V$, its restriction to $V$ can be written as $g=g_{1} \oplus \cdots \oplus g_{n}$, with $g_{i}$ in $\mathcal{A}_{i}$. The matrix of left multiplication by $g$ on $\mathcal{A}_{1} \oplus \cdots \oplus \mathcal{A}_{n}$ is the diagonal matrix with entries $g_{i}$, so

$$
\begin{equation*}
\operatorname{tr}(g)=g_{1}+\cdots+g_{n} \tag{8.6.3}
\end{equation*}
$$

8.6.4. Lemma. Let $Y \xrightarrow{\pi} X$ be a branched covering of smooth curves, let $p$ be a point of $X$, let $q_{1}, \ldots, q_{k}$ be the fibre over $p$, and let $e_{i}$ be the ramification index at $q_{i}$. If a rational function $g$ on $Y$ is regular at the points $q_{1}, \ldots, q_{k}$, its trace is regular at $p$, and its value at $p$ is $[\operatorname{tr}(g)](p)=e_{1} g\left(q_{1}\right)+\cdots+e_{k} g\left(q_{k}\right)$.
proof. The regularity was discussed above. If $p$ isn't a branch point, we will have $k=n$ and $e_{i}=1$ for all $i$. In this case, the lemma follows by evaluating (8.6.3). It follows by continuity for any point $p$. As a point $p^{\prime}$ approaches $p, e_{i}$ points $q^{\prime}$ of $Y$ approach $q_{i}$ 8.1.8). For each such point, the limit of $g\left(q^{\prime}\right)$ will be $g\left(q_{i}\right)$.
traced

The structure sheaf is naturally contravariant. A branched covering $Y \xrightarrow{\pi} X$ corresponds to an $\mathcal{O}_{X^{-}}$ module homomorphism $\mathcal{O}_{X} \rightarrow \mathcal{O}_{Y}$. The trace map for functions is a homomorphism of $\mathcal{O}_{X}$-modules in the opposite direction: $\mathcal{O}_{Y} \xrightarrow{\mathrm{tr}} \mathcal{O}_{X}$.

Differentials are also naturally contravariant. A morphism $Y \rightarrow X$ induces an $\mathcal{O}_{X}$-module homomorphism $\Omega_{X} \rightarrow \Omega_{Y}$ that sends a differential $d x$ on $X$ to a differential on $Y$ that we denote by $d x$ too (8.5.4) (ii). As is true for functions, there is a trace map for differentials in the opposite direction. It is defined below, in 8.6.7, and will be denoted by $\tau$ :

$$
\Omega_{Y} \xrightarrow{\tau} \Omega_{X}
$$

First, a lemma about the natural contravariant map $\Omega_{X} \rightarrow \Omega_{Y}$ :
8.6.6. Lemma. (i) Let $p$ be the image in $X$ of a point $q$ of $Y$, let $x$ and $y$ be local generators for the maximal ideals of $X$ and $Y$ at $p$ and $q$, respectively, and let e be the ramification index at $q$. Then $d x=v y^{e-1} d y$, where $v$ is a local unit at $q$.
(ii) The canonical homomorphism $\Omega_{X} \rightarrow \Omega_{Y}$ is injective.
proof. (i) As we have noted before, $x=u y^{e}$, where $u$ is a local unit. Since $d y$ generates $\Omega_{Y}$ locally, there is a rational function $z$ that is regular at $q$, such that $d u=z d y$. Then $d x=d\left(u y^{e}\right)=y^{e} z d y+e y^{e-1} u d y=$ $v y^{e-1} d y$, where $v=y z+e u$. Since $y z$ is zero at $q$ and $e u$ is a local unit there, $v$ is a local unit.
(ii) See 8.2.4.

To define the trace for differentials, we begin with differentials of the functions fields. Let $F$ and $K$ be the function fields of $X$ and $Y$, respectively. Because the $\mathcal{O}_{Y}$-module $\Omega_{Y}$ is invertible, the module $\Omega_{K}$ of $K$-differentials, which is the localization $\Omega_{Y} \otimes_{\mathcal{O}} K$, is a free $K$-module of rank one. Any nonzero differential will form a $K$-basis. We choose as basis a nonzero $F$-differential $\alpha$. Its image in $\Omega_{K}$, which we denote by $\alpha$ too, will be a $K$-basis for $\Omega_{K}$. We can, for example, take $\alpha=d x$, where $x$ is a local coordinate function on $X$.

Since $\alpha$ is a basis, an element $\beta$ of $\Omega_{K}$ can be written uniquely, as

$$
\beta=g \alpha
$$

where $g$ is an element of $K$. The trace $\Omega_{K} \xrightarrow{\tau} \Omega_{F}$ is defined by

$$
\begin{equation*}
\tau(\beta)=\operatorname{tr}(g) \alpha \tag{8.6.7}
\end{equation*}
$$

deftrdif
where $\operatorname{tr}(g)$ is the trace of the function $g$. Since the trace for functions is $F$-linear, $\tau$ is also an $F$-linear map.
We need to check that $\tau$ is independent of the choice of $\alpha$. If $\alpha^{\prime}$ is another nonzero $F$-differential, then $f \alpha^{\prime}=\alpha$ for some nonzero element $f$ of $F$, and $g \alpha=g f \alpha^{\prime}$. Since $\operatorname{tr}$ is $F$-linear, $\operatorname{tr}(g f)=f \operatorname{tr} g$, and

$$
\operatorname{tr}(g f) \alpha^{\prime}=\operatorname{tr}(g) f \alpha^{\prime}=\operatorname{tr}(g) \alpha
$$

Using $\alpha^{\prime}$ in place of $\alpha$ gives the same value for the trace.
A differenial of the function field $K$ will be called a rational differential. A rational differential $\beta$ is regular at a point $q$ of $Y$ if there is an affine open neighborhood $Y^{\prime}=\operatorname{Spec} B$ of $q$ such that $\beta$ is an element of $\Omega_{B}$. If $y$ is a local generator for the maximal ideal $\mathfrak{m}_{q}$ and $\beta=g d y$, the differential $\beta$ is regular at $q$ if the rational function $g$ is regular there.
at $q$ of a regular
Let $p$ be a point of $X$. Working locally at $p$, we may suppose that $X$ and $Y$ are affine, $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$, that the maximal ideal at $p$ is a principal ideal, generated by an element $x$ of $A$, and that the differential $d x$ generates $\Omega_{A}$. Let $q_{1}, \ldots, q_{k}$ be the points of $Y$ that lie over $p$, and let $e_{i}$ be the ramification index at $q_{i}$.
8.6.8. Corollary. With notation as above,
(i) When viewed as a differential on $Y$, dx has zeros of orders $e_{i}-1$ at $q_{i}$.
(ii) If a differential $\beta$ on $Y$ is regular at the points $q_{i}, \ldots, q_{k}$, then $\beta=g d x$, where $g$ is a rational function with poles of orders at most $e_{i}-1$ at $q_{i}$.

This follows from Lemma 8.6.6(i).
8.6.9. Main Lemma. Let $Y \xrightarrow{\pi} X$ be a branched covering, let $p$ be a point of $X$, and let $q_{1}, \ldots, q_{k}$ be the points of $Y$ that lie over $p$. Also, let $\beta$ be a rational differential on $Y$.
(i) If $\beta$ is regular at the points $q_{1}, \ldots, q_{k}$, its trace $\tau(\beta)$ is regular at $p$.
(ii) If $\beta$ has a simple pole at $q_{i}$ and is regular at $q_{j}$ for all $j \neq i$, then $\tau(\beta)$ is not regular at $p$.
proof. (i) Corollary 8.6 .8 tells us that $\beta=g d x$, where $g$ has poles of orders at most $e_{i}-1$ at the points $q_{i}$. Since $x$ has a zero of order $e_{i}$ at $q_{i}$, the function $x g$ is regular at $q_{i}$, and its value there is zero. Then $\operatorname{tr}(x g)$ is regular at $p$, and its value at $p$ is zero 8.6.4. So $x^{-1} \operatorname{tr}(x g)$ is a regular function at $p$. Since $\operatorname{tr}$ is $F$-linear and $x$ is in $F, x^{-1} \operatorname{tr}(x g)=\operatorname{tr}(g)$. Therefore $\operatorname{tr}(g)$ and $\tau(\beta)=\operatorname{tr}(g) d x$ are regular at $p$.
(ii) With $\beta=g d x$, the function $x g$ will be regular at $p$. Its value at $q_{j}$ will be zero when $j \neq i$, and not zero when $j=i$. Then $\operatorname{tr}(x g)$ will be regular at $p$, but not zero there 8.6.4. Therefore $\tau(\beta)=x^{-1} \operatorname{tr}(x g) d x$ won't be regular at $p$.
8.6.10. Corollary. The trace map defines a homomorphism of $\mathcal{O}_{X}$-modules $\Omega_{Y} \xrightarrow{\tau} \Omega_{X}$.
8.6.11. Example. Let $Y$ be the locus $y^{e}=x$ in $\mathbb{A}_{x, y}^{2}$. Multiplication by $\zeta=e^{2 \pi i / e}$ permutes the sheets of $Y$ over $X$. The trace of a power $y^{k}$ is

$$
\begin{equation*}
\operatorname{tr}\left(y^{k}\right)=\sum_{j} \zeta^{k j} y^{k} \tag{8.6.12}
\end{equation*}
$$

The sum $\sum \zeta^{k j}$ is zero unless $k \equiv 0$ modulo $e$. So $d y=y^{1-e} d x / e$, and $\tau(d y)=\operatorname{tr}\left(y^{1-e}\right) d x / e=0$. But $y^{-1} d y=y^{-e} d x / e=x^{-1} d x / e$, and $\tau\left(y^{-1} d y\right)=\operatorname{tr}\left(x^{-1}\right) d x / e=x^{-1} \operatorname{tr}(1) d x / e=d x / x$. This isn't regular at $x=0$.

Let $Y$ be a smooth curve, and let $Y \xrightarrow{\pi} X$ be a branched covering. As is true for any $\mathcal{O}_{Y}$-module, the module of differentials $\Omega_{Y}$ is isomorphic to the module of homomorphisms $\mathcal{O}_{Y}\left(\mathcal{O}_{Y}, \Omega_{Y}\right)$. The homomorphism $\mathcal{O}_{Y} \rightarrow \Omega_{Y}$ that corresponds to a section $\beta$ of $\Omega_{Y}$ on an open set $U$ sends a regular function $f$ on $U$ to $f \beta$. We will denote that homomorphism by $\beta$ too: $\mathcal{O}_{Y} \xrightarrow{\beta} \Omega_{Y}$.
8.6.13. Lemma. Composition with the trace defines a homomorphism of $\mathcal{O}_{X}$-modules

$$
\Omega_{Y} \text { approx }_{c} o_{Y}\left(\mathcal{O}_{Y}, \Omega_{Y}\right) \xrightarrow{\tau} \mathcal{O}_{X}\left(\mathcal{O}_{Y}, \Omega_{X}\right)
$$

This is true because $\tau$ is $\mathcal{O}_{X}$-linear. An $\mathcal{O}_{Y}$-linear map becomes an $\mathcal{O}_{X}$-linear map by restriction of scalars. So when we compose an $\mathcal{O}_{Y}$-linear map $\beta$ with $\tau$, the result will be $\mathcal{O}_{X}$-linear. It is a homomorphism of $\mathcal{O}_{X}$-modules.
8.6.14. Theorem. (i) The map (??) is bijective.
(ii) More generally, if $\mathcal{M}$ is any locally free $\mathcal{O}_{Y}$-module, composition with the trace defines a bijection

$$
\begin{equation*}
\mathcal{O}_{Y}\left(\mathcal{M}, \Omega_{\mathcal{O}_{Y}}\right) \xrightarrow{\tau \circ} \mathcal{O}_{X}\left(\mathcal{M}, \Omega_{\mathcal{O}_{X}}\right) \tag{8.6.15}
\end{equation*}
$$

This theorem follows from the Main Lemma 8.6.9, when one looks closely.
Remark. The domain and range of the map 8.6.15) are to be interpreted as modules on $X$. For example, $\mathcal{O}_{Y}$ denotes the direct image on $X$ of the structure sheaf on $Y$. When we insert the symbols Hom and $\pi_{*}$ into the notation, 8.6.15 becomes an isomorphism

$$
\pi_{*}\left(\underline{\operatorname{Hom}}_{\mathcal{O}_{Y}}\left(\mathcal{M}, \Omega_{Y}\right)\right) \xrightarrow{\tau \circ} \underline{\operatorname{Hom}}_{\mathcal{O}_{X}}\left(\pi_{*} \mathcal{M}, \Omega_{X}\right)
$$

Because the theorem is about modules on $X$, we can verify it locally on $X$. In particular, we may suppose that $X$ and $Y$ are affine, $X=\operatorname{Spec} A$ and $Y=\operatorname{Spec} B$. When we state the theorem in terms of algebras and modules, the statement of the theorem for the affine varieties becomes this:
8.6.16. Theorem. Let $Y \rightarrow X$ be a branched covering, with $Y=\operatorname{Spec} B \rightarrow X=\operatorname{Spec} A$.
(i) The trace map $\Omega_{B}={ }_{B}\left(B, \Omega_{B}\right) \xrightarrow{\tau \circ} A\left(B, \Omega_{A}\right)$ is bijective.
(ii) For any locally free $B$-module $M$, composition with the trace defines a bijection $B\left(M, \Omega_{B}\right) \xrightarrow{\tau \circ} A\left(M, \Omega_{A}\right)$.
8.6.17. Lemma. Let $A \subset B$ be rings, let $M$ be a B-module, and let $N$ be an $A$-module. Then ${ }_{A}(M, N)$ has the structure of a $B$-module.

In the theorem and the lemma, when we write $A_{A}\left(M, \Omega_{B}\right)$ and ${ }_{A}(M, N)$, we are interpreting the $B$-modules $\Omega_{B}$ and $N$ as $A$-modules by restriction of scalars.
proof of the lemma We must define scalar multiplication of a homomorphism $M \xrightarrow{\varphi} N$ of $A$-modules by an element $b$ of $B$. The definition of $b \varphi$ is $[b \varphi](m)=\varphi(b m)$. Here one must show that this map $[b \varphi]$ is a homomorphism of $A$-modules $M \rightarrow N$, and that the axioms for a $B$-module are true for ${ }_{A}(M, N)$. You will be able to check these things.
proof of Theorem 8.6.14 (i). We use the algebra notation of Theorem 8.6.16 Since the theorem is local, we are still allowed to localize $A$. As $A$-modules, both $B$ and $\Omega_{B}$ are torsion-free, and therefore locally free. Localizing as needed, we may assume that they are free $A$-modules, and that $\Omega_{A}$ is a free module of rank one with basis $d x$. Then ${ }_{A}\left(B, \Omega_{A}\right)$ will be a free $A$-module too.

Let's denote ${ }_{A}\left(B, \Omega_{A}\right)$ by $\Theta$. Lemma 8.6 .17 tells us that $\Theta$ is a $B$-module. Because $B$ and $\Omega_{A}$ are locally free $A$-modules, $\Theta$ is a locally free $A$-module and a locally free $B$-module. Since $\Omega_{A}$ has $A$-rank 1 , the $A$ rank of $\Theta$ is the same as the $A$-rank of $B$. So the $B$-rank of $\Theta$ is 18.10 (ii). Therefore $\Theta$ is an invertible $B$-module.

The trace map $\Omega_{B} \xrightarrow{\tau} \Theta$ isn't the zero map because, if $x$ is a local coordinate on $X$, then $\tau d x \neq 0$. Since domain and range are invertible $B$-modules, $\tau$ is an injective homomorphism. Its image, which is isomorphic to $\Omega_{B}$, is an invertible submodule of the invertible $B$-module $\Theta$. Therefore $\Theta$ is isomorphic to an invertible module $\Omega_{B}(D)$ for some effective divisor $D$ 8.2.7. To complete the proof of the theorem, we show that the divisor $D$ is zero.

Suppose that $D>0$, let $q$ be a point in the support of $D$. Then $\Omega_{B}(q) \subset \Omega_{B}(D) \approx \Theta$. Let $p$ be the image of $q$ in $X$. We choose a rational differential $\beta$ in $\Omega_{K}$ that has a simple pole at $q$, and is regular at the other points of $Y$ in the fibre over $p$. The Chinese Remainder Theorem allows us to do this. According to Proposition 8.2.11, the trace $\tau(\beta)$ isn't regular at $p$. It isn't in $\Theta$.
proof of Theorem 8.6.14(ii). We go back to the statement in terms of $\mathcal{O}$-modules. We are to show that if $\mathcal{M}$ is a locally free $\mathcal{O}_{Y}$-module, composition with the trace defines a bijective map $\mathcal{O}_{Y}\left(\mathcal{M}, \Omega_{\mathcal{O}_{Y}}\right) \rightarrow \mathcal{O}_{X}\left(\mathcal{M}, \Omega_{\mathcal{O}_{X}}\right)$. Part (i) of the theorem tells us that this is true in when $\mathcal{M}=\mathcal{O}_{Y}$, and therefore it is also true when $\mathcal{M}$ is a free module $\mathcal{O}_{Y}^{k}$. And, since (ii) is a statement about $\mathcal{O}_{X}$-modules, it suffices to prove it locally on $X$. So it suffices to prove that a locally free $\mathcal{O}_{Y}$-module is free on the inverse image of an open set in $X$.
8.6.18. Lemma. Let $q_{1}, \ldots, q_{k}$ be points of a smooth curve $Y$, and let $\mathcal{M}$ be a locally free $\mathcal{O}_{Y}$-module. There is an open set $V$ that contains the points $q_{1}, \ldots, q_{k}$, such that $\mathcal{M}$ is free on $V$.

We assume the lemma and complete the proof of the theorem. Let $\left\{q_{1}, \ldots, q_{k}\right\}$ be the fibre over a point $p$ of $X$ and let be $V$ as in the lemma. The complement $D=Y-V$ is a finite set whose image $C$ in $X$ is a finite set that doesn't contain $p$. If $U$ is the complement of $C$ in $X$, its inverse image $W$ will be a subset of $V$ that contains the fibre and on which $\mathcal{M}$ is free.
proof of the lemma We may assume that $Y$ is affine, $Y=\operatorname{Spec} B$, and that the $\mathcal{O}$-module $\mathcal{M}$ corresponds to a locally free $B$-module $M$.

We go back to Lemma 8.1.3 Let $\mathfrak{m}_{i}$ be the maximal ideal of $B$ at $q_{i}$, and let $\bar{B}_{i}=B / \mathfrak{m}_{i}^{e_{i}}$. The quotient $\bar{B}=B / x B$ is isomorphic to the product $\bar{B}_{1} \times \cdots \times \bar{B}_{k}$. Since $\mathcal{M}$ is locally free, it is free in a neighborhood of each point $q_{i}$. Therefore $\bar{M}_{i}=M / \mathfrak{m}_{i} M$ is a free $\bar{B}_{i}$-module. whose dimension is the rank $r$ of $\mathcal{M}$.

Let $r$ be the rank of $M$. It follows from the Chinese Remainder Theorem that there are elements $m_{1}, \ldots, m_{r}$ in $M$ whose residues form a basis of $\bar{M}_{i}$ for every $i$. Let $V$ be the free $B$-module with basis $v_{1}, \ldots, v_{r}$, and let $V \xrightarrow{\varphi} M$ be the map defined by $\varphi\left(v_{i}\right)=m_{i}$, let $C$ be the cokernel of $\varphi$, and let $\bar{V}_{i}=V / \mathfrak{m}_{i} V$. The map $\varphi$ induces a bijection $\bar{V}_{i} \rightarrow \bar{M}_{i}$, and therefore $\bar{C}_{i}=C / \mathfrak{m}_{i} C$ is the zero module: $C=\mathfrak{m}_{i} C$. The Nakayama Lemma tells us that $C$ is zero at $q_{i}$. Since $C$ is a finite $B$-module, it is is supported on a finite set $F$. When
locallyfreeonX
we localize to delete this finite set from $X$, the elements $m_{1}, \ldots, m_{k}$ generate $M$, and since $M$ has rank $r$, they form a basis. Then $M$ is free.

Note. Theorem 8.6 .14 is subtle, and though its proof is understandable, it doesn't give much insight as to why the theorem is true. I don't like that. To get more insight, we would need a better understanding of differentials. As my father Emil Artin said,
"One doesn't really understand differentials, but one can learn to work with them."

### 8.7 The Riemann-Roch Theorem II

rroch serredual

## (8.7.1) the Serre dual

Let $Y$ be a smooth projective curve, and let $\mathcal{M}$ be a locally free $\mathcal{O}_{Y}$-module. The Serre dual of $\mathcal{M}$, which we will denote by $\mathcal{M}^{\#}$, is the module

$$
\begin{equation*}
\mathcal{M}^{\#}={ }_{Y}\left(\mathcal{M}, \Omega_{Y}\right) \quad\left(=\underline{\operatorname{Hom}}_{\mathcal{O}_{Y}}\left(\mathcal{M}, \Omega_{Y}\right)\right) \tag{8.7.2}
\end{equation*}
$$

Since the invertible module $\Omega_{Y}$ is locally isomorphic to $\mathcal{O}_{Y}$, the Serre dual $\mathcal{M}^{\#}$ will be locally isomorphic to the ordinary dual $\mathcal{M}^{*}$. It will be a locally free module with the same rank as $\mathcal{M}$, and the bidual $\left(\mathcal{M}^{\#}\right)^{\#}$ will be isomorphic to $\mathcal{M}$. This follows from Corollary 8.1.19, because $\mathcal{M}^{\#} \approx \mathcal{M}^{*} \otimes_{\mathcal{O}} \Omega_{Y}$. To spell this out,

$$
\left(\mathcal{M}^{\#}\right)^{\#} \approx\left(\mathcal{M}^{*} \otimes_{\mathcal{O}} \Omega_{Y}\right)^{*} \otimes_{\mathcal{O}} \Omega_{Y} \approx \mathcal{M}^{* *} \otimes_{\mathcal{O}} \Omega_{Y}^{*} \otimes_{\mathcal{O}} \Omega_{Y} \approx \mathcal{M}^{* *} \approx \mathcal{M}
$$

For example, $\mathcal{O}_{Y}^{\#}=\Omega_{Y}$ and $\Omega_{Y}^{\#}=\mathcal{O}_{Y}$.
8.7.3. Riemann-Roch Theorem, version 2. Let $\mathcal{M}$ be a locally free $\mathcal{O}_{Y}$-module on a smooth projective curve $Y$, and let $\mathcal{M}^{\#}$ be its Serre dual. Then $\mathbf{h}^{0} \mathcal{M}=\mathbf{h}^{1} \mathcal{M} \#$ and $\mathbf{h}^{1} \mathcal{M}=\mathbf{h}^{0} \mathcal{M}^{\#}$.

Because $\mathcal{M}$ and $\left(\mathcal{M}^{\#}\right)^{\#}$ are isomorphic, the two assertions of the theorem are equivalent.
For example, $\mathbf{h}^{1} \Omega_{Y}=\mathbf{h}^{0} \mathcal{O}_{Y}=1$ and $\mathbf{h}^{0} \Omega_{Y}=\mathbf{h}^{1} \mathcal{O}_{Y}=p_{a}$. If $\mathcal{M}$ is a locally free $\mathcal{O}_{Y}$-module, then

$$
\begin{equation*}
\chi(\mathcal{M})=\mathbf{h}^{0} \mathcal{M}-\mathbf{h}^{0} \mathcal{M}^{\#} \tag{8.7.4}
\end{equation*}
$$

A more precise statement of the Riemann-Roch Theorem is that $H^{1}(Y, \mathcal{M})$ and $H^{0}\left(Y, \mathcal{M}^{\#}\right)$ are dual vector spaces in a canonical way. We omit the proof of this. The fact that their dimensions are equal is enough for many applications. The canonical isomorphism becomes important when one wants to apply the theorem to a cohomology sequence.

Our plan is to prove Theorem 8.7.3 directly for the projective line. The structure of locally free modules on $\mathbb{P}^{1}$ is very simple, so this will be easy. Following Grothendieck, we derive it for an arbitrary smooth projective curve $Y$ by projection to $\mathbb{P}^{1}$.

Let $Y$ be a smooth projective curve, let $X=\mathbb{P}^{1}$, and let $Y \xrightarrow{\pi} X$ be a branched covering. Let $\mathcal{M}$ be a locally free $\mathcal{O}_{Y}$-module, and let the Serre dual of $\mathcal{M}$, as defined in 8.7.2, be

$$
\mathcal{M}_{1}^{\#}={ }_{Y}\left(\mathcal{M}, \Omega_{Y}\right)
$$

The direct image of $\mathcal{M}$ is a locally free $\mathcal{O}_{X}$-module that we are denoting by $\mathcal{M}$ too, and we can form the Serre dual on $X$. Let

$$
\mathcal{M}_{2}^{\#}={ }_{x}\left(\mathcal{M}, \Omega_{X}\right)
$$

8.7.5. Corollary. The direct image $\pi_{*} \mathcal{M}_{1}^{\#}$, which we denote by $\mathcal{M}_{1}^{\#}$, is isomorphic to $\mathcal{M}_{2}^{\#}$.
proof. This is Theorem ??.
The corollary allows us to drop the subscripts from $\mathcal{M} \#$. Because a branched covering $Y \xrightarrow{\pi} X$ is an affine morphism, the cohomology of $\mathcal{M}$ and of its Serre dual $\mathcal{M}{ }^{\#}$ can be computed, either on $Y$ or on $X$. If $\mathcal{M}$ is a locally free $\mathcal{O}_{Y}$-module, then $H^{q}(Y, \mathcal{M}) \approx H^{q}(X, \mathcal{M})$ and $H^{q}\left(Y, \mathcal{M}^{\#}\right) \approx H^{q}\left(X, \mathcal{M}^{\#}\right)$ (see (7.4.25)).

Thus it is enough to prove Riemann-Roch for the projective line.

## (8.7.6) Riemann-Roch for the projective line

The Riemann-Roch Theorem for the projective line $X=\mathbb{P}^{1}$ is a consequence of the Birkhoff-Grothendieck Theorem tells us that a locally free $\mathcal{O}_{X}$-module $\mathcal{M}$ on $X$ is a direct sum of twisting modules. To prove Riemann-Roch for the projective line, it suffices to prove it for the twisting modules $\mathcal{O}_{X}(k)$
8.7.7. Lemma. The module of differentials $\Omega_{X}$ on $X$ is isomorphic to the twisting module $\mathcal{O}_{X}(-2)$.
proof. Since $\Omega_{X}$ is invertible, the Birkhoff-Grothendieck Theorem tells us that it is isomorphic to the twisting module $\mathcal{O}_{X}(k)$ for some $k$. We need only identify the integer $k$.

Let $\mathbb{U}^{0}=\operatorname{Spec} \mathbb{C}[x]$, and $\mathbb{U}^{1}=\operatorname{Spec} \mathbb{C}[z]$ be the standard open subsets of $\mathbb{P}^{1}$, with $z=x^{-1}$. On $\mathbb{U}^{0}$, the module of differentials is free, with basis $d x$, and $d x=d\left(z^{-1}\right)=-z^{-2} d z$ describes the differential $d x$ on $\mathbb{U}^{1}$. Since the point $p$ at infinity is $\{z=0\}, d x$ has a pole of order 2 there. It is a global section of $\Omega_{X}(2 p)$, and as a section of that module, it isn't zero anywhere. So multiplication by $d x$ defines an isomorphism $\mathcal{O} \rightarrow \Omega_{X}(2 p)$ that sends 1 to $d x$. Tensoring with $\mathcal{O}(-2 p)$, we find that $\Omega_{X}$ is isomorphic to $\mathcal{O}(-2 p)$.
8.7.8. Lemma. Let let $\mathcal{M}$ and $\mathcal{N}$ be locally free $\mathcal{O}$-modules on the projective line $X$. Then ${ }_{X}(\mathcal{M}(r), \mathcal{N})$ is canonically isomorphic to $x_{x}(\mathcal{M}, \mathcal{N}(-r))$.
proof. When we tensor a homomorphism $\mathcal{M}(r) \xrightarrow{\varphi} \mathcal{N}$ with $\mathcal{O}(-r)$, we obtain a homomorphism $\mathcal{M} \rightarrow$ $\mathcal{N}(-r)$. Tensoring with $\mathcal{O}(r)$ is the inverse operation.

The Serre dual $\mathcal{O}(n)^{\#}$ of $\mathcal{O}(n)$ is therefore

$$
\mathcal{O}(n)^{\#}={ }_{x}(\mathcal{O}(n), \mathcal{O}(-2)) \approx \mathcal{O}(-2-n)
$$

To prove Riemann-Roch for $X=\mathbb{P}^{1}$, we must show that

$$
\mathbf{h}^{0} \mathcal{O}(n)=\mathbf{h}^{1} X, \mathcal{O}(-2-n) \quad \text { and } \quad \mathbf{h}^{1} \mathcal{O}(n)=\mathbf{h}^{0} \mathcal{O}(-2-n)
$$

This follows from (Theorem7.5.4), which computes the cohomology of the twisting modules. As we've noted before, these two assertions are equivalent, so it suffices to verify the first one. If $n<0$, then $-2-n>0$. In this case $\mathbf{h}^{0} \mathcal{O}(n)=\mathbf{h}^{1} \mathcal{O}(-2-n)=0$. If $n \geq 0$, Theorem 7.5.4 asserts that $\mathbf{h}^{0} \mathcal{O}(n)=n+1$ and that $\mathbf{h}^{1} \mathcal{O}(-2-n)=(2+n)-1=n+1$.

### 8.8 Using Riemann-Roch

## (8.8.1) genus

There are three closely related numbers associated to a smooth projective curve $Y$ : its topological genus $g$, its arithmetic genus $p_{a}=\mathbf{h}^{1} \mathcal{O}_{Y}$, and the degree $\delta$ of the module of differentials $\Omega_{Y}$ (see 8.2.
8.8.2. Theorem. Let $Y$ be a smooth projective curve. The topological genus $g$ and the arithmetic genus $p_{a}$ of $Y$ are equal, and the degree $\delta$ of the module $\Omega_{Y}$ is $2 p_{a}-2$, which is equal to $2 g-2$.
dualityfor pone

Thus the Riemann-Roch Theorem8.3.3 can we written as

$$
\chi(\mathcal{O}(D))=\operatorname{deg} D+1-g
$$

We'll write it this way in what follows.
proof. Let $Y \xrightarrow{\pi} X$ be a branched covering with $X=\mathbb{P}^{1}$. The topological Euler characteristic $e(Y)$, which is $2-2 g$, can be computed in terms of the branching data for the covering, as in 1.7.27). Let $q_{i}$ be the ramification points in $Y$, and let $e_{i}$ be the ramification index at $q_{i}$. Then $e_{i}$ sheets of the covering come together at $q_{i}$. One might say that $e_{i}-1$ points are missing. If the degree of $Y$ over $X$ is $n$, then since $e(X)=2$,

$$
\begin{equation*}
2-2 g=e(Y)=n e(X)-\sum\left(e_{i}-1\right)=2 n-\sum\left(e_{i}-1\right) \tag{8.8.3}
\end{equation*}
$$

We compute the degree $\delta$ of $\Omega_{Y}$ in two ways. First, the Riemann-Roch Theorem tells us that $\mathbf{h}^{0} \Omega_{Y}=$ $\mathbf{h}^{1} \mathcal{O}_{Y}=p_{a}$ and $\mathbf{h}^{1} \Omega_{Y}=\mathbf{h}^{0} \mathcal{O}_{Y}=1$. So $\chi\left(\Omega_{Y}\right)=-\chi\left(\mathcal{O}_{Y}\right)=p_{a}-1$. The Riemann-Roch Theorem also tells us that $\chi\left(\Omega_{Y}\right)=\delta+1-p_{a}$ (??). Therefore

$$
\begin{equation*}
\delta=2 p_{a}-2 \tag{8.8.4}
\end{equation*}
$$

Next, we compute $\delta$ by computing the divisor of the differential $d x$ on $Y, x$ being a coordinate in $X$. Let $q_{i}$ be one of the ramification points in $Y$, and let $e_{i}$ be the ramification index at $q_{i}$. Then $d x$ has a zero of order $e_{i}-1$ at $q_{i}$. On $X, d x$ has a pole of order 2 at $\infty$. Let's suppose that the point at infinity isn't a branch point. Then there will be $n$ points of $Y$ at which $d x$ has a pole of order $2, n$ being the degree of $Y$ over $X$, as above. The degree of $\Omega_{Y}$ is therefore

$$
\begin{equation*}
\delta=\text { zeros }- \text { poles }=\sum\left(e_{i}-1\right)-2 n \tag{8.8.5}
\end{equation*}
$$

Combining 8.8.5 with 8.8.3, one sees that $\delta=2 g-2$. Since we also have $\delta=2 p_{a}-2$, we conclude that $g=p_{a}$.
Honezero 8.8.6. Corollary. Let $D$ be a divisor on a smooth projective curve $Y$ of genus $g$. If $\operatorname{deg} D>2 g-2$, then $\mathbf{h}^{1} \mathcal{O}(D)=0$. If $\operatorname{deg} D \leq g-2$, then $\mathbf{h}^{1} \mathcal{O}(D)>0$.
proof. This follows from Corollary 8.3.10 (v).

## (8.8.7) curves of genus zero

Let $Y$ be a smooth projective curve $Y$ whose genus $g$ is zero, and let $p$ be a point of $Y$. The exact sequence

$$
0 \rightarrow \mathcal{O}_{Y} \rightarrow \mathcal{O}_{Y}(p) \rightarrow \epsilon \rightarrow 0
$$

where $\epsilon$ is a one-dimensional module supported at $p$, gives us an exact cohomology sequence

$$
0 \rightarrow H^{0}\left(Y, \mathcal{O}_{Y}\right) \rightarrow H^{0}\left(Y, \mathcal{O}_{Y}(p)\right) \rightarrow H^{0}(Y, \epsilon) \rightarrow 0
$$

The zero on the right is due to the fact that $\mathbf{h}^{1} \mathcal{O}_{Y}=g=0$. We also have $\mathbf{h}^{0} \mathcal{O}_{Y}=1$ and $\mathbf{h}^{0} \epsilon=1$, so $\mathbf{h}^{0} \mathcal{O}_{Y}(p)=2$. We choose a basis $(1, x)$ for $H^{0}\left(Y, \mathcal{O}_{Y}(p)\right), 1$ being the constant function and $x$ being a nonconstant function with a single pole of order 1 at $p$. This basis defines a point of $\mathbb{P}^{1}$ with values in the function field $K$ of $Y$, and therefore a morphism $Y \xrightarrow{\varphi} \mathbb{P}^{1}$. Because $x$ has just one pole of order 1, it takes every value exactly once. Therefore $\varphi$ is bijective. It is a map of degree 1 , and therefore an isomorphism 8.1.5.
8.8.8. Corollary. Every smooth projective curve of genus zero is isomorphic to the projective line $\mathbb{P}^{1}$.

A curve, smooth or not, whose function field is isomorphic to the field $\mathbb{C}(t)$ of rational functions in one variable is called a rational curve. A smooth projective curve of genus zero is a rational curve.

## (8.8.9) curves of genus one

A smooth projective curve of genus $g=1$ is called an elliptic curve. The Riemann-Roch Theorem tells us that on an elliptic curve $Y$,

$$
\chi(\mathcal{O}(D))=\operatorname{deg} D
$$

Since $\mathbf{h}^{0} \Omega_{Y}=\mathbf{h}^{1} \mathcal{O}_{Y}=1, \Omega_{Y}$ has a nonzero global section $\omega$. Since $\Omega_{Y}$ has degree zero 8.8.2, $\omega$ doesn't vanish anywhere. Multiplication by $\omega$ defines an isomorphism $\mathcal{O} \rightarrow \Omega_{Y}$. So $\Omega_{Y}$ is a free module of rank one.

The next lemma follows from Riemann-Roch.
8.8.10. Lemma. Let p be a point of an elliptic curve $Y$. For any $r>0, \mathbf{h}^{0} \mathcal{O}(r p)=r$, and $\mathbf{h}^{1} \mathcal{O}(r p)=0$. $\qquad$
hrpforel-
liptic

Since $H^{0}\left(Y, \mathcal{O}_{Y}\right) \subset H^{0}\left(Y, \mathcal{O}_{Y}(p)\right)$, and since both spaces have dimension one, they are equal. So (1) is a basis for $H^{0}\left(Y, \mathcal{O}_{Y}(p)\right)$. We choose a basis $(1, x)$ for the two-dimensional space $H^{1}\left(Y, \mathcal{O}_{Y}(2 p)\right)$. Then $x$ isn't a section of $\mathcal{O}(p)$. It has a pole of order precisely 2 at $p$ and no other pole. Next, we choose a basis $(1, x, y)$ for $H^{1}\left(Y, \mathcal{O}_{Y}(3 p)\right)$. So $y$ has a pole of order 3 at $p$, and no other pole. The point $(1, x, y)$ of $\mathbb{P}^{2}$ with values in $K$ determines a morphism $Y \xrightarrow{\varphi} \mathbb{P}^{2}$.

Let $u, v, w$ be coordinates in $\mathbb{P}^{2}$. The map $\varphi$ sends a point $q$ distinct from $p$ to $(u, v, w)=(1, x(q), y(q))$. Since $Y$ has dimension one, $\varphi$ is a finite morphism. Its image will be a closed subvariety of $\mathbb{P}^{2}$ of dimension one.

To determine the image of the point $p$, we multiply $(1, x, y)$ by $\lambda=y^{-1}$ to normalize the second coordinate to 1 , obtaining the equivalent vector $\left(y^{-1}, x y^{-1}, 1\right)$. The rational function $y^{-1}$ has a zero of order 3 at $p$, and $x y^{-1}$ has a simple zero there. Evaluating at $p$, we see that the image of $p$ is the point $(0,0,1)$.

Let $Y^{\prime}$ be the image of $Y$, which is a curve in $\mathbb{P}^{2}$. The map $Y \rightarrow \mathbb{P}^{2}$ restricts to a finite morphism $Y \rightarrow Y^{\prime}$. Let $\ell$ be a generic line $\{a u+b v+c w=0\}$ in $\mathbb{P}^{2}$. The rational function $a+b x+c y$ on $Y$ has a pole of order 3 at $p$ and no other pole. It takes every value, including zero, three times, and the set of points $q$ of $Y$ at which $a+b x+c y$ is zero is the inverse image of the intersection $Y^{\prime} \cap \ell$. The only possibilities for the degree of $Y^{\prime}$ are 1 and 3 . Since $1, x, y$ are independent, they don't satisfy any homogeneous linear equation. So $Y^{\prime}$ isn't a line. The image $Y^{\prime}$ is a cubic curve (see Corollary 1.3.9).

To determine the image, we look for a cubic relation among the functions $1, x, y$ on $Y$. The seven monomials $1, x, y, x^{2}, x y, x^{3}, y^{2}$ have poles at $p$ of orders $0,2,3,4,5,6,6$, respectively, and no other poles. They are sections of $\mathcal{O}_{Y}(6 p)$. Riemann-Roch tells us that $\mathbf{h}^{0} \mathcal{O}_{Y}(6 p)=6$. So those seven functions are linearly dependent. The dependency relation gives us a cubic equation among $x$ and $y$, which we may write in the form

$$
c y^{2}+\left(a_{1} x+a_{3}\right) y+\left(a_{0} x^{3}+a_{2} x^{2}+a_{4} x+a_{6}\right)=0
$$

There can be no linear relation among functions whose orders of pole at $p$ are distinct. So, when we delete either $x^{3}$ or $y^{2}$ from the list of monomials, we obtain an independent set of six functions. They form a basis for the six-dimensional space $H^{0}(Y, \mathcal{O}(6 p))$. So, in the cubic relation, the coefficients $c$ and $a_{0}$ aren't zero. We can scale $y$ and $x$ to normalize $c$ and $a_{0}$ to 1 . We eliminate the linear term in $y$ from this relation by substituting $y-\frac{1}{2}\left(a_{1} x+a_{3}\right)$ for $y$. Next, we eliminate the quadratic term in $x$. by substituting $x-\frac{1}{3} a_{2}$ for $x$. Bringing the terms in $x$ to the other side of the equation, we are left with a cubic relation of the form

$$
y^{2}=x^{3}+a_{4} x+a_{6}
$$

The coefficients $a_{4}$ and $a_{6}$ have changed, of course.
The cubic curve $Y^{\prime}$ defined by the homogenized equation $y^{2} z=x^{3}+a_{4} x z^{2}+a_{6} z^{3}$ is the image of $Y$. This curve meets a generic line $a x+b y+c z=0$ in three points and, as we saw above, its inverse image in $Y$ consists of three points too. Therefore the morphism $Y \xrightarrow{\varphi} Y^{\prime}$ is generically injective, and $Y$ is the normalization of $Y^{\prime}$. Corollary 7.6.3 computes the cohomology of $Y^{\prime}: \mathbf{h}^{0} \mathcal{O}_{Y^{\prime}}=\mathbf{h}^{1} \mathcal{O}_{Y^{\prime}}=1$. This tells us that $\mathbf{h}^{q} \mathcal{O}_{Y^{\prime}}=\mathbf{h}^{q} \mathcal{O}_{Y}$ for all $q$. Let's denote the direct image of $\mathcal{O}_{Y}$ by the same symbol $\mathcal{O}_{Y}$, and let $\mathcal{F}$ be the $\mathcal{O}_{Y^{\prime}}$-module $\mathcal{O}_{Y} / \mathcal{O}_{Y^{\prime}}$. The exact sequence $0 \rightarrow \mathcal{O}_{Y^{\prime}} \rightarrow \mathcal{O}_{Y} \rightarrow \mathcal{F} \rightarrow 0$ shows that $\mathbf{h}^{0} \mathcal{F}=0$. Since $Y$ is the normalization of $Y^{\prime}, \mathcal{F}$ is a torsion module with no global sections. So $\mathcal{F}=0$, and $Y \approx Y^{\prime}$.
8.8.11. Corollary. Every elliptic curve is isomorphic to a cubic curve in $\mathbb{P}^{2}$.

## (8.8.12) the group law on an elliptic curve

The points of an elliptic curve form an abelian group, once one chooses a point to be the identity element.
We choose a point of an elliptic curve $Y$ and label it $o$. We'll write the law of composition in the group as $p \oplus q$, using the symbol $\oplus$ to distinguish the sum in the group, which is a point of $Y$, from the divisor $p+q$.

Let $p$ and $q$ be points of $Y$. To define $p \oplus q$, we compute the cohomology of $\mathcal{O}_{Y}(p+q-o)$. It follows from Riemann-Roch that $\mathbf{h}^{0} \mathcal{O}_{Y}(p+q-o)=1$ and that $\mathbf{h}^{1} \mathcal{O}_{Y}(p+q-o)=0$. There is a nonzero function $f$, unique up to scalar factor, with simple poles at $p$ and $q$ and a zero at $o$. This function has exactly one other zero. That zero is defined to be the sum $p \oplus q$ in the group. In terms of linearly equivalent divisors, $s=p \oplus q$ is the unique point such that $p+q$ is linearly equivalent to $o+s$.
grplaw
maptoP
mapcurve
basept
degnobasept
8.8.13. Proposition. The law of composition $\oplus$ defined above makes an ellipic curve into an abelian group.

The proof is an exercise.

## (8.8.14) interlude: maps to projective space

Let $Y$ be a smooth projective curve. We have seen that any set $\left(f_{0}, \ldots, f_{n}\right)$ of rational functions on $Y$, not all zero, defines a morphism $Y \xrightarrow{\varphi} \mathbb{P}^{n}$ (5.3.3). As a reminder, let $q$ be a point of $Y$ and let $g_{j}=f_{j} / f_{i}$, where $i$ is an index such that $f_{i}$ has the minimum value $\mathrm{v}_{q}\left(f_{i}\right)$. Then $g_{j}$ are regular at $q$ for all $j$, and the morphism $\varphi$ sends the point $q$ to is $\left(g_{0}(q), \ldots, g_{n}(q)\right)$. For example, the inverse image $\varphi^{-1}\left(\mathbb{U}^{0}\right)$ of the standard open set $\mathbb{U}^{0}$ is the set of points of $Y$ at which the functions $g_{j}=f_{j} / f_{0}$ are regular. If $q$ is such a point, then $\varphi(q)=\left(1, g_{1}(q), \ldots, g_{n}(q)\right)$.
8.8.15. Lemma. Let $Y$ be a smooth projective curve, and let $Y \xrightarrow{\varphi} \mathbb{P}^{n}$ be the morphism to projective space that is defined by some set $\left(f_{0}, \ldots, f_{n}\right)$ of rational functions on $Y$ that aren't all zero.
(i) If the subspace of the functionfield of $Y$ that is spanned by $\left\{f_{0}, \ldots, f_{n}\right\}$ has dimension at least two, then $\varphi$ is not a constant function.
(ii) If $\left\{f_{0}, \ldots, f_{n}\right\}$ are linearly independent, the image isn't contained in any hyperplane.

The degree $d$ of a nonconstant morphism $Y \xrightarrow{\varphi} \mathbb{P}^{n}$ from a projective curve $Y$, smooth or not, to projective space is the number of points of the inverse image $\varphi^{-1} H$ of a generic hyperplane $H$ in $\mathbb{P}^{n}$. We check that this number is well-defined. Say that $H$ is the locus $h(x)=0$, where $h=\sum a_{i} x_{i}$, and that a second generic hyperplane $G$ is the locus $g(x)=0$, where $g=\sum b_{i} x_{i}$. Let $f(x)=h / g$, and let $\widetilde{f}=f \circ \varphi$. The divisor of $\tilde{f}$ on $Y$ is $\varphi^{-1} H-\varphi^{-1} G$. It has degree zero.

## (8.8.16) base points

Let $D$ be a divisor on the smooth projective curve $Y$ such that $\mathbf{h}^{0} \mathcal{O}(D)=k>1$. A basis $\left(f_{0}, \ldots, f_{k}\right)$ of global sections of $\mathcal{O}(D)$ defines a morphism $Y \rightarrow \mathbb{P}^{k-1}$. This is the most common way to construct a morphism to projective space, though one could use any set of rational functions.

If a global section of $\mathcal{O}(D)$ vanishes at a point $p$ of $Y$, it is a global section of $\mathcal{O}(D-p)$. A point $p$ is a base point of $\mathcal{O}(D)$ if every global section of $\mathcal{O}(D)$ vanishes at $p$. A base point can be described in terms of the usual exact sequence

$$
0 \rightarrow \mathcal{O}(D-p) \rightarrow \mathcal{O}(D) \rightarrow \epsilon \rightarrow 0
$$

The point $p$ is a base point if $\mathbf{h}^{0} \mathcal{O}(D-p)=\mathbf{h}^{0} \mathcal{O}(D)$, or if $\mathbf{h}^{1} \mathcal{O}(D-p)=\mathbf{h}^{1} \mathcal{O}(D)-1$.
8.8.17. Lemma. Let $D$ be a divisor on a smooth projective curve $Y$, and suppose that $\mathbf{h}^{0} \mathcal{O}(D)>1$. Let $Y \xrightarrow{\varphi} \mathbb{P}^{n}$ be the morphism defined by a basis of global sections.
(i) The image of $\varphi$ isn't contained in any hyperplane.
(ii) If $\mathcal{O}(D)$ has no base points, the degree $r$ of the morphism $\varphi$ is equal to degree of $D$. If there are base points, the degree is lower.

## (8.8.18) canonical divisors

Because the module $\Omega_{Y}$ of differentials on a smooth curve $Y$ is invertible, it is isomorphic to $\mathcal{O}(K)$ for some divisor $K$. Such a divisor $K$ is called a canonical divisor. The degree of $K$ is $2 g-2$. It is often convenient to represent $\Omega_{Y}$ as a module $\mathcal{O}(K)$, though the canonical divisor $K$ isn't unique. It is determined only up to linear equivalence (see (8.2.13)).

When written in terms of a canonical divisor $K$, the Serre dual of an invertible module $\mathcal{O}(D)$ will be $\mathcal{O}(D)^{\#}=\mathcal{O}(\mathcal{O}(D), \mathcal{O}(K)) \approx \mathcal{O}(K-D)$. With this notation, the Riemann-Roch Theorem for $\mathcal{O}(D)$ becomes

$$
\begin{equation*}
\mathbf{h}^{0} \mathcal{O}(D)=\mathbf{h}^{1} \mathcal{O}(K-D) \quad \text { and } \quad \mathbf{h}^{1} \mathcal{O}(D)=\mathbf{h}^{0} \mathcal{O}(K-D) \tag{8.8.19}
\end{equation*}
$$

8.8.20. Proposition. Let $K$ be a canonical divisor on a smooth projective curve $Y$ of genus $g>0$.
(i) $\mathcal{O}(K)$ has no base point.
(ii) Every point p of $Y$ is a base point of $\mathcal{O}(K+p)$.
proof. (i) Let $p$ be a point of $Y$. We apply Riemann-Roch to the exact sequence

$$
0 \rightarrow \mathcal{O}(K-p) \rightarrow \mathcal{O}(K) \rightarrow \epsilon \rightarrow 0
$$

where $\epsilon$ denotes a one-dimensional module supported on a point $p$. The Serre duals of $\mathcal{O}(K)$ and $\mathcal{O}(K-p)$ are $\mathcal{O}$ and $\mathcal{O}(p)$, respectively. They form an exact sequence

$$
0 \rightarrow \mathcal{O} \rightarrow \mathcal{O}(p) \rightarrow \epsilon^{\prime} \rightarrow 0
$$

When $Y$ has positive genus, there is no rational function on $Y$ with just one simple pole. So $\mathbf{h}^{0} \mathcal{O}(p)=1$. Riemann-Roch tells us that $\mathbf{h}^{1} \mathcal{O}(K-p)=1$. Also, $\mathbf{h}^{1} \mathcal{O}(K)=1$. The cohomology sequence

$$
0 \rightarrow H^{0}(\mathcal{O}(K-p)) \rightarrow H^{0}(\mathcal{O}(K)) \rightarrow[1] \rightarrow H^{1}(\mathcal{O}(K-p)) \rightarrow H^{1}(\mathcal{O}(K)) \rightarrow 0
$$

shows that $\mathbf{h}^{0} \mathcal{O}(K-p)=\mathbf{h}^{0} \mathcal{O}(K)-1$. So $p$ is not a base point.
(ii) Here, the relevant sequence is

$$
0 \rightarrow \mathcal{O}(K) \rightarrow \mathcal{O}(K+p) \rightarrow \epsilon_{3} \rightarrow 0
$$

The Serre dual of $\mathcal{O}(K+p)$ is $\mathcal{O}(-p)$, which has no global section. Therefore $\mathbf{h}^{1} \mathcal{O}(K+p)=0$, while $\left.\mathbf{h}^{1} \mathcal{O}(K)=\mathbf{h}^{0} \mathcal{O}\right)=1$. The cohomology sequence

$$
0 \rightarrow \mathbf{h}^{0} \mathcal{O}(K) \rightarrow \mathbf{h}^{0} \mathcal{O}(K+p) \rightarrow[1] \rightarrow \mathbf{h}^{1} \mathcal{O}(K) \rightarrow \mathbf{h}^{1} \mathcal{O}(K+p) \rightarrow 0
$$

shows that $H^{0}(\mathcal{O}(K+p))=H^{0}(\mathcal{O}(K))$. So $p$ is a base point of $\mathcal{O}(K+p)$.

## (8.8.21) hyperelliptic curves

A hyperelliptic curve $Y$ is a smooth projective curve of genus $g>1$ that can be represented as a branched double covering of the projective line. So $Y$ is hyperelliptic if there is a morphism $Y \xrightarrow{\pi} X$ of degree two, with $X=\mathbb{P}^{1}$. Justification for the strange term 'hyperelliptic' is that every elliptic curve can be represented as a double cover of $\mathbb{P}^{1}$, by the map to $\mathbb{P}^{1}$ defined by the global sections of $\mathcal{O}(2 p)$.

The topological Euler characteristic of a hyperelliptic curve $Y$ can be computed in terms of the covering $Y \rightarrow X$, which will be branched at a finite set, say of $n$ points. Since $\pi$ has degree two, the multiplicity of a branch point will be 2 . The Euler characteristic is therefore $e(Y)=2 e(X)-n=4-n$. Since we know that $e(Y)=2-2 g$, the number $n$ of branch points is $2 g+2$. When $g=3, n=8$.

It would take some experimentation to guess that the next remarkable theorem might be true, and to find a proof.
8.8.22. Theorem. Let $K$ be a canonical divisor on a hyperelliptic curve $Y$, and let $Y \xrightarrow{\pi} X=\mathbb{P}^{1}$ be the associated branched covering of degree 2. Let $Y \xrightarrow{\kappa} \mathbb{P}^{g-1}$ be the morphism defined by the global sections of $\Omega_{Y}=\mathcal{O}(K)$. This morphism $\kappa$ factors through $X$ : There is a morphism $X \xrightarrow{u} \mathbb{P}^{g-1}$ such that $\kappa$ is the composition $u \circ \pi$ :

onedblcover
8.8.23. Corollary. A curve of genus $g \geq 2$ can be presented as a branched covering of $\mathbb{P}^{1}$ of degree 2 in at most one way.

## proof of Theorem 8.8.22

Let $x$ be an affine coordinate in $X$, so that the standard affine open subset $\mathbb{U}^{0}$ of $X$ is $\operatorname{Spec} \mathbb{C}[x]$. We suppose that the point $p_{\infty}$ at infinity of $X$ isn't a branch point of the covering. Let $Y^{0}=\pi^{-1} \mathbb{U}^{0}$. Then $Y^{0}$ will have an equation of the form

$$
y^{2}=f(x)
$$

where $f$ is a polynomial with $n=2 g+2$ simple roots. There will be two points of $Y$ above the point $p_{\infty}$, which are interchanged by the automorphism $y \rightarrow-y$. Let's call those points $q_{1}$ and $q_{2}$.

We start with the differential $d x$, which we view as a differential on $Y$. Then $2 y d y=f^{\prime}(x) d x$. Since $f$ has simple roots, $f^{\prime}$ doesn't vanish at any of them. Therefore $d x$ has simple zeros on $Y$ above the roots of $f$, the points at which $y=0$. We also have a regular function on $Y^{0}$ with simple roots at those points, namely the function $y$. Therefore the differential $\omega=\frac{d x}{y}$ is regular and nowhere zero on $Y^{0}$. Because the degree of a differential on $Y$ is $2 g-2$, $\omega$ has a total of $2 g-2$ zeros at infinity. By symmetry, $\omega$ has zeros of order $g-1$ at the each of two points $q_{1}$ and $q_{2}$. So $K=(g-1) q_{1}+(g-1) q_{2}$ is a canonical divisor on $Y$, i.e., $\Omega_{Y} \approx \mathcal{O}_{Y}(K)$.

Since $K$ has zeros of order $g-1$ at infinity, the rational functions $1, x, x^{2}, \ldots, x^{g-1}$, when viewed as functions on $Y$, are among the global sections of $\mathcal{O}_{Y}(K)$. They are independent, and there are $g$ of them. Since $\mathbf{h}^{0} \mathcal{O}_{Y}(K)=g$, they form a basis of $H^{0}\left(\mathcal{O}_{Y}(K)\right)$. The map $Y \rightarrow \mathbb{P}^{g-1}$ defined by the global sections of $\mathcal{O}_{Y}(K)$ evaluates these powers of $x$, so it factors through $X$.

## (8.8.24) canonical embedding

Let $Y$ be a smooth projective curve of genus $g \geq 2$, and let $K$ be a canonical divisor on $Y$. Since $\mathcal{O}(K)$ has no base point (??), its global sections define a morphism $Y \rightarrow \mathbb{P}^{g-1}$ that is called the canonical map. Let's denote the canonical map by $\kappa$. The degree of $\kappa$ is the degree $2 g-2$ of the canonical divisor.
8.8.25. Theorem. Let $Y$ be a smooth projective curve of genus $g$ at least two. If $Y$ isn't hyperelliptic, the canonical map embeds $Y$ as a closed subvariety of projective space $\mathbb{P}^{g-1}$.
proof. We show first that, if the canonical map $Y \xrightarrow{\kappa} \mathbb{P}^{g-1}$ isn't injective, then $Y$ is hyperelliptic. Let $p$ and $q$ be distinct points of $Y$ such that $\kappa(p)=\kappa(q)$. We choose an effective canonical divisor such that $p$ and $q$ aren't in its support. We inspect the global sections of $\mathcal{O}(K-p-q)$. Since $\kappa(p)=\kappa(q)$, any global section of $\mathcal{O}(K)$ that vanishes at $p$ vanishes at $q$ too. Therefore $\mathcal{O}(K-p)$ and $\mathcal{O}(K-p-q)$ have the same global sections, and $q$ is a base point of $\mathcal{O}(K-p)$. We've computed the cohomology of $\mathcal{O}(K-p)$ before: $\mathbf{h}^{0} \mathcal{O}(K-p)=g-1$ and $\mathbf{h}^{1} \mathcal{O}(K-p)=1$. Then $\mathbf{h}^{0} \mathcal{O}(K-p-q)=g-1$ and $\mathbf{h}^{1} \mathcal{O}(K-p-q)=2$. The Serre dual of $\mathcal{O}(K-p-q)$ is $\mathcal{O}(p+q)$, so by Riemann-Roch, $\mathbf{h}^{0} \mathcal{O}(p+q)=2$. If $D$ is a divisor of degree one on a curve of positive genus, then $\mathbf{h}^{0} \mathcal{O}(D) \leq 1$ (Proposition ??). Therefore $\mathcal{O}(p+q)$ has no base point. Its global sections define a morphism $Y \rightarrow \mathbb{P}^{1}$ of degree 2 . So $Y$ is hyperelliptic.

If $Y$ isn't hyperelliptic, the canonical map is injective, so $Y$ is mapped bijectively to its image $Y^{\prime}$ in $\mathbb{P}^{g-1}$. This almost proves the theorem. But: Can $Y^{\prime}$ have a cusp? We must show that the bijective map $Y \xrightarrow{\kappa} Y^{\prime}$ is an isomorphism.

We go over the computation made above for a pair of points $p, q$, this time taking $q=p$. The computation is the same. Since $Y$ isn't hyperelliptic, $p$ isn't a base point of $\mathcal{O}_{Y}(K-p)$. Therefore $\mathbf{h}^{0} \mathcal{O}_{Y}(K-2 p)=$ $\mathbf{h}^{0} \mathcal{O}_{Y}(K-p)-1$. This tells us that there is a global section $f$ of $\mathcal{O}_{Y}(K)$ that has a zero of order exactly 1 at $p$. When properly interpreted, this fact shows that $\kappa$ doesn't collapse any tangent vectors to $Y$, and that $\kappa$ is an isomorphism. Since we haven't discussed tangent vectors, we prove this directly.

## \#\#\# reread this\#\#\#

Since $\kappa$ is bijective, the function fields of $Y$ and its image $Y^{\prime}$ are equal, and $Y$ is the normalization of $Y^{\prime}$. Moreover, $\kappa$ is an isomorphism except on a finite set. We work locally at a point $p$ of $Y^{\prime}$, and we denote the unique point of $Y$ that maps to $Y^{\prime}$ by $p$ too. When we restrict the global section $f$ of $\mathcal{O}_{Y}(K)$ found above to the image $Y^{\prime}$, we obtain an element of the maximal ideal $\mathfrak{m}_{p}^{\prime}$ of $\mathcal{O}_{Y^{\prime}}$ at $p$, that we denote by $x$. On $Y$, this element has a zero of order one at $p$, and therefore it is a local generator for the maximal ideal $\mathfrak{m}_{p}$ of $\mathcal{O}_{Y}$. Let $\mathcal{O}^{\prime}$ and $\mathcal{O}$ be the local rings at $p$. We apply the Local Nakayama Lemma 5.1.19, regarding $\mathcal{O}$ as an $\mathcal{O}^{\prime}$-module.

We substitute $V=\mathcal{O}$ and $M=\mathfrak{m}_{p}^{\prime}$ into the statement of that lemma. Since $x$ is in $\mathfrak{m}_{p}^{\prime}, V / M V=\mathcal{O} / \mathfrak{m}_{p}^{\prime} \mathcal{O}$ is the residue field $k(p)$ of $\mathcal{O}$, which is spanned, as $\mathcal{O}^{\prime}$-module, by the element 1 . The Local Nakayama Lemma tells us that $\mathcal{O}$ is spanned, as $\mathcal{O}^{\prime}$-module, by 1 , and this shows that $\mathcal{O}=\mathcal{O}^{\prime}$.

## some curves of low genus

Here $Y$ will denote a smooth projective curve of genus $g$.

## curves of genus 2 .

When the genus of a smooth projective curve $Y$ is 2 , then $2 g-2=2$. The canonical map $\kappa$ is a map of degree 2 from $Y$ to $\mathbb{P}^{1}$. Every smooth projective curve of genus 2 is hyperelliptic.

## curves of genus 3.

Let $Y$ be a smooth projective curve of genus $g=3$. The canonical map $\kappa$ is a morphism of degree 4 from $Y$ to $\mathbb{P}^{2}$. If $Y$ isn't hyperelliptic, its image will be a plane curve of degree 4 , isomorphic to $Y$. The genus of a smooth projective curve of degree 4 is $\binom{3}{2}=3$ 1.7.29, which checks.

There is a second way to arrive at the same result. We go through it because the same method can be used for curves of genus 4 or 5 .

Riemann-Roch determines the dimension of the space of global sections of $\mathcal{O}(d K)$ :

$$
\mathbf{h}^{1} \mathcal{O}(d K)=\mathbf{h}^{0} \mathcal{O}((1-d) K)=0
$$

when $d>1$. For such $d$,

$$
\begin{equation*}
\mathbf{h}^{0} \mathcal{O}(d K)=\operatorname{deg}(d K)+1-g=d(2 g-2)-(g-1)=(2 d-1)(g-1) \tag{8.8.27}
\end{equation*}
$$

OdK
In our case $g=3$, so $\mathbf{h}^{0} \mathcal{O}(d K)=4 d-2$ when $d>1$.
The number of monomials of degree $d$ in $n+1$ variables $x_{0}, \ldots, x_{n}$ is $\binom{n+d}{d}$. Here $n=2$, so the number is $\binom{d+2}{2}$.

We assemble this information into a table:

| $d$ | 0 | 1 | 2 | 3 | 4 | 5 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\operatorname{monos} \operatorname{deg} d$ | 1 | 3 | 6 | 10 | 15 | 21 |
| $\mathbf{h}^{0} \mathcal{O}(d K)$ | 1 | 3 | 6 | 10 | 14 | 18 |

Now, if $\left(\alpha_{0}, \ldots, \alpha_{2}\right)$ is a basis of $H^{0} \mathcal{O}(K)$, the products $\alpha_{i_{1}} \cdots \alpha_{i_{d}}$ of length $d$ are global sections of $\mathcal{O}(d K)$. In fact, they generate the space $H^{0} \mathcal{O}(d K)$ of global sections. This isn't very important here, so we omit the proof. What we see from the table is that there is at least one homogeneous polynomial $f\left(x_{0}, \ldots, x_{2}\right)$ of degree 4 such that $f(\alpha)=0$. This means that the curve $Y$ lies in the zero locus of that polynomial, which is a quartic curve.

You will be able to show that, in fact, $Y$ is this quartic curve, so $f$ is, up to scalar factor, the only homogeneous quartic that vanishes on $Y$. Therefore the monomials of degree 4 in $\alpha$ span a space of dimension 14, and therefore they span $H^{0} \mathcal{O}(4 K)$. This is one case of the general fact that was stated above.

The table also shows that there are (at least) three independent polynomials of degree 5 that vanish on $Y$. They don't give new relations because we know three such polynomials, namely $x_{o} f, x_{1} f, x_{2} f$.

## curves of genus 4 .

When $Y$ is a smooth projective curve of genus 4 and is not hyperelliptic, the canonical map embeds $Y$ as a curve of degree $2 g-2=6$ in $\mathbb{P}^{3}$. Let's leave the analysis of this case as an exercise.

## curves of genus 5 .

With genus 5, things start to become more complicated.

Let $Y$ be a smooth projective curves of genus 5 that isn't hyperelliptic. The canonical map embeds $Y$ as a curve of degree 8 in $\mathbb{P}^{4}$. We make a computation analogous to what was done for genus 3 .

For $d>1$, the dimension of the space of global sections of $\mathcal{O}(d K)$ is

$$
\mathbf{h}^{0} \mathcal{O}(d K)=(2 d-1)(g-1)=8 d-4
$$

The number of monomials of degree $d$ in 5 variables is $\binom{d+4}{4}$.
We form a table:

| $d$ | 0 | 1 | 2 | 3 |
| :---: | :---: | :---: | :---: | :---: |
| monos deg $d$ | 1 | 5 | 15 | 35 |
| $\mathbf{h}^{0} \mathcal{O}(d K)$ | 1 | 5 | 12 | 20 |

This table predicts that there are (at least) three independent homogeneous quadratic polynomials $q_{1}, q_{2}, q_{3}$ that vanish on the curve $Y$.

Let $H_{1}, H_{2}, H_{3}$ be hypersurfaces in $\mathbb{P}^{4}$, of degrees $r_{1}, r_{2}, r_{3}$, respectively. If the intersection $Z=H_{1} \cap$ $H_{2} \cap H_{3}$ has dimension 1, it is called a complete intersection. Bézout's Theorem has a generalization that applies here. If $Z$ is a complete intersection, its degree, the number of intersections of $Z$ with a generic hyperplane, will be the product $r_{1} r_{2} r_{3}$. We set the proof up below, leaving details as an exercise. It is analogous to the proof of the usual Bézout's Theorem.

Thus, when $Q_{i}$ is the quadric $q_{i}=0$, the intersection $Z=Q_{1} \cap Q_{2} \cap Q_{3}$ will have degree $2^{3}=8$ and it will contain $Y$, also of degree 8 . Then $Y$ will be equal to $Z$, and will be complete intersection.

However, it is possible that the three quadrics intersect in a subset of dimension 2, and in this case the canonically embedded curve $Y$ isn't a complete intersection.

A curve which can be obtained as a three-sheeted covering of $\mathbb{P}^{1}$ is called a trigonal curve, another peculiar term.
8.8.28. Proposition. A trigonal curve of genus 5 is not isomorphic to an intersection of three quadrics in $\mathbb{P}^{4}$.
proof. A trigonal curve $Y$ will have a degee three morphism to the projective line: $Y \rightarrow X=\mathbb{P}^{1}$. Let's suppose that the point at infinity of $X$ isn't a branch point. Let the fibre over the point at infinity be $\left\{p_{1}, p_{2}, p_{3}\right\}$. With coordinates $\left(x_{0}, x_{1}\right)$ on $X$, the rational function $u=x_{1} / x_{0}$ on $X$ has poles $D=\sum p_{i}$ on $Y$, so $H^{0}(Y, \mathcal{O}(D))$ contains 1 and $u$, and therefore $\mathbf{h}^{0} \mathcal{O}(D) \geq 2$. By Riemann-Roch, $\chi \mathcal{O}(D)=3+1-g=-1$. so $\mathbf{h}^{1} \mathcal{O}(D)=\mathbf{h}^{0} \mathcal{O}(K-D) \geq 3$. There are three independent global sections of $\mathcal{O}(K)$ that vanish on $D$. Let them be $\alpha_{0}, \alpha_{1}, \alpha_{2}$. Then, when $Y$ is embedded into $\mathbb{P}^{4}$ by a basis $\left(\alpha_{0}, \ldots, \alpha_{4}\right)$, the three planes $\left\{x_{i}=0\right\}$, $i=0,1,2$ contain $D$. The intersection of these planes is a line $L$ that contains the three points $D$.

We go back to the three quadrics $Q_{1}, Q_{2}, Q_{3}$ that contain $Y$. Since they contain $Y$, they contain $D$. A quadric $Q$ intersects the line $L$ in two points unless $L \subset Q$. Therefore each of the quadrics $Q_{i}$ contains $L$, and $Q_{1} \cap Q_{2} \cap Q_{3}$ contains $L$ as well as $Y$. According to Bézout, the intersection of three quadrics, if it is one-dimensional, will have degree $2^{3}=8$. But $Y$ has degree 8 and $L$ has degree 1 , and $8+1=9$ is too big. Therefore $Q_{1} \cap Q_{2} \cap Q_{3}$ has dimension 2 .

It can be shown that this is the only exceptional case. A curve of genus 5 is either hyperelliptic, or trigonal, or else it is a complete intersecton of three quadrics in $\mathbb{P}^{4}$. However, this requires more work.
beztwo 8.8.30. Theorem. Let $H_{1}, H_{2}, H_{3}$ be hypersurfaces in $\mathbb{P}^{r}$ of degrees $\ell . m, n$, such that the intersection $Z=H_{1} \cap H_{2} \cap H_{3}$ has dimension $r-3$. Then the degree of $Z$ is the product $\ell m n$.

We slice $\mathbb{P}^{r}, H_{i}$, and $Z$ with a generic linear subspace $L$ of $\mathbb{P}^{r}$ of dimension 3. So $L$ is a projective space of dimension 3. The degree of $Z$ is defined to be the number of points of the intersection $L \cap Z$, and $L \cap H_{i}$ are surfaces in $\mathbb{P}^{3}$, of degrees $\ell, m, n$. We replace $\mathbb{P}^{r}$ by $L, H_{i}$ by $H_{i} \cap Z$, and $Z$ by $L \cap Z$. This reduces us to proving the theorem when $r=3$. We restate the theorem for that case.
8.8.31. Theorem. Let $H_{1}, H_{2}, H_{3}$ be surfaces in $\mathbb{P}^{3}$ of degrees $\ell . m, n$, respectivly, such that the intersection $Z=H_{1} \cap H_{2} \cap H_{3}$ has dimension zero. The number of intersecton points, counted with multiplicity, is the product $\ell m n$.

Setup for the proof. Say that the surfaces $H_{i}$ are the zero loci of the homogenous polynomials $f, g, h$, respectively. Let $\mathcal{O}$ be the structure sheaf on $\mathbb{P}^{3}$. We form the quotient $\overline{\mathcal{O}}=\mathcal{O} /(f, g, h)$. Its suport will be the finite set $Z$. We want to show that the dimension of $\overline{\mathcal{O}}$ as vector space is the product $\ell m n$.

The notation $\overline{\mathcal{O}}=\mathcal{O} /(f, g, h)$ isn't precise. We should really write $\mathcal{O}=\mathcal{O} / \mathcal{I}$, where $\mathcal{I}$ is the ideal $\mathcal{O}(-\ell) f+\mathcal{O}(-m) g++\mathcal{O}(-n) h$.

We form a resolution of $\overline{\mathcal{O}}$ analogous to the resolution 7.8.3:
$0 \rightarrow \mathcal{O}(-\ell-m-n) \xrightarrow{A} \mathcal{O}(-m-n) \oplus \mathcal{O}(-\ell-n) \oplus \mathcal{O}(-\ell-m) \xrightarrow{B} \mathcal{O}(-\ell) \oplus \mathcal{O}(-m) \oplus \mathcal{O}(-n) \xrightarrow{C} \mathcal{O} \rightarrow \overline{\mathcal{O}} \rightarrow 0$ where $A=\left(\begin{array}{l}f \\ g \\ h\end{array}\right), B=\left(\begin{array}{ccc}0 & h & -g \\ -h & 0 & f \\ g & -f & 0\end{array}\right)$, and $C=\left(\begin{array}{lll}f & g & h\end{array}\right)$.
The maps are given by matrix multiplication, with $A, B, C$ operating on the left, and with the direct sums represented as columns.

It is an exercise to show that this sequence is exact provided that the intersection has dimension zero. It shows that $\chi(\overline{\mathcal{O}})$ is independent of $f, g, h$. When $f, g, h$ are generic linear polynomials, there is just one point of intersection, the theorem is true in that case, and it follows when $f, g, h$ are products of generic linear polynomials. Therefore it is true in all cases. (Or, one can evaluate the alternating sum of the Euler characteristic of the terms.)

## Index

$S$-fractions, 51, 105
$\mathcal{O}$-module, 123
Bézout, Etienne (1730-1783), 157
Birkhoff, George David (1884-1944), 170
Chevalley, Claude (1909-1984), 94
Grothendieck, Alexander (1928-2014), 170, 177
Artin, Emil (1898-1962), 177
Bézout, Etienne (1730-1783), 32
Betti, Enrico (1823-1892), 26
Borel, Emile (1871-1956), 60
Dürer, Albrecht (1471-1528), 8
Desargues, Girard (1591-1661), 8
Grassmann, Hermann (1809-1877),75
Hausdorff, Felix (1868-1942), 60
Heine, Eduard (1821-1881), 60
Hensel, Kurt 1861-1941, 29
Hesse, Otto (1811-1877), 15
Hilbert, David (1862-1943), 39,44
Jacobi, Carl Gustav Jacob 1804-1851, 30
Möbius, August Ferdinand (1790-1868), 8
Nakayama, Tadashi (1912-1964), 81
Noether, Emmy (1882-1935), 38, 42, 83
Plücker, Julius (1801-1868), 35
Schelter, William (1947-2001), 96
Segre, Corrado (1863-1924), 61
Serre, Jean-Pierre (1926- ), 141, 154, 177
Veronese, Giuseppe (1854-1917), 61
Zariski, Oscar (1899-1986),40
affine cone, 62
affine double plane, 97
affine hypersurface, 44
affine morphism, 149
affine open subset, 74
affine plane, 5
affine plane curve, 5
affine space, [5
affine variety, 42,73
algebaic dimension, 6
algebra, 17, 37
algebra generators, 37
algebraically dependent, independent, 16
annihilator, 134
arithmetic genus, 153, 154
ascending, descending chain conditions, 42
basis for a topology, 49
bidual, 20
bihomogeneous polynomial, 65
bilinear relations, 104
Birkhoff-Grothendieck Theorem, 170
blowup, 54, 71
blowup map, 28
branch locus, 97,99
branch point, 25, 26,162
branched covering, 25, 162
canonical divisor, 181
canonical map, 183
center of projection, 25, 71
characteristic properties of cohomology, 144
Chinese Remainder Theorem, 38, 134
classical topology, 12
closd set, open set, 42
coboundary map, 142
cohomological functor, 142,143
cohomology of a complex, 142
cohomology of an $\mathcal{O}$-module, 141
cohomology sequence, 141, 143
cokernel, 123
comaximal ideals, 38
commutative diagram, 38
commuting matrices, 46
compact space, 59
complex, 142
conic, 6,9
connected space, 42
constructible function, 119
constructible set, 113
contracted ideal, 50,106
contravariant functor, 122
coordinate algebra, 43,46
Correspondence Theorem, 38
covering diagram, 125
curve, 110
cusp, 28
decomposable element, 76
degree, 61
degree of a morphism to $\mathbb{P}^{n}, 181$
degree of an affine plane curve, 5
degree of an invertible module, 167
derivation, 171
diagonal, 70
differential, 171
direct image, 131
direct limit, 130
directed set, 130
discrete valuation, 107
discriminant, 23
divisor, 157, 165
divisor of a function, 165
divisor of a polynomial, 11
domain, 17, 37
double point, 28
dual $\mathcal{O}$-module, 164
dual curve, 18
dual plane, 18
effective divisor, 165
eigenvector, 81
elliptic curve, 179
Euler characteristic, 26, 168
Euler characteristic of a finite $\mathcal{O}$-module, 157
exact sequence, 103,123
extended ideal,50, 106
extension by zero, 132
extension of domains, 82
extension of scalars, 105
exterior algebra, 75
Fermat curve, 14
fibre dimension, 119
fibre of a map, 6
fibred product, 116
finite $\mathcal{O}$-module, 123
finite module, 39
finite morphism, 83, 94
finite-type algebra, 37
finitely generated ideal, 39
flex point, 14
function field, 66,122
function field module, 130
function field of an affine variety,52
general position, generic, 30
generators of an $\mathcal{O}$-module, 137
genus, 26, 168, 178
geometric genus, 154
good point, 68
graded algebra, 75
graph of a morphism, 70, 117
Grassmanian, 75
Hausdorff space, 60
Heine-Borel Theorem, 60
Hessian determinant, 16
Hessian divisor, 33
Hessian matrix, 15
Hilbert Basis Theorem, 39
Hilbert Nullstellensatz, 44, 45
homogeneous fraction, 72
homogeneous ideal, 62
homogeneous parts, 10
homogeneous polynomial, 9,96
homogenize, dehomogenize, 13
homomorphism of $\mathcal{O}$-modules, 123
hyperelliptic curve, 182
hypersurface, 61
ideal generators, 40
ideal in the structure sheaf, 127
image, 123
increasing, strictly increasing chain, 39
induced topology, 41
integral closure, 85
integral extension, 82
integral morphism, 83,94
intersection multiplicity, 11, 158
invariant element, 55
invertible $\mathcal{O}$-module, 166
irreducible polynomial, 6
irreducible space, 42
irregularity, 154
irrelevant ideal, 62
isolated point, 13
isomorphism, 69
isomorphism of affine varieties, 52
kernel, 123
Laurent, Pierre Alphonse (1813-1854),49
line, 6, 7, 41
line at infinity, 7
linear map, 38
linear subspace, 60
linearly equivalent divisors, 166
local domain of dimension one, 107
local property, 109
local ring, 106
local ring at a point, 107, 109
local unit, 162
localization of an algebra, 49
locally closed set, 113
locally free module, 109
locally principal ideal, 109
locus of zeros, 5
maximal ideal, 37,127
maximal member, 39
module of differentials, 171
morphism, 68
morphism of affine varieties, 52
morphism of families, 125
multiplicative system, 51, 105
multiplicity, 10
multiplicity of zero,27
Nagata, Masayoshi (1927-2008),44
nilpotent ideal,48
nilradical,48
node, 28
noetherian ring, 39
noetherian space, 42
normal domain, normal variety, 84
normalization, 84
open covering, 60
order of zero, pole, 107
order of zero, pole of a differential, 174
ordinary bitangent, 34
ordinary curve, 35
orientability, 25
Pierre de Fermat (1607-1665), 14
Plücker Formulas, 35
plane projective curve, 11
pofint with values in a field, 67
point at infinity, 7
presentation of a module, 163
prime ideal, 37
principal divisor, 166
product equations, 29
product ideal, 37
product topology, 64
projection, 25,71
projective double plane, 99
projective line, 7
projective plane, 7
projective space, 7
projective variety, 59
proper variety, 117
quadric, 60
quasicompact, 60
radical of an ideal, 40
Rainich, George Yuri (1886-1968), 45
ramification index, 162
rank, 109
rational curve, 98,179
rational function on a variety, 66
rational function on an affine variety, 52
real projective plane, 8
reducible curve, 11
regular differential, 174
regular function, 52, 66, 122
regular function on an affine variety, 47
regular homogeneous fraction, 134
residue field, 44
residue field module, 127
residue field of a local ring, 106
resolution, acyclic resolution, 147
restriction of a divisor to an open set, 165
restriction of a morphism, 94
restriction of an $\mathcal{O}$-module, 133
restriction of scalars, 106
resultant, 22
resultant matrix, 22
scaling, 9
section, 122, 123
sections equal on an open set, 124
Segre embedding, 61
semicontinuous function, 119
Serre dual, 177
short exact sequence, 103
simple localization, 51, 105
smooth curve, singular curve, 14
smooth point, singular point, 14,110
special line, 27
special linear group, 44
spectrum of an algebra, 46
square-free polynomial, 86
standard affine open set, 8
standard cusp, 28
Strong Nullstellensatz, 45
structure sheaf, 121
support of a divisor, 165
support of a module, 134
tacnode, 28
tangent line, 14
tensor algebra, 78
tensor product module, 104
torsion, 50, 127
torsion-free, 50
trace of a differential, 174
trace of a field extension, 86
transcendence basis, 17
transcendental element, 17
transversal intersection, 32
trefoil knot, 28
triangulation, 26
trigonal curve, 185
twist of an $\mathcal{O}$-module, 136
twisted cubic, 62
twisting module, 136
unit ideal, 12
valuation, 107
valuation ring, 107
value of a function, 47
variety, 63
weighted degree, 23
weighted projective space, 99
Zariski closed set, 59
Zariski closed, Zariski open,40
Zariski topology, 41, 59
zero of a polynomial, 11,40


[^0]:    ${ }^{1}$ While writing a paper, the mathematician Nagata decided that the English language needed this unusual word, and he managed to find it in a dictionary.

