# MATH3195/5195M: Commutative rings and algebraic geometry 

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## Introduction

Some useful books:

- Miles Reid - Undergraduate algebraic geometry, LMS Student Texts 12, CUP, 1988.
- Miles Reid - Undergraduate commutative algebra, LMS Student Texts 29, CUP, 1995.
- M.F. Atiyah and I.G. MacDonald - Introduction to commutative algebra, Westview Press, 1994
- David Cox, John Little, and Donal O'Shea - Ideals, Varieties, and Algorithms, UTM Springer, Fourth Edition, 2015.
- Rodney Sharp - Steps in commutative algebra 2nd Ed, LMS Student Texts 51, CUP, 2000.
- Robin Hartshorne - Algebraic Geometry, Springer Verlag, 1997. (First chapter only)
- W. Fulton - Algebraic Curves.


## A bit of history

Commutative algebra has its origins in number theory and geometry. On the other hand, it is the foundation of modern algebraic geometry and complex analytic geometry.

The most basic commutative rings are the integers $\mathbb{Z}$ and the polynomial ring $k[x]$ over a field $k$. We will also encounter these rings frequently.

Commutative algebra was probably started by Dedekind, who coined the notion of an ideal in $\mathbb{Z}$ (around 1870). Ideals are a generalization of prime elements. David Hilbert introduced the notion of a ring. A few years later, 1890, he proved his famous basis theorem, that says that every ideal in polynomial ring (over a field) is finitely generated (this will be proven in the course). Later on, in the 1920s, Emmy Noether studied the ascending chain condition on commutative rings (we will work a lot with Noetherian rings). This was in some sense the birth of modern abstract algebra. The 1930s saw developments of dimension theory of commutative rings, as well as the concepts of localization and completion (mostly by the German mathematician Felix Krull).

In the 1940 geometry enters the picture, with work by Claude Chevalley and Oscar Zariski: they applied the formal language of modern abstract algebra to algebraic geometry. The next milestone for algebraic geometry came already in the 1960s, when Alexander Grothendieck develeoped the language of schemes that revolutionized our understanding of algebraic geometry.

Since then, there are many different directions of research in commutative algebra and algebraic geometry, from the abstract (homological methods) to computational commutative algebra (Gröbner bases techniques). We mention a few more important results: Heisuke Hironaka proved resolution of singularities in 1964, Michael Artin proved the approximation theorem 1969. From the 1970s on homological methods became popular (e.g via work of Auslander, Buchsbaum, Northcott, Rees, Eisenbud, and of course, Serre). Melvin Hochster formulated the homological conjectures in 1970, which are still a major object of study. One of them, the direct summand conjecture, was only recently proven in 2017 by Yves André using the machinery of Scholze's perfectoid spaces.

## Some examples

Here we give a few examples of typical problems in commutative algebra and algebraic geometry.

## Divisibility

Example. Consider the polynomial ring $\mathbb{C}[x]$ in one variable with coefficients in the complex numbers and let

$$
f(x)=x^{4}-2 x^{2}+1 \quad \text { and } g(x)=x^{3}+x^{2}-2 x .
$$

Question: Do $f(x)$ and $g(x)$ have a common factor? (Equivalently: Do the two functions have a common 0? or: Do the two subsets $\{x \in \mathbb{C}: f(x)=0\}$ and $\{x \in \mathbb{C}: g(x)=0\}$ of $\mathbb{C}$ have nonempty intersection?)
Solution: Using Euclid's algorithm, we can write $f(x)=g(x)(x-1)+\left(x^{2}-2 x+1\right)$, and further with $r_{1}(x):=x^{2}-2 x+1$ we get $g(x)=r_{1}(x)(x+3)+(3 x-3)$. We find that $r_{2}(x):=3 x-3$ divides $r_{1}(x)$ and thus $(x-1)$ is a common factor of $f(x)$ and $g(x)$. Geometrically, this means that the set $X:=\{x \in \mathbb{C}: f(x)=g(x)=0\}$ in $\mathbb{C}$ is nonempty, more precisely $X=\{1\}$.

Example. Consider the three polynomials in $\mathbb{C}[x, y]$ :

$$
f(x, y):=x^{3}-y^{2}, g(x, y):=x+y \text { and } h(x, y):=x-y .
$$

Do these three polynomials have a common "factor"? It is not quite clear how to factor polynomials in several variables, but we can still ask the geometric question: do the zerosets $X_{1}:=\{(x, y) \in$ $\left.\mathbb{C}^{2}: f(x, y)=0\right\}, X_{2}:=\left\{(x, y) \in \mathbb{C}^{2}: g(x, y)=0\right\}$ and $X_{3}:=\left\{(x, y) \in \mathbb{C}^{2}: h(x, y)=0\right\}$ have a nonempty intersection in $\mathbb{C}^{2}$ ? More compactly: is $X_{1} \cap X_{2} \cap X_{3}=\left\{(x, y) \in \mathbb{C}^{2}: f(x, y)=\right.$ $g(x, y)=h(x, y)=0\}$ equal to $\varnothing$ ? We can even become more greedy and ask to find all solutions of this system of equations in $\mathbb{C}^{2}$, or ask about the size of solutions (which has to be defined in a suitable way!).
In this example one easily finds that $X_{1} \cap X_{2} \cap X_{3}=\{(0,0)\}$ has only one solution. In due course we will learn about some more general algebraic techniques, involving so-called Gröbner bases, to solve systems of polynomial equations in several variables. Algebraically these will be questions about ideals in polynomial rings.

## Geometry

In the last example we have already alluded to the fact that solutions of systems of polynomial equations can be interpreted geometrically. Let us introduce some terminology: Let $K$ be a field and $K\left[x_{1}, \ldots, x_{n}\right]$ be the polynomial ring over $K$ in $n$ variables. A polynomial $P\left(x_{1}, \ldots, x_{n}\right)=$ $\sum_{\alpha \in \mathbb{N}^{n}} a_{\alpha} \underline{x}^{\alpha}$, where $\underline{x}^{\alpha}=x_{1}^{\alpha_{1}} \cdots x_{n}^{\alpha_{n}}$ and $a_{\alpha} \in K$, gives a function $P: K^{n} \rightarrow K,\left(a_{1}, \ldots, a_{n}\right) \mapsto$ $P\left(a_{1}, \ldots, a_{n}\right)$. For example, the polynomial $P(x, y):=x^{3}-y^{2} \in \mathbb{R}[x, y]$ evaluates to 0 for the points $(0,0),(1,1),\left(\frac{1}{4}, \frac{1}{8}\right)$.

Given polynomials $P_{1}(\underline{x}), \ldots, P_{k}(\underline{x}) \in K\left[x_{1}, \ldots, x_{n}\right]$, one defines

$$
\mathbb{V}\left(P_{1}, \ldots, P_{k}\right)=\left\{\left(a_{1}, \ldots, a_{n}\right) \in K^{n}: P_{i}\left(a_{1}, \ldots, a_{n}\right)=0 \text { for all } i=1, \ldots, k\right\}
$$

Sets $X \subseteq K^{n}$ of the form $X:=\mathbb{V}\left(P_{1}, \ldots, P_{k}\right)$ are called algebraic sets and we will see in the course that they reflect algebraic properties of so-called ideals in the polynomial ring $K\left[x_{1}, \ldots, x_{n}\right]$. For example, we have $\mathbb{V}\left(P_{1}, P_{2}\right)=\mathbb{V}\left(P_{1}\right) \cap \mathbb{V}\left(P_{2}\right)$ and $\mathbb{V}\left(P_{1} \cdot P_{2}\right)=\mathbb{V}\left(P_{1}\right) \cup \mathbb{V}\left(P_{2}\right)$.
In particular: finding $\mathbb{V}\left(P_{1}, \ldots, P_{k}\right)$ is equivalent to solving the system of polynomial equations $\left\{\underline{x} \in K^{n}: P_{1}(\underline{x})=\cdots=P_{k}(\underline{x})=0\right\}$.

## Solutions over different rings

Example 0.1. What are the integer solutions of $X^{2}+Y^{2}=Z^{2}$ ? Solutions certainly exist, for example

$$
(X, Y, Z)=(3,4,5)
$$

but are there others? Can we find all of them?
Note first that if $Z=0$ then both $X$ and $Y$ must also be zero. Assuming henceforth then that $Z \neq 0$ we substitute $x=\frac{X}{Z}$ and $y=\frac{Y}{Z}$ we can re-frame the question as

What are the rational solutions of $x^{2}+y^{2}=1$ ?
Think of the solution set as a circle in $\mathbb{R}^{2}$ :


Consider a line of slope $t$ through $(0,1)$ that rotates about $(0,1)$. We can then find all solutions of our new equation by using $t$ as a new parameter. Note that we will obtain any rational point on the circle except $(0,-1)$, which would correspond to $t=\infty$.


So we want rational solutions of

$$
\left.\begin{array}{r}
y-t x=1 \\
x^{2}+y^{2}=1
\end{array}\right\}
$$

Substituting the first of these into the second we see

$$
\begin{aligned}
x^{2}+(t x+1)^{2}=1 & \Longrightarrow x^{2}+t^{2} x^{2}+2 t x+1=1 \\
& \Longrightarrow x^{2}\left(t^{2}+1\right)+2 t x=0 \\
& \Longrightarrow x\left(x\left(t^{2}+1\right)+2 t\right)=0
\end{aligned}
$$

This gives two solutions, $x=0$ and $x=\frac{-2 t}{t^{2}+1}$. The first solution for $x$ gives $y=1$, and the second gives

$$
y=\frac{-2 t^{2}}{t^{2}+1}+1=\frac{1-t^{2}}{1+t^{2}}
$$

Note that also $t=0$ gives $y=1$. All rational points on the circle are therefore

$$
(x, y)=(0,-1) \text { and }(x, y)=\left(\frac{-2 t}{1+t^{2}}, \frac{1-t^{2}}{1+t^{2}}\right) \quad(t \in \mathbb{R})
$$

We need to see which values of $t$ give rational values of $x$ and $y$. A bit of checking shows that $x, y \in \mathbb{Q} \Longleftrightarrow t \in \mathbb{Q}$. So let $t=\frac{m}{n}$, where $m$ and $n$ are coprime integers. Then

$$
x=\frac{-2 m n}{m^{2}+n^{2}} \text { and } y=\frac{n^{2}-m^{2}}{m^{2}+n^{2}}
$$

Returning to our original variables $X, Y$ and $Z$ we see that integer solutions to $X^{2}+Y^{2}=Z^{2}$ can be given by

$$
\begin{gathered}
Y=2 m n, Y=n^{2}-m^{2}, Z=m^{2}+n^{2}, \quad m, n \in \mathbb{Z}, m, n \text { coprime, or } \\
X=m n, Y=\frac{n^{2}-m^{2}}{2}, Z=\frac{m^{2}+n^{2}}{2} \quad \text { if both } m \text { and } n \text { are odd. }
\end{gathered}
$$

For instance, $m=1, n=3$ gives $X=3, Y=4, Z=5$.

## Parametrization of algebraic varieties

Similar to linear algebra, where one finds parametrizations of linear sub-spaces of $K^{n}$, one may want to find a parametrization of an algebraic set $X \subseteq K^{n}$. This may not be possible in general, but sometimes one can use known parametrizations of so-called smooth algebraic varieties to construct parametrizations of more complicated ones.
For example, let $X:=\mathbb{V}\left(x^{4}+y^{2}-x^{2}\right) \subseteq \mathbb{R}^{2}$. This curve, determined by $f(x, y)=x^{4}+y^{2}-x^{2}$ is called lemniscate. How can we find a parametrization of $X$ ?
We make the following observation: consider the map $\pi: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}$ sending a point $(a, b) \mapsto$ $(a, a b)$. Then

$$
f(\pi(x, y))=f(x, x y)=x^{4}+x^{2} y^{2}-x^{2}=x^{2}\left(x^{2}+y^{2}-1\right) .
$$

Using that $\mathbb{V}\left(x^{2}\left(x^{2}+y^{2}-1\right)\right)=\mathbb{V}\left(x^{2}\right) \cup \mathbb{V}\left(x^{2}+y^{2}-1\right)$, we see that $\mathbb{V}(f(\pi(x, y))$ is the union of a circle and a line. Of course one can parametrize the circle with the standard parametrization $x=\cos t$ and $y=\sin t$ (or, if one wants to avoid transcendental functions, with the rational parametrization from the example above). But then we see, that $X$ is parametrized by all points $(x, x y)$ that satisfy $f(x, x y)=0$, and thus we obtain the parametrization $(\cos t, \cos t \cdot \sin t)$ of $X$.

Although this construction seems ad-hoc, $\pi$ is an example of a blow up map, that actually is a so-called resolution of singularities of $X$.

## Part I

## Commutative Algebra

## 1 Revision of rings

Definition 1.1. A ring is a triple $(R,+, \cdot)$ of a set $R$ and two binary operations

$$
\begin{aligned}
+: R \times R \longrightarrow R & \text { (addition) } \\
\cdot: R \times R \longrightarrow R & \text { (multiplication) }
\end{aligned}
$$

such that the following hold:
(i) $(R,+)$ is an abelian group, with identity $0=0_{R}$;
(ii) there is an element $1=1_{R}$ such that $1 \cdot r=r \cdot 1=r$ for all $r \in R$;
(iii) $\cdot$ is associative, i.e. $(r \cdot s) \cdot t=r \cdot(s \cdot t)$ for all $r, s, t \in R$;
(iv) $\cdot$ distributes over + , i.e. $r \cdot(s+t)=r \cdot s+r \cdot t$ and $(s+t) \cdot r=s \cdot r+t \cdot r$ for all $r, s, t \in R$.

We will often abbreviate the triple $(R,+, \cdot)$ to just $R$ with the operations implicit, and moreover the multiplication $r \cdot s$ to just $r s$.

Definition 1.2. A ring $R$ is called commutative if $r s=s r$ for all $r, s \in R$.
Remark. In this course all rings will be commutative rings, and so hereafter we will take "ring" to mean "commutative ring".
Example 1.3. (i) $\mathbb{Z}$, the set of integers.
(ii) $\mathbb{Z}_{n}=\mathbb{Z} / n \mathbb{Z}$, the integers modulo $n$.
(iii) $\mathbb{R}$, the set of real numbers.
(iv) $\mathbb{C}$, the set of complex numbers.
(v) $\mathcal{C}[0,1]$, the set of continuous functions on $[0,1]$.
(vi) Gaussian integers $\mathbb{Z}[i]=\{a+b i: a, b \in \mathbb{Z}\}$.
(vii) Let $X$ be any set, and define $\mathfrak{F}_{X}=\mathbb{R}^{X}=\{$ functions $f: X \longrightarrow \mathbb{R}\}$. Define $+, \cdot: \mathfrak{F}_{X} \times \mathfrak{F}_{X} \longrightarrow$ $\mathfrak{F}_{X}$ by

$$
\begin{aligned}
(f+g): X & \rightarrow \mathbb{R} \\
x & \mapsto f(x)+g(x) \\
(f \cdot g): X & \rightarrow \mathbb{R} \\
x & \mapsto f(x) g(x)
\end{aligned}
$$

Then $\mathfrak{F}_{X}$ is a commutative ring, with additive identity $0_{\mathfrak{F}_{X}}: x \mapsto 0$ and multiplicative identity $1_{\tilde{F}_{X}}: x \mapsto 1$.
(viii) We can also construct new rings from old ones. Let $R$ be any commutative ring, and define

$$
R[x]=\{\text { polynomials in } x \text { with coefficients in } R\}=\left\{\sum_{i=0}^{n} r_{i} x^{i}: n \in \mathbb{N} \text { and } r_{i} \in R \forall i\right\}
$$

This is also a commutative ring. We can then define $R\left[x_{1}, \ldots, x_{n}\right]$ inductively by

$$
R\left[x_{1}, \ldots, x_{n}\right]=R\left[x_{1}, \ldots, x_{n-1}\right]\left[x_{n}\right] .
$$

This is just polynomials in the variables $x_{1}, \ldots, x_{n}$ with coefficients in $R$.
(ix) $R[[x]]=\{$ formal power series in $x$ with coefficients in $R\}=\left\{\sum_{i=0}^{\infty} r_{i} x^{i}: r_{i} \in R \forall i\right\}$. Note that these are formal objects, not necessarily functions from $R$ to $R$. For instance, $\sum_{i=0}^{\infty} x^{i}$ is an element of $\mathbb{R}[[x]]$, but we cannot evaluate this at $x=1$ so it does not define a function $\mathbb{R} \rightarrow \mathbb{R}$.

Definition 1.4. A field is a ring $K$ where every element other than $0_{K}$ has a multiplicative inverse. Formally, for each $r \in K \backslash\{0\}$ there exists an $r^{-1} \in K \backslash\{0\}$ such that $r r^{-1}=r^{-1} r=1_{K}$.

Example 1.5. (i) Familiar fields are $\mathbb{C}, \mathbb{R}, \mathbb{Q}$. Another example is $\mathbb{Z}_{p}=\mathbb{Z} / p \mathbb{Z}$ for any prime $p$.
(ii) $\mathbb{Z}$ itself is not a field, nor is the set $\mathbb{Z}[i]$ of Gaussian integers. For instance, $2+0 i$ has no inverse. In fact the units of $\mathbb{Z}[i]$ are $\pm 1, \pm i$.

We will now see another way of constructing rings and fields from old ones:
Example 1.6. Let $R, S$ be rings. The Cartesian product $R \times S=(R \times S,+, \cdot)$ of $R$ and $S$ is also a ring, where we define

$$
\begin{aligned}
\left(r_{1}, s_{1}\right)+\left(r_{2}, s_{2}\right) & =\left(r_{1}+r_{2}, s_{1}+s_{2}\right) \\
\left(r_{1}, s_{1}\right) \cdot\left(r_{2}, s_{2}\right) & =\left(r_{1} r_{2}, s_{1} s_{2}\right) .
\end{aligned}
$$

for all $r_{1}, r_{2} \in R, s_{1}, s_{2} \in S$. We have $0_{R \times S}=\left(0_{R}, 0_{S}\right)$ and $1_{R \times S}=\left(1_{R}, 1_{S}\right)$. Note that if $K$ and $L$ are fields then $K \times L$ is not a field, for instance $(0,1)$ has no multiplicative inverse.

Definition 1.7. A subset $S \subseteq R$ of a ring $R$ is called a subring if $(S,+)$ is a subgroup of $(R,+)$, $1_{R} \in S$ and $S$ is closed under multiplication. Similarly, if $K$ is a field then a subset $L \subseteq K$ is called a subfield if it is a subring of $K$ and $r^{-1} \in L$ for all non-zero $r \in L$.

Example 1.8. Let $R=\mathbb{R}$ and $S=\{a+b \sqrt{5}: a, b \in \mathbb{Z}\}$. Clearly $0=0+\sqrt{5}, 1=1+0 \sqrt{5} \in S$, so we will check that it is additively and multiplicatively closed. For all $a, b, c, d \in \mathbb{R}$, we have

$$
\begin{aligned}
(a+b \sqrt{5})+(c+d \sqrt{5}) & =(a+c)+(c+d) \sqrt{5} \in S \\
(a+b \sqrt{5})(c+d \sqrt{5}) & =a c+a d \sqrt{5}+b c \sqrt{5}+5 b d \\
& =(a c+5 b d)+(a d+b c) \sqrt{5} \in S
\end{aligned}
$$

Similarly if $R=\mathbb{C}$, then $S=\{a+b \sqrt{-5}: a, b \in \mathbb{Z}\}$ is a subring. Rings like these play an important role in areas of number theory.

Definition 1.9. Let $R, S$ be rings. A ring homomorphism from $R$ to $S$ is a map $\varphi: R \rightarrow S$ such that for all $r_{1}, r_{2} \in R$ :
(i) $\varphi\left(r_{1}+r_{2}\right)=\varphi\left(r_{1}\right)+\varphi\left(r_{2}\right)$;
(ii) $\varphi\left(r_{1} r_{2}\right)=\varphi\left(r_{1}\right) \varphi\left(r_{2}\right)$;
(iii) $\varphi\left(1_{R}\right)=1_{S}$.

If $\varphi$ is bijective then we say $\varphi$ is an isomorphism.
Exercise (Exercise sheet 0). If $\varphi: R \rightarrow S$ is a ring isomorphism, prove that $\varphi^{-1}: S \rightarrow R$ is a ring homomorphism (and hence also an isomorphism).

Definition 1.10. Let $\varphi: R \rightarrow S$ be a ring homomorphism. The kernel of $\varphi$, denoted $\operatorname{Ker} \varphi$, is the set

$$
\operatorname{Ker} \varphi=\left\{r \in R: \varphi(r)=0_{S}\right\} .
$$

The image of $\varphi$, denoted $\operatorname{Im} \varphi$, is the set

$$
\operatorname{Im} \varphi=\{\varphi(r): r \in R\} .
$$

The proof of the following proposition is left as an easy exercise:
Proposition 1.11. (i) $\operatorname{Im} \varphi$ is a subring of $S$.
(ii) $\operatorname{Ker} \varphi$ is not necessarily a subring of $R$.

Proof. Exercise.

## 2 Revision of ideals

That $\operatorname{Ker} \varphi$ is not a subring of $R$ causes us problems if we wish to introduce quotient rings like we introduced quotient groups. Note that if $H$ is a subgroup of $G$ then $G / H$ does not necessarily exist. Note also that dealing with commutative groups circumvents this problem, but that is not the case when dealing with rings. The "correct" notion of a substructure that allows us to take quotients is that of an ideal.

Definition 2.1. Let $R$ be a ring. A subset $I \subseteq R$ is called an ideal if:
(i) $I \neq \varnothing$;
(ii) for all $x, y \in I, x-y \in I$;
(iii) for all $x \in I$ and $r \in R, r x \in I$.

We write $I \subseteq R$ to mean $I$ is an ideal of the ring $R$.
If $I \neq R$, then we say that $I$ is a proper ideal of $R$.
Example 2.2. (i) Let $R$ be a ring. Then $\left\{0_{R}\right\}$ and $R$ are both ideals of $R$, usually referred to as trivial ideals.
(ii) For any $n \in \mathbb{Z}, n \mathbb{Z}$ is an ideal of $\mathbb{Z}$.
(iii) For a ring homomorphism $\varphi: R \rightarrow S, \operatorname{Ker} \varphi$ is an ideal of $R$. Indeed let $x, y \in \operatorname{Ker} \varphi$ and $r \in R$, then

$$
\begin{aligned}
\varphi(0) & =0 \text { so } 0 \in \operatorname{Ker} \varphi \quad(\operatorname{Ker} \varphi \neq \varnothing) \\
\varphi(x+y)=\varphi(x)+\varphi(y)=0+0 & =0 \text { so } x+y \in \operatorname{Ker} \varphi \\
\varphi(r x)=\varphi(r) \varphi(x)=\varphi(r) 0 & =0 \text { so } r x \in \operatorname{Ker} \varphi
\end{aligned}
$$

(iv) A crucial example for algebraic geometry, and one we will encounter many times later in the course, is the following. Let $K$ be a field (usually $\mathbb{R}$ or $\mathbb{C}$ ), $V \subseteq K^{n}$ be a set and $R=$ $K\left[X_{1}, \ldots, X_{n}\right]$. Then

$$
I(V)=\{f \in R: f(v)=0 \text { for all } v \in V\}
$$

is an ideal of $R$.

Definition 2.3. Let $A$ be a non-empty subset of a ring $R$. The ideal generated by $A$, denoted $\langle A\rangle$, is the set of all elements

$$
\langle A\rangle=\left\{\sum_{i=1}^{n} r_{i} a_{i}: n \in \mathbb{N}, r_{1}, \ldots, r_{n} \in R, a_{1}, \ldots, a_{n} \in A\right\}
$$

We say an ideal $I$ is finitely generated if there exists a finite subset $A \subseteq R$ such that $I=\langle A\rangle$. If $I=\langle a\rangle$ is generated by one element, then $I$ is called a principal ideal.

Example 2.4. Let $R=K[x, y, z]$, and $I=\langle x, y, z\rangle$. Then $I$ consists of all polynomials in $K[x, y, z]$ without constant term. One can show that $I=J$, where $J=\left\langle x+y, y+z^{2}, z\right\rangle$.

We can also perform operations on ideals as per the following proposition.
Proposition 2.5. Let $I, J$ be ideals of a ring $R$. The following are then also ideals of $R$ :
(i) $I \cap J=\{x: x \in I$ and $x \in J\}$, the intersection of $I$ and $J$;
(ii) $I J=\langle\{x y: x \in I, y \in J\}\rangle$, the product of $I$ and $J$;
(iii) $I+J=\langle I \cup J\rangle$, the sum of $I$ and $J$;
(iv) $(I: J)=\{r \in R: r J \subseteq I\}$, the ideal quotient of I and $J$.

Proof. Exercise. See Exercise Sheet 1.
In algebraic geometry the following type of ideals will play an important role:
Definition 2.6. Let $I \subseteq R$ be an ideal in a ring. Then

$$
\sqrt{I}:=\left\{x \in R: \text { there exists an } n \in \mathbb{N} \text { such that } x^{n} \in I\right\}
$$

is an ideal, called the radical of $I$. If $I=\sqrt{I}$, then $I$ is called a radical ideal.
See exercise sheet 1 for a proof that $\sqrt{I}$ is an ideal in $R$.
Example 2.7. (1) Let $I=288 \mathbb{Z}$ in $\mathbb{Z}$. Then $\sqrt{I}=6 \mathbb{Z}$ (see this from $288=2^{5} 3^{2}$ ), and so $I$ is not a radical ideal.
(2) Let $I=\left\langle x^{2}, y^{2}\right\rangle$ in $K[x, y]$. It is clear that $\sqrt{I} \supseteq\langle x, y\rangle$. For the other inclusion note that a polynomial $P(x, y)$ is in $\sqrt{I}$ if and only if there exists an $n$, such that $P^{n}(x, y)$ is in $I$, that is $P^{n}$ does not have a constant term. But $P(0,0)^{n}=0$ if and only if $P(0,0)=0$, thus $P$ itself must be without nonconstant term, thus $P(x, y) \in I$.

We will now move on to quotient rings.
Definition 2.8. Let $I$ be an ideal of a ring $R$. A coset of $I$ in $R$ is a set

$$
r+I=\{r+x: x \in I\}
$$

for some $r \in R$. This may also be denoted by $\bar{r}$, and we denote by $R / I$ the set of cosets of $I$ in $R$.
The following proposition is straightforward:
Proposition 2.9. (i) Two cosets are either equal or disjoint, and the union of all cosets is $R$. We say that the cosets partition $R$.
(ii) Cosets $r+I$ and $s+I$ are equal if and only if $r-s \in I$.
(iii) We can define multiplication and addition on $R / I$ by setting $(r+I)+(s+I)=(r+s)+I$ and $(r+I)(s+I)=r s+I$.
(iv) The additive and multiplicative identities of $R / I$ are $0+I=I$ and $1+I$ respectively.

This proposition shows that we have a ring structure on $R / I$, with much of the structure inherited from the ring structure on $R$.

Proposition 2.10. Let I be an ideal of a ring $R$. Define $\varphi: R \rightarrow R / I$ by $\varphi(r)=r+I$. Then:
(i) $\varphi$ is a ring homomorphism (called the quotient homomorphism);
(ii) $\operatorname{Ker} \varphi=I$;
(iii) there is a bijection between ideals of $R / I$ and the ideals of $R$ which contain $I$, given by

$$
\begin{aligned}
J \subseteq R / I & \longmapsto \varphi^{-1}(J)=\{r \in R: r+I \in J\} \\
I \subseteq K \subseteq R & \longmapsto \varphi(K)=\{r+I: r \in K\} .
\end{aligned}
$$

Proof. (i) See Exercise Sheet 1.
(ii) See Exercise Sheet 1.
(iii) For an ideal $K$ such that $I \subseteq K \subseteq R$, we first show that $\varphi(K)$ is an ideal of $R / I$ (note that this may not be true for any $\varphi$ ). Clearly $\varphi(K) \neq \varnothing$, as $\varphi(I)=I \in \varphi(K)$. For any two cosets $r+I, s+I \in \varphi(K)$ we have $r, s \in K$, and since $K$ is an ideal then $r-s \in K$. Hence $(r+I)-(s+I)=(r-s)+I \in \varphi(K)$. If now we also choose any $t+I \in R / I$ then $(t+I)(r+I)=t r+I \in \varphi(K)$, since $t r \in K$ again due to $K$ being an ideal of $R$.
We now show that the assignment $K \longmapsto \varphi(K)$ is injective. Suppose $K \neq K^{\prime}$ are both ideals of $R$ containing $I$, then without loss of generality there is some $r \in K$ such that $r \notin K^{\prime}$. We clearly have $r+I \in \varphi(K)$. We will show that $r+I \notin \varphi\left(K^{\prime}\right)$, thus $\varphi(K) \neq \varphi\left(K^{\prime}\right)$. Assume for a contradiction that $r+I \in \varphi\left(K^{\prime}\right)$, then $r+I=s+I$ for some $s \in K^{\prime}$. By the equality rule for cosets, we have $r-s \in I \subseteq K^{\prime}$, and hence $(r-s)+s=r \in K^{\prime}$, a contradiction.
Finally, we show the map $K \longmapsto \varphi(K)$ is surjective. Given an ideal $J \subseteq R / I$ we clearly have $\varphi\left(\varphi^{-1}(J)\right)=J$, so we must show that $\varphi^{-1}(J)$ is an ideal of $R$ containing $I$. The containment is easy, since $I=\varphi^{-1}(0) \subseteq \varphi^{-1}(J)$. If now $r, s \in \varphi^{-1}(J)$, then $r+I, s+I \in J$ and hence $(r-s)+I \in J$. Therefore $r-s \in \varphi^{-1}(J)$. Similarly if $t \in R$ then $t+I \in R / I$ and $(t+I)(r+$ $I)=\operatorname{tr}+I \in J$, hence $\operatorname{tr} \in \varphi^{-1}(J)$.

Theorem 2.11. Let $\varphi: R \rightarrow S$ be a ring homomorphism. Then $\bar{\varphi}: R / \operatorname{Ker} \varphi \rightarrow \operatorname{Im} \varphi$ given by $\bar{\varphi}(r+$ $\operatorname{Ker} \varphi)=\varphi(r)$ is an isomorphism.
Proof. See Exercise Sheet 1 (remember to check that this is well defined!).

## 3 Prime ideals

Definition 3.1. An ideal $\mathfrak{p}$ of $R$ is called a prime ideal if;
(i) $\mathfrak{p} \neq R$;
(ii) $x y \in P \Longrightarrow x \in \mathfrak{p}$ or $y \in P$.

The first example below explains the name of these ideals.
Example 3.2. (i) The ideal $n \mathbb{Z}$ of $\mathbb{Z}$ is prime if and only if either $n$ is prime or $n=0$ (Exercise).
(ii) The ideal $\langle f\rangle$ of $\mathbb{C}[x]$ is prime if and only if either $f=0$ or $f$ is irreducible, i.e. $f$ cannot be written as the product of two polynomials of positive degree.

Proposition 3.3. Let $\varphi: R \rightarrow S$ be a ring homomorphism. If $\mathfrak{p} \subseteq S$ is a prime ideal, then $\varphi^{-1}(\mathfrak{p}) \subseteq R$ is a prime ideal.

Proof. Let $x, y \in R$ be such that $x y \in \varphi^{-1}(\mathfrak{p})$, i.e. $\varphi(x y) \in \mathfrak{p}$. Now $\varphi(x y)=\varphi(x) \varphi(y)$, and since $\mathfrak{p}$ is prime we therefore have either $\varphi(x) \in \mathfrak{p}$ or $\varphi(y) \in \mathfrak{p}$. Hence either $x \in \varphi^{-1}(\mathfrak{p})$ or $y \in \varphi^{-1}(\mathfrak{p})$.
Proposition 3.4. Let $I$ be an ideal of a ring $R$. If $\mathfrak{p}$ is a prime ideal of $R$ containing $I$, then the image of $\mathfrak{p}$ in $R / I$ is also prime.

Proof. Denote by $\overline{\mathfrak{p}}$ the image of $\mathfrak{p}$ in $R / I$. Suppose $x+I, y+I \in R / I$ are such that $(x+I)(y+I) \in$ $\overline{\mathfrak{p}}$. Then $x y+I \in \overline{\mathfrak{p}}$, so there is some $p \in \mathfrak{p}$ such that $x y-p \in I \subseteq \mathfrak{p}$. Therefore $x y \in \mathfrak{p}$, so either $x \in \mathfrak{p}$ or $y \in \mathfrak{p}$ as $\mathfrak{p}$ is prime, thus either $x+I \in \overline{\mathfrak{p}}$ or $y+I \in \overline{\mathfrak{p}}$.

Remark 3.5. These two propositions show that the bijection between ideals of $R / I$ and ideals of $R$ containing $I$ restricts to a bijection between prime ideals of $R / I$ and prime ideals of $R$ containing I.

Definition 3.6. A ring $R$ is an integral domain if:
(i) $R \neq\{0\}$;
(ii) for all $r, s \in R, r s=0 \Longrightarrow r=0$ or $s=0$, i.e. there are no non-zero zero divisors.

Example 3.7. (i) $\mathbb{Z}$ and $K[x]$ are integral domains.
(ii) $R=K[x] /\left\langle x^{2}\right\rangle$ is not an integral domain, since $\bar{x} \neq \overline{0}$ in $R$ but $\bar{x} \cdot \bar{x}=\overline{0}$.
(iii) $\mathbb{Z}_{4}$ is not an integral domain, as $(2+4 \mathbb{Z})(2+4 \mathbb{Z})=4+4 \mathbb{Z}=0$.
(iv) $\mathbb{R}[x] /\left\langle x^{2}+1\right\rangle$ is an integral domain but $\mathbb{C}[x] /\left\langle x^{2}+1\right\rangle$ is not. (Why?)
(v) $\mathbb{R}[x, y] /\left\langle x^{2}-y^{2}\right\rangle$ is not an integral domain. Geometrically, $V\left(\left\langle x^{2}-y^{2}\right\rangle\right)$ corresponds to two crossing lines in $\mathbb{R}^{2}$. The ring $\mathbb{R}[x, y] /\left\langle x^{2}-y^{2}\right\rangle$ is an integral domain. Geometrically, $V\left(\left\langle x^{2}-y^{2}\right\rangle\right)$ is a cusp in $\mathbb{R}^{2}$, an irreducible curve (see later about the connection between irreducible algebraic varieties and prime ideals).

Theorem 3.8. Let $I \subsetneq R$ be an ideal. Then $I$ is prime if and only if $R / I$ is an integral domain.
Proof. Suppose $I$ is prime. Then since $I \neq R$ we have $R / I \neq\{0\}$. Now suppose $a+I$ is non-zero in $R / I$ and there is some $b+I \in R / I$ such that $(a+I)(b+I)=I$. Then $a b+I=I$ and $a b \in I$. Since $I$ is prime we have either $a \in I$ or $b \in I$, but since $a+I \neq I$ this forces $b \in I$. Hence $b+I=0$ in $R / I$, and $R / I$ is an integral domain.
Suppose now that $R / I$ is an integral domain. Since $R / I \neq\{0\}$ we must have $I \neq R$. Now let $a b \in I$ for some $a, b \in R$, then $a b+I=(a+I)(b+I)=I$. Since $R / I$ is an integral domain, we must have either $a+I=I$ or $b+I=I$, and hence either $a \in I$ or $b \in I$. Therefore $I$ is prime.

Theorem 3.9. Let $R$ be a ring, $I_{1}, \ldots, I_{n} \subseteq R$ be ideals, and $\mathfrak{p} \subseteq R$ be a prime ideal. Then the following are equivalent:
(i) $I_{j} \subseteq \mathfrak{p}$ for some $1 \leqslant j \leqslant n$;
(ii) $I_{1} \cap \cdots \cap I_{n} \subseteq \mathfrak{p}$;
(iii) $I_{1} \ldots I_{n} \subseteq \mathfrak{p}$.

Proof. $(i) \Longrightarrow(i i) \Longrightarrow(i i i)$ are trivial.
(iii) $\Longrightarrow \quad(i)$ : Assume that $I_{1} \ldots I_{n} \subseteq \mathfrak{p}$ but for all $1 \leqslant j \leqslant n$ we can choose $a_{j} \in I_{j} \backslash \mathfrak{p}$. Then $a_{1} \ldots a_{n} \in I_{1} \ldots I_{n} \backslash \mathfrak{p}$ as $\mathfrak{p}$ is prime, a contradiction.

## 4 Maximal ideals

Definition 4.1. An ideal $I$ of a ring $R$ is called a maximal ideal if:
(i) $I \neq R$;
(ii) there is no ideal $J$ of $R$ such that $I \subsetneq J \subsetneq R$.

Example 4.2. (i) $p \mathbb{Z} \subseteq \mathbb{Z}$ is a maximal ideal for $p$ prime (we will see a proof of this soon).
(ii) $\langle X\rangle \subseteq R[X, Y]$ is not maximal, as $\langle X\rangle \subsetneq\langle X, Y\rangle \subsetneq R[X, Y]$.

Theorem 4.3. Maximal ideals are prime.
Proof. Let $\mathfrak{m}$ be a maximal ideal of a ring $R$ and suppose $a b \in \mathfrak{m}$ for some $a, b \in R$. If neither $a$ nor $b$ are in $\mathfrak{m}$ then both $\langle a\rangle+\mathfrak{m}$ and $\langle b\rangle+\mathfrak{m}$ are strictly bigger than $\mathfrak{m}$. As $\mathfrak{m}$ is maximal, we must then have $\langle a\rangle+\mathfrak{m}=\langle b\rangle+\mathfrak{m}=R$. But now

$$
\begin{aligned}
R & =R R \\
& =(\langle a\rangle+\mathfrak{m})(\langle b\rangle+\mathfrak{m}) \\
& =\mathfrak{m}^{2}+\langle a\rangle \mathfrak{m}+\langle b\rangle \mathfrak{m}+\langle a b\rangle \\
& \subseteq \mathfrak{m} \neq R,
\end{aligned}
$$

which is a contradiction.
Proposition 4.4. Let $R$ be a ring. Then:
(i) $R$ is a field iff $\{0\}$ and $R$ are the only ideals of $R$;
(ii) an ideal $I \subseteq R$ is maximal if and only if $R / I$ is a field.

Proof. (i) Assume $R$ is a field and let $I \subseteq R$ be a non-zero ideal. Choose $r \in I \backslash\{0\}$, then $r$ has an inverse $r^{-1} \in R$. Hence $r^{-1} r=1 \in I$, so $I=R$.

Conversely suppose $\{0\}$ and $R$ are the only ideals of $R$, and choose $r \in R \backslash\{0\}$. Then $\langle r\rangle=R$ and so there exists some $s \in R$ such that $s r=1$, i.e. $r$ has an inverse $r^{-1}=s$. Therefore $R$ is a field.
(ii) If $I$ is maximal then by Proposition $2.10, R / I$ has no ideals other than $\{I\}$ and $R / I$. Therefore $R / I$ is a field by (i).
If now $R / I$ is a field then again by Proposition 2.10 and (i), any ideal of $R$ which contains $I$ must either be $I$ or $R$, so $I$ is maximal.

Remark. Let $\varphi: R \rightarrow S$ be a ring homomorphism. Unlike the situation with prime ideals, $\mathfrak{m} \subseteq S$ maximal does not imply that $\varphi^{-1}(\mathfrak{m})$ is maximal. For instance, let $\varphi: \mathbb{Z} \rightarrow \mathbb{Q}$ be the inclusion map. Then $\left\{0_{\mathbb{Q}}\right\} \subseteq \mathbb{Q}$ is maximal as $\mathbb{Q}$ is a field, but $\varphi^{-1}\left(\left\{0_{\mathbb{Q}}\right\}\right)=\left\{0_{\mathbb{Z}}\right\} \subsetneq 2 \mathbb{Z} \subsetneq \mathbb{Z}$, so $\varphi^{-1}\left(\left\{0_{\mathbb{Q}}\right\}\right)$ is not maximal.

However we do have the following result which is analogous to Remark 3.5:
Proposition 4.5. The bijection between ideals of $R / I$ and ideals of $R$ containing I restricts to a bijection between maximal ideals of $R / I$ and maximal ideals of $R$ containing $I$.

Proof. Exercise.

We will soon show that every proper ideal is contained in some maximal ideal. In order to prove this however, we must take a brief diversion into set theory.

A partially ordered set or poset $(\Sigma, \leqslant)$ is a set $\Sigma$ and a binary relation $\leqslant \subseteq \Sigma \times \Sigma$ which is:
(i) reflexive, i.e. $x \leqslant x \forall x \in \Sigma$;
(ii) transitive, i.e. $x \leqslant y$ and $y \leqslant z \Longrightarrow x \leqslant z \forall x, y, z \in \Sigma$;
(iii) antisymmetric, i.e. $x \leqslant y$ and $y \leqslant x \Longrightarrow x=y \forall x, y \in \Sigma$.

A subset $S \subseteq \Sigma$ is totally ordered if for all $s, t \in S$ we have either $s \leqslant t$ or $t \leqslant s$ (or both).
Given a subset $S \subseteq \Sigma$, an element $u \in \Sigma$ is an upper bound for $S$ if $s \leqslant u$ for all $s \in S$.
A maximal element of $\Sigma$ is an element $m \in \Sigma$ such that there is no $s \in \Sigma$ with $m \leqslant s$ and $m \neq s$.
Example. A poset without a maximal element is the set $(\mathbb{Z}, \leqslant)$.
Theorem (Zorn's Lemma). Suppose that $(\Sigma, \leqslant)$ is a non-empty poset and that any totally ordered subset $S \subseteq \Sigma$ has an upper bound in $\Sigma$. Then $\Sigma$ has a maximal element.

This is equivalent to the Axiom of Choice, and we take it as an axiom in ZFC (where we generally do maths).

We can now prove the following:
Proposition 4.6. Let $R$ be a non-zero ring. Then every proper ideal $I$ is contained in a maximal ideal.
Proof. Let $\Sigma$ be the set of ideals $J \subsetneq R$ containing $I$, ordered by inclusion $\subseteq$. Then $(\Sigma, \subseteq)$ is a nonempty poset, since $I \in \Sigma$. If $\left\{J_{\lambda}: \lambda \in \Lambda\right\}$ is a totally ordered subset of $\Sigma$ then clearly $J^{*}=\cup_{\lambda \in \Lambda}$ is a proper ideal of $R$ containing $I$, and moreover $J^{*}$ is an upper bound for $\left\{J_{\lambda}: \lambda \in \Lambda\right\}$. By Zorn's Lemma, $\Sigma$ then has a maximal element. But a maximal element of $\Sigma$ is an ideal $\mathfrak{m} \neq R$ containing $I$ with no proper ideals $J$ containing it, so is a maximal ideal containing $I$.

This proposition shows that we usually have lots of maximal ideals, even if they can be hard to find.

Example 4.7. Let $K$ be a field, $R=K\left[x_{1}, \ldots, x_{n}\right]$ and $a_{1}, \ldots, a_{n} \in K$. Then $\mathfrak{m}=\left\langle x_{1}-a_{1}, \ldots, x_{n}-\right.$ $\left.a_{n}\right\rangle$ is a maximal ideal. If it wasn't, then there would exist a polynomial $f \in R$ such that $f \neq \mathfrak{m}$ and $\langle f\rangle+\mathfrak{m} \subsetneq R$. Applying the division algorithm $n$ times gives

$$
f=f_{1}\left(x_{1}-a_{1}\right)+\cdots+f_{n}\left(x_{n}-a_{n}\right)+b
$$

where $f_{i} \in K\left[x_{i}, x_{i+1}, \ldots, x_{n}\right] \subseteq R$ for each $1 \leqslant i \leqslant n$ and $b \in K$. Since $f \notin \mathfrak{m}$, we must have $b \neq 0$ and so $b$ has an inverse $b^{-1}$. Therefore $1=b^{-1}\left(f-f_{i}\left(x_{1}-a_{1}\right)-\cdots-f_{n}\left(x_{n}-a_{n}\right)\right) \in\langle f\rangle+\mathfrak{m}$ and so $\langle f\rangle+\mathfrak{m}=R$, a contradiction.

Are these the only maximal ideals of $K\left[x_{1}, \ldots, x_{n}\right]$ ? The answer is yes when $K$ is algebraically closed, but we need a bit more theory in order to prove this.
In some cases, there are far fewer maximal ideals.
Definition 4.8. A ring $R$ is called a local ring if it has precisely one maximal ideal $\mathfrak{m}$. We usually denote this ring by the pair $(R, \mathfrak{m})$.

Example 4.9. (1) If $K$ is a field, then $K$ is a local ring, with maximal ideal $\{0\}$.
(2) The formal power series ring $K[[x]]$ is local with maximal ideal $\langle x\rangle$ (Exercise!).

In order to talk about the prime and maximal ideals in a ring, we introduce the following notions, which will play a crucial role in algebraic geometry, since they allow to define the Zariski topology (see later!).

Definition 4.10. Let $R$ be a ring, then

$$
\operatorname{Spec}(R)=\{\mathfrak{p} \subseteq R: \mathfrak{p} \text { is a prime ideal in } R\}
$$

is called the spectrum of $R$. The set of all maximal ideals of $R$ is called the maximal spectrum of $R$ and denoted by maxSpec $(R)$.

Example 4.11. Let $R=K[x]$ the polynomial ring in one variable over a field $K$. Then $R$ is a principal ideal ring, and an ideal $I \subseteq R$ is maximal if and only if $I$ is prime if and only if $I$ is generated by an irreducible polynomial $P(x)$. Thus we have

$$
\operatorname{Spec}(R)=\operatorname{maxSpec}(R)=\{\langle P(x)\rangle \subseteq K[x]: P(x) \text { is irreducible }\}
$$

If $K$ is algebraically closed, then $P(x) \in K[x]$ is irreducible if and only if $\operatorname{deg}(P(x))=1$, that is, $P(x)$ can be written as $P(x)=x-\lambda$, where $\lambda \in K$. Thus we get

$$
\operatorname{Spec}(R)=\{\langle x-\lambda\rangle: \lambda \in K\} .
$$

This means that elements in $\operatorname{Spec}(R)$ are in bijection with elements of $K$, or said differently, with points in $\mathbb{A}_{K}^{1}$, the affine line.
More generally, one can show that elements of maxSpec $\left(K\left[x_{1}, \ldots, x_{n}\right]\right)$ for $K$ algebraically closed are in bijection with points in $\mathbb{A}_{K}^{n}=K^{n}$. (cf. example 4.7)

## 5 Polynomial ring $K\left[x_{1}, \ldots, x_{n}\right]$

We have already defined the polynomial ring in $n$ variables over a field $K$ via: $K\left[x_{1}, \ldots, x_{n}\right]=$ $\left(K\left[x_{1}, \ldots, x_{n-1}\right]\right)\left[x_{n}\right]$. In the following we study some properties of these rings and in particular define monomial orderings, that will be useful when dealing with the question on defining a division algorithm on $K\left[x_{1}, \ldots, x_{n}\right]$.
First note that the elements of $K\left[x_{1}, \ldots, x_{n}\right]$ are finite sums of the form $P\left(x_{1}, \ldots, x_{n}\right)=\sum_{\alpha \in \mathbb{N}^{n}} a_{\alpha} \underline{x}^{\alpha}$. (We sometimes write short $K[\underline{x}]$ for $K\left[x_{1}, \ldots, x_{n}\right]$ and $\underline{x}^{\alpha}$ for $x_{1}^{\alpha_{1}} \cdots x_{n}^{\alpha_{n}}$ ). An element $\underline{x}^{\alpha}$ of $K[\underline{x}]$ is called a monomial. The $a_{\alpha}$ in $P(\underline{x})=\sum_{\alpha \in \mathbb{N}^{n}} a_{\alpha} \underline{x}^{\alpha}$ are called coefficients of $P$.

One can distinguish between polynomials $P(\underline{x})$ as elements of the polynomial ring $K[\underline{x}]$ or as polynomial maps, that is, any $P$ gives a map

$$
P: K^{n} \rightarrow K,\left(a_{1}, \ldots, a_{n}\right) \mapsto P\left(a_{1}, \ldots, a_{n}\right) .
$$

Given polynomials $P_{1}(\underline{x}), \ldots, P_{m}(\underline{x}) \in K[\underline{x}]$ one defines

$$
V\left(P_{1}, \ldots, P_{m}\right)=\left\{\left(a_{1}, \ldots, a_{n}\right) \in K^{n}: P_{i}\left(a_{1}, \ldots, a_{n}\right)=0 \text { for all } i=1, \ldots, m\right\}
$$

the vanishing set (or zero-set) of $P_{1}, \ldots, P_{m}$ in $K^{n}$. One writes $\mathbb{A}_{K}^{n}:=K^{n}=\left\{\left(a_{1}, \ldots, a_{n}\right) \in K^{n}\right\}$ for the affine $n$-space over $K$. If $X \subseteq \mathbb{A}_{K}^{n}$ is of the form $X=V\left(P_{1}, \ldots, P_{m}\right)$, then $X$ is called an algebraic set and the $P_{1}, \ldots, P_{m}$ define $X$. If $X \subseteq \mathbb{A}_{K}^{n}$ is an algebraic set, then

$$
I(X)=\left\{P(\underline{x}) \in K\left[x_{1}, \ldots, x_{n}\right]: P\left(a_{1}, \ldots, a_{n}\right)=0 \text { for all }\left(a_{1}, \ldots, a_{n}\right) \in X\right\}
$$

is an ideal in $K\left[x_{1}, \ldots, x_{n}\right]$, the defining ideal of $X$. Later we will study the relation between ideals in $K\left[x_{1}, \ldots, x_{n}\right]$ and algebraic sets in $\mathbb{A}_{K}^{n}$.

Example 5.1. (1) $X=V\left(x^{3}-y^{2}\right) \subseteq \mathbb{A}_{\mathbb{R}}^{2}$ defines a cusp. This is an irreducible curve in the real plane.
(2) $X=V\left(x^{2}+y^{2}\right) \subseteq \mathbb{A}_{\mathbb{R}}^{2}$ is the point $\{(0,0)\}$. However, $V\left(x^{2}+y^{2}\right) \subseteq \mathbb{A}_{\mathbb{C}}^{2}$ consists of the two lines $\{x+i y=0\}$ and $\{x-i y=0\}$.
(3) Consider $J=\left\langle x^{3}, x y, y^{2}, z\right\rangle \subseteq K[x, y, z]$. Then one can see that $V(J)=\{(0,0,0)\}$, but $I(V(J))=$ $\langle x, y, z\rangle \supsetneq J$.

Consider the polynomial ring $K\left[x_{1}, \ldots, x_{n}\right]$. We define the (total) degree of a monomial $x_{1}^{\alpha_{1}} \cdots x_{n}^{\alpha_{n}}$ as $|\alpha|=\alpha_{1}+\cdots \alpha_{n}$. Consequently, the degree of a polynomial $P\left(x_{1}, \ldots, x_{n}\right)=\sum_{\alpha \in \mathbb{N}^{n}} a_{\alpha} \underline{x}^{\alpha}$ is $\operatorname{deg}(P)=\max \left\{|\alpha|: a_{\alpha} \neq 0\right\}$. The order of $P$ is $\operatorname{ord}(P)=\min \left\{|\alpha|: a_{\alpha} \neq 0\right\}$.
We can write $P(\underline{x})=\sum_{d} P^{(d)}$, where $P^{(d)}$ is the sum of all monomials in $P(\underline{x})$ with $\operatorname{deg}\left(\underline{x}^{\alpha}\right)=d$. If $P \neq 0$, then we say that $P(\underline{x})$ is homogeneous of degree $d$ if $P(\underline{x})=P^{(d)}$.
Example 5.2. (1) $P: \mathbb{R}^{3} \rightarrow \mathbb{R}:(x, y, z) \mapsto x^{2} y+x y z+x^{2} y^{2}-\sqrt{2} z^{3}$ corresponds to the polynomial $P \in \mathbb{R}[x, y, z]$ with $\operatorname{deg}(P)=4$, ord $(P)=3$ and $P=P^{(3)}+P^{(4)}$, with $P^{(3)}=x^{2} y+x y z-\sqrt{2} z^{3}$ and $P^{(4)}=x^{2} y^{2}$.
(2) $P(x, y, z)=x^{3} y z-x y^{4}$ is homogeneous of degree 4 .

Remark 5.3. We can decompose $K[\underline{x}]$ into graded components, where each graded component is a finite-dimensional $K$-vector space:

$$
K\left[x_{1}, \ldots, x_{n}\right]=\bigoplus_{d=0}^{\infty} K\left[x_{1}, \ldots, x_{n}\right]_{d}
$$

where $K\left[x_{1}, \ldots, x_{n}\right]_{d}:=\{$ homogeneous polynomials of degree $d\}$. Each $K\left[x_{1}, \ldots, x_{n}\right]_{d}$ is a finite dimensional $K$-vector space with basis all monomials of degree $d$ (What is its dimension?). For example, for $n=2$ we have $K[x, y]_{0}=K, K[x, y]_{1}=K x \oplus K y \cong K^{2}, K[x, y]_{2}=K x^{2} \oplus K x y \oplus K y^{2} \cong$ $K^{3}, \ldots$.

Next we consider ring homomorphisms from $K[\underline{x}]$. In particular important are evaluation homomorphisms: Let $a \in K^{n}$, and define

$$
\varepsilon_{a}: K\left[x_{1}, \ldots, x_{n}\right] \rightarrow K: P \mapsto P\left(a_{1}, \ldots, a_{n}\right) .
$$

$\varepsilon_{a}$ is a ring homomorphism and in particular, if $a=(0, \ldots, 0)$, then $\varepsilon_{0}(P)=P(0)$ yields the constant term of $P$.
More generally, define substitution homomorphisms: let $f \in K\left[x_{1}, \ldots, x_{n}\right]$ and $g_{1}, \ldots g_{n} \in K\left[y_{1}, \ldots, y_{m}\right]$. Then $f\left(g_{1}, \ldots, g_{n}\right)$ is an element of $K\left[y_{1}, \ldots, y_{m}\right]$. This can be described by the homomorphism

$$
g^{*}: K\left[x_{1}, \ldots x_{n}\right] \rightarrow K\left[y_{1}, \ldots, y_{m}\right]: f \mapsto g^{*}(f)=f\left(g_{1}, \ldots, g_{n}\right)
$$

The evaluation homomorphism $\varepsilon_{a}$ is a special case, that is, set $g_{i}=a_{i}$ in $K$, then $g^{*}=\varepsilon_{a}$.

## Monomial orderings of $K[\underline{x}]$

If $n=1$, then the degree gives a total order on the set of monomials in $K[x]: x^{\alpha}<x^{\beta}$ if and only if $\alpha<\beta$. However, if $n \geqslant 2$, the degree only yields a partial order on the set of monomials, e.g., for $n=2$, both monomials $x_{1} x_{2}$ and $x_{1}^{2}$ have the same degree. In order to get a total order on monomials, we introduce the following:
Definition 5.4. A monomial ordering $>_{\varepsilon}$ on $K\left[x_{1}, \ldots, x_{n}\right]$ (or, equivalently, on $\mathbb{N}^{n}$ ) is a total order on the set of monomials $\underline{x}^{\alpha}, \alpha \in \mathbb{N}^{n}$ of $K\left[x_{1}, \ldots, x_{n}\right]$ (that is, either $\underline{x}^{\alpha}>_{\varepsilon} \underline{x}^{\beta}, \underline{x}^{\alpha}=\underline{x}^{\beta}$, or $\underline{x}^{\alpha}<_{\varepsilon} \underline{x}^{\beta}$ ) such that
(i) If $\alpha>_{\varepsilon} \beta$ and $\gamma \in \mathbb{N}^{n}$, then $\alpha+\gamma>_{\varepsilon} \beta+\gamma$.
(ii) $>_{\varepsilon}$ is a well-ordering on $\mathbb{N}^{n}$ (this means that every non-empty subseteq of $\mathbb{N}^{n}$ has a smallest element with respect to $>_{\varepsilon}$ ).
We write $\alpha \geqslant_{\varepsilon} \beta$ if $\alpha>_{\varepsilon} \beta$ or $\alpha=\beta$.
Example 5.5. (1) The lexicographic order $>_{\text {lex }}$ is a monomial order (see homework for a proof!) defined (on $\mathbb{N}^{n}$ ) as follows: $\alpha>_{\text {lex }} \beta: \Leftrightarrow$ there exists a $j \leqslant n$ such that $\alpha_{i}=\beta_{i}$ for all $i<j$ and $\alpha_{j}>\beta_{j}$.
(2) The degree lexicographic order $>_{\text {deglex }}$ is defined as:

$$
\alpha>_{\text {deglex }} \beta: \Leftrightarrow\left\{\begin{array}{l}
|\alpha|>|\beta| ; \text { or } \\
|\alpha|=|\beta| \text { and } \alpha>_{\text {lex }} \beta
\end{array}\right.
$$

(3) The reverse lexicographic order $>_{\text {revlex }}: \alpha>_{\text {revlex }} \beta: \Leftrightarrow$ there exists a $j \geqslant 1$ such that $\alpha_{i}=\beta_{i}$ for all $i>j$ and $\alpha_{j}>\beta_{j}$.
Example 5.6. More generally, one can define a linear order $>_{\lambda}$ : Let $\lambda \in \mathbb{R}_{+}^{n}$ be a vector with $\mathbb{Q}$ linearly independent components. Then $\lambda$ induces a linear map $\lambda: \mathbb{N}^{n} \rightarrow \mathbb{R}_{\geqslant 0}, \alpha \mapsto\langle\alpha, \lambda\rangle=$ $\sum_{i=1}^{n} \alpha_{i} \lambda_{i}$. Then $\alpha>_{\lambda} \beta: \Leftrightarrow\langle\alpha, \lambda\rangle>\langle\beta, \lambda\rangle$.

Example 5.7. For $n=2$, consider $>_{l e x}$ : Then $x_{1}^{2} x_{2}^{3}>_{l e x} x_{1}^{2} x_{2}$, because $(2,3)$ is greater than $(2,1)$ in the lexicographic order. Also $x_{1}^{2}>_{\text {lex }} x_{2}^{3}$.
For $>_{\text {deglex }}$ we similarly compute $x_{1}^{2} x_{2}^{3}>_{\text {lex }} x_{1}^{2} x_{2}$ but $x_{1}^{2}<_{\text {deglex }} x_{2}^{3}$
Definition 5.8. Let $f(\underline{x})=\sum_{\alpha \in \mathbb{N}^{n}} a_{\alpha} \underline{x}^{\alpha} \in K\left[x_{1}, \ldots, x_{n}\right]$ and let $>_{\varepsilon}$ be a monomial order. Then $\operatorname{deg}_{\varepsilon}(f)=\max _{>_{\varepsilon}}\left(\alpha \in \mathbb{N}^{n}: a_{\alpha} \neq 0\right)$ is called the $>_{\varepsilon}$-degree of $f$. The leading coefficient $l_{c_{\varepsilon}}(f)$ is $a_{\operatorname{deg}_{\varepsilon}(f)} \in K$. The leading monomial of $f$ is $\operatorname{lm}(f)=x^{\operatorname{deg}_{\varepsilon}(f)}$. The leading term of $f$ is $l t_{\varepsilon}(f)=$ $l c_{\varepsilon}(f) \cdot \operatorname{lm}_{\varepsilon}(f)$.

Remark 5.9. This is already enough to define an Euclidean division on $K\left[x_{1}, \ldots, x_{n}\right]$ (see later in Section 18 on Gröbner bases).

## 6 Localisation

We can construct $\mathbb{Q}$ from $\mathbb{Z}$ by inverting all non-zero elements. Formally this is done by viewing $\mathbb{Q}$ as a set of equivalence classes in $\mathbb{Z} \times(\mathbb{Z} \backslash\{0\})$ via the relation

$$
(r, a) \sim(s, b) \Longleftrightarrow a s=b r
$$

We then write $\frac{r}{a}$ for the equivalence class of $(r, a)$. Addition and multiplication of equivalence classes is defined by

$$
\begin{equation*}
\frac{r}{a}+\frac{s}{b}=\frac{a s+b r}{a b} \text { and } \frac{r}{a} \frac{s}{b}=\frac{r s}{a b} \tag{*}
\end{equation*}
$$

We also have $0_{\mathbb{Q}}=\frac{0}{1}$ and $1_{\mathbb{Q}}=\frac{1}{1}$. It is easy to check that provided $r \neq 0, \frac{a}{r}$ is a multiplicative inverse for $\frac{r}{a}$.
We wish to repeat the above for a general ring $R$. Notice from $(*)$ that if we invert $a$ and $b$ then we have also inverted $a b$. This motivates the following.
Definition 6.1. Let $R$ be a ring and $A \subseteq R$ be a subset. We say $A$ is multiplicatively closed if:
(i) $1_{R} \in A$;
(ii) $a, b \in A \Longrightarrow a b \in A$.

Example 6.2. (1) For any ring, $R$ itself is multiplicatively closed. If $R=K$, then $K^{*}=K \backslash\{0\}$ is multiplicatively closed.
(2) If $f \in R=K\left[x_{1}, \ldots, x_{n}\right]$ is a nonzero element, then $A=\left\{1, f, f^{2}, f^{3}, \ldots\right\}$ is a multiplicatively closed set.

Definition 6.3. Let $R$ be a ring and $A \subseteq R$ be multiplicatively closed. The localisation of $R$ at $A$, denoted $A^{-1} R$ or $R\left[A^{-1}\right]$ or $R_{A}$, is the set of equivalence classes of $R \times A$ under the equivalence relation

$$
(r, a) \sim(s, b) \Longleftrightarrow \text { there exists a } c \in A \text { such that } c(a s-b r)=0
$$

We will again usually write the equivalence class of $(r, a)$ as $\frac{r}{a}$, with addition and multiplication defined as in $(*)$.

Lemma 6.4. Let $R$ be a ring and $A \subseteq R$ a multiplicatively closed subset. Then the localisation $A^{-1} R$ of $R$ at $A$ is also a ring via the sum and product $(*)$, and $0_{A^{-1} R}=\frac{0_{R}}{1_{R}}$ and $1_{A^{-1} R}=\frac{1_{R}}{1_{R}}$. Moreover there is a ring homomorphism

$$
\begin{aligned}
i: R & \rightarrow A^{-1} R \\
r & \mapsto \frac{r}{1}
\end{aligned}
$$

with kernel Ker $i=\{r \in R: r a=0$ for some $a \in A\}$.
In some cases, such as the construction of $Q$ above, we wish to invert as many things as possible.
Definition 6.5. Let $R$ be an integral domain. The quotient field or field of fractions of $R$, denoted Quot $(R)$, is the localisation

$$
\operatorname{Quot}(R)=(R \backslash\{0\})^{-1} R
$$

Example 6.6. In each of the following, $A$ is a multiplicatively closed subset of a ring $R$.
(i) $R_{A}$ is the zero ring if and only if $0 \in A$.
(ii) Let $a \in A$. We write $R_{a}$ for the localisation of $R$ at the set $\left\{a^{n}: n \geqslant 0\right\}$.
(iii) Let $\mathfrak{p}$ be a prime ideal of $R$. Then $A=R \backslash \mathfrak{p}$ is multiplicatively closed and we write $R_{\mathfrak{p}}$ for $A^{-1} R$. (Careful here! The "correct" way to write this would be $R_{R \backslash \mathfrak{p}}$ ).
(iv) Let $p \in \mathbb{Z}$ be prime. Then

$$
\begin{aligned}
\mathbb{Z}_{p} & =\left\{\frac{a}{b} \in \mathbb{Q}: b \text { is a power of } p\right\}, \\
\mathbb{Z}_{\langle p\rangle} & =\left\{\frac{a}{b} \in \mathbb{Q}: p \nmid b\right\}, \\
\operatorname{Quot}(\mathbb{Z}) & =\mathbb{Q} .
\end{aligned}
$$

Since $A^{-1} R$ is a ring, we can talk about its ideals and how they relate to the ideals of $R$.
Definition 6.7. Given an ideal $I$ of $R$, we define the localisation of the ideal $I$ to be the set

$$
A^{-1} I=\left\{\frac{x}{a}: x \in I, a \in A\right\} .
$$

Proposition 6.8. Let $R$ be a ring, $A \subseteq R$ a multiplicatively closed subset, and $I \subseteq R$ an ideal.
(i) $A^{-1} I$ is an ideal of $A^{-1} R$. Moreover, if $I$ is generated by a set $X$, then $A^{-1} I$ is generated by $\left\{\frac{x}{1}: x \in X\right\}$.
(ii) We have $\frac{x}{a} \in A^{-1}$ I if and only if there is some $b \in A$ with $x b \in I$.
(iii) $A^{-1} I=A^{-1} R$ if and only if $I \cap A \neq \varnothing$.
(iv) The map $I \mapsto A^{-1}$ I commutes with forming finite sums, products and intersections, and quotients.

Proof. See Homework Sheet.
This leads to a correspondence theorem for between ideals of $R$ and ideals of $A^{-1} R$.
Theorem 6.9. There is a bijection
$\left\{\right.$ ideals $\left.J \subseteq A^{-1} R\right\} \leftrightarrow\{$ ideals $I \subseteq R$ such that no element of $A$ is a zero divisor in $R / I\}$,
sending $J \mapsto i^{-1}(J)$ and $I \mapsto A^{-1} I$, where $i^{-1}(J)$ is the preimage of $J$ under the homomorphism from Lemma 6.4.
Moreover, this restricts to a bijection

$$
\left\{\text { prime ideals } Q \subseteq A^{-1} R\right\} \leftrightarrow\{\text { prime ideals } P \subseteq R \text { with } P \cap A=\varnothing\}
$$

Proof. Suppose $J \subseteq A^{-1} R$ is an ideal. Then $i^{-1}(J)$ is an ideal, being the preimage of an ideal under a ring homomorphism. By definition we have

$$
i^{-1}(J)=\left\{x \in R: \frac{x}{1} \in J\right\}
$$

and therefore $A^{-1}\left(i^{-1}(J)\right) \subseteq J$ (see Definition 6.7). Conversely if $\frac{x}{a} \in J$ then $\frac{x}{1}=\frac{a}{1} \frac{x}{a} \in J$, so $x \in i^{-1}(J)$. Thus $\frac{x}{a} \in A^{-1}\left(i^{-1}(J)\right)$ hence $J \subseteq A^{-1}\left(i^{-1}(J)\right)$, and therefore $J=A^{-1}\left(i^{-1}(J)\right)$.
We have shown that the maps are inverses to one another, so we must determine the image of $J \mapsto i^{-1}(J)$. We claim that $I$ is in the image if and only if $I=i^{-1}\left(A^{-1} I\right)$. Indeed, such an ideal is certainly in the image of $i^{-1}$, whereas if $I=i^{-1}(J)$ then $A^{-1} I=A^{-1}\left(i^{-1}(J)\right)=J$, and so $i^{-1}\left(A^{-1} I\right)=i^{-1}(J)=I$.
Now we always have $I \subseteq i^{-1}\left(A^{-1} I\right)$, so $I \neq i^{-1}\left(A^{-1} I\right)$ if and only if there is some $x \notin I$ such that $\frac{x}{1} \in A^{-1} I$. By Proposition 6.8(ii), this is equivalent to there being some $x \notin I$ and $b \in A$ with $x b \in I$. That is, there exists $b \in A$ and $x+I \neq I=0_{R / I}$ in $R / I$ with $(b+I)(x+I)=I=0_{R / I}$, i.e. some element of $A$ is a zero divisor in $R / I$.
For the second part, observe first that if $P \subseteq R$ is prime then $R / P$ is an integral domain (Theorem 3.8), so $A$ contains a zero divisor in $R / P$ if and only if $A \cap P \neq \varnothing$. It is therefore enough to show that prime ideals always map to prime ideals. Recall from Proposition 3.3 that if $Q \subseteq A^{-1} R$ is prime, then $i^{-1}(Q) \subseteq R$ is prime. On the other hand if $P \subseteq R$ is prime and $P \cap A=\varnothing$, then $R / P$ is an integral domain and $\bar{A} \subseteq R / P$ does not contain $0_{R / P}$, so by Proposition 6.8(iv) we have

$$
A^{-1} R / A^{-1} P \cong \bar{A}^{-1}(R / P) \subseteq \operatorname{Quot}(R / P)
$$

Since $\operatorname{Quot}(R / P)$ is a field, it contains no non-zero zero divisors. Therefore as a subring neither does $A^{-1} R / A^{-1} P$, i.e. it is an integral domain, and so $A^{-1} P \subseteq A^{-1} R$ is a prime ideal.
The following corollary then gives an insight into the name "localisation".
Corollary 6.10. Let $\mathfrak{p} \subseteq R$ be a prime ideal. Then the prime ideals of $R_{\mathfrak{p}}$ are in bijection with the prime ideals of $R$ contained in $\mathfrak{p}$. In particular $R_{\mathfrak{p}}$ has a unique maximal ideal $P_{\mathfrak{p}}$, and hence $\left(R_{\mathfrak{p}}, \mathfrak{p}_{\mathfrak{p}}\right)$ is a local ring.

Proof. By Theorem 6.9, the prime ideals of $R_{\mathfrak{p}}$ are in bijection with the prime ideals $\mathfrak{p}^{\prime}$ of $R$ that do not intersect $R \backslash \mathfrak{p}$. But this is precisely the condition that $\mathfrak{p}^{\prime} \subseteq \mathfrak{p}$.
The maximality and uniqueness of $\mathfrak{p}_{\mathfrak{p}}$ follows from the fact that the bijection is inclusion preserving. In particular if $Q_{1} \subseteq Q_{2}$ are ideals of $R_{\mathfrak{p}}$ then $i^{-1}\left(Q_{1}\right) \subseteq i^{-1}\left(Q_{2}\right)$, and if $P_{1} \subseteq P_{2}$ are ideals of $R$ then $\left(P_{1}\right)_{\mathfrak{p}} \subseteq\left(P_{2}\right)_{\mathfrak{p}}$. The largest prime ideal of $R$ contained in $\mathfrak{p}$ is $\mathfrak{p}$ itself, and this is the unique ideal with this property, therefore $\mathfrak{p}_{\mathfrak{p}}$ is the unique maximal ideal of $R_{\mathfrak{p}}$.

Theorem 6.11 (Universal property of the localisation). Let $R$ be a ring and $A \subseteq R$ be a multiplicatively closed set. Let $\varphi: R \rightarrow A^{-1} R, r \mapsto \frac{r}{1}$ the ring homomorphism from above (note here: $\varphi(A) \subseteq A^{-1} R$ is invertible in the localisation $A^{-1} R$ ). Let $f: R \rightarrow B$ be a ring homomorphism such that $f(a)$ is a unit in $B$ for all $a \in A$. Then there exists a unique ring homomorphism $h: A^{-1} R \rightarrow B$ such that $f=h \circ \varphi$ :


Proof. (1) We show uniqueness first: If $h$ satisfies the conditions of the theorem, then $h\left(\frac{r}{1}\right)=$ $h \circ \varphi(r)=f(r)$ for all $r \in R$. For any $a \in A$ we have $h\left(\frac{1}{a}\right)=h\left(\left(\frac{a}{1}\right)^{-1}\right)=h\left(\frac{a}{1}\right)^{-1}$ (check this!), and this is equal to $f(a)^{-1}$. Therefore $h\left(\frac{r}{a}\right)=h\left(\frac{r}{1} \cdot \frac{1}{a}\right)=h\left(\frac{r}{1}\right) h\left(\frac{1}{a}\right)=f(r) f(a)^{-1}$. This means that $h$ is uniquely determined by $f$.
(2) For the existence we first define $h\left(\frac{r}{a}\right):=f(r) f(a)^{-1}$. Then we have to show that $h$ is a welldefined ring homomorphism: for the well-definedness, assume that $\frac{r}{a}=\frac{r^{\prime}}{a^{\prime}}$. Then there exists a
$c \in A$ such that $c r a^{\prime}=c r^{\prime} a$. Thus $f(0)=f\left(c r a^{\prime}-c r^{\prime} a\right)=f(c)\left(f(r) f\left(a^{\prime}\right)-f\left(r^{\prime}\right) f(a)\right)$ since $f$ is a ring homomorphism. Since $c \in A$, by assumption $f(c)$ is a unit in $B$, thus $f(r) f\left(a^{\prime}\right)=f\left(r^{\prime}\right) f(a)$ and this implies that

$$
f(r) f(a)^{-1}=f\left(a^{\prime}\right)^{-1} f\left(r^{\prime}\right)
$$

and the left hand side of this equation is equal to $h\left(\frac{r}{a}\right)$, whereas the right hand side to $h\left(\frac{r^{\prime}}{a^{\prime}}\right)$. Showing that $h$ is a ring homomorphism is an exercise.

Remark 6.12. This theorem shows that the localisation $A^{-1} R$ is uniquely determined by the following conditions: if $f: R \rightarrow B$ is any ring homomorphism such that
(i) $a \in A$ implies that $f(a)$ is a unit in $B$,
(ii) $f(r)=0$ implies that $r a=0$ for some $a \in A$,
(iii) every element of $B$ is of the form $f(r) f(a)^{-1}$,
then there exists a unique ring isomorphism $h: A^{-1} \rightarrow B$ such that $f=h \circ \varphi$.

## 7 The radical, nilradical and Jacobson radical

Recall that an element $x$ in a ring $R$ is called zero-divisor if there exists a $y \neq 0$ in $R$ such that $x \cdot y=0$.

Example 7.1. (1) $0 \in R$ is always a zero-divisor.
(2) $\mathbb{Z}, K\left[x_{1}, \ldots, x_{n}\right]$, and more generally, any integral domain $R$ does not have nonzero zerodivisors.
(3) In $K[x, y] /\langle x y\rangle$ every element contained in the maximal ideal $\langle\bar{x}, \bar{y}\rangle$ is a zero-divisor.

Definition 7.2. Let $R$ be a ring. An element $r \in R$ is nilpotent if there exists an integer $n \geqslant 1$ such that $r^{n}=0$.

Example 7.3. (1) In an integral domain $R$ are no nonzero nilpotent elements.
(2) In the ring $K[x, y] /\langle x y\rangle$ there are no nonzero nilpotent elements.
(3) The ring $K[x] /\langle x\rangle \cong K$, so does not contain any nonzero nilpotent elements. But in $K[x] /\left\langle x^{k}\right\rangle$ for $k \geqslant 2$, ever $x^{i}, 1 \leqslant i \leqslant k$ is nilpotent.
(4) A noncommutative example: In the ring $M_{2}(\mathbb{R})$ of $2 \times 2$ real matrices,

$$
\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right)^{2}=\left(\begin{array}{ll}
0 & 0 \\
0 & 0
\end{array}\right)
$$

Definition 7.4. The nilradical of a ring $R$, denoted $\operatorname{nil}(R)$, is the set of all nilpotent elements of $R$.
Theorem 7.5. Let $R$ be a ring. Then $\operatorname{nil}(R)$ is an ideal of $R$, and moreover is the intersection of all prime ideals of $R$.

Proof. If $r, s \in \operatorname{nil}(R)$ then there exist $n, m \in \mathbb{N}$ such that $r^{n}=s^{m}=0$. By the binomial theorem we have

$$
(r+s)^{n+m}=\sum_{i=0}^{n+m}\binom{n+m}{i} r^{i} s^{n+m-i}
$$

and for all $0 \leqslant i \leqslant n+m$ we have either $i \geqslant n$ or $n+m-i \geqslant m$, so either $r^{i}=0$ or $s^{n+m-i}=0$. Hence $(r+s)^{n+m}=0$ and $r+s \in \operatorname{nil}(R)$. Now for $t \in R,(t r)^{n}=t^{n} r^{n}=0$. Finally $0 \in \operatorname{nil}(R)$ so $\operatorname{nil}(R) \neq \varnothing$, and $\operatorname{nil}(R)$ is an ideal of $R$.
We now show that nil $(R) \subseteq P$ for all prime ideals $P$, therefore giving containment one way. Indeed, let $P$ be a prime ideal. Then for any $r \in \operatorname{nil}(R)$ there exists some $n \in \mathbb{N}$ such that $r^{n}=0 \in P$, but since $P$ is prime we must then have $r \in P$.

Finally, we show that the intersection of all prime ideals is contained in the nilradical. In fact, we will prove the contrapositive. Suppose $r$ is not nilpotent. Then $0 \notin\left\{r^{i}: i \geqslant 1\right\}$ and the set

$$
S=\left\{I \subseteq R: I \text { is an ideal and } r^{i} \notin I \text { for all } i \geqslant 1\right\}
$$

is non-empty as $\{0\} \in S$. We turn $S$ into a poset by inclusion, and then any totally ordered subset of $S$ has an upper bound, namely the union of all its elements (cf. proof of Proposition 4.6). By Zorn's Lemma, there is a maximal element $J \in S$. That $J$ is an ideal is immediate, so we now prove that it is prime. Suppose $a b \in J$ but $a \notin J$ and $b \notin J$. Then $\langle a\rangle+J$ and $\langle b\rangle+J$ are strictly greater than $J$, so $r^{m} \in\langle a\rangle+J$ and $r^{n} \in\langle b\rangle+J$ for some $m, n \in \mathbb{N}$. Thus $r^{n+m} \in(\langle a\rangle+J)(\langle b\rangle+J) \subseteq J$, contradicting the choice of $J$. Therefore $J$ is a prime ideal and moreover $r \notin J$ (set $i=1$ in the above), so $r \notin \bigcap_{P \text { prime }} P$.

Recall the notion of radical ideal: Let $I$ be an ideal of a ring $R$. The radical of $I$, denoted $\sqrt{I}$, is the set $\left\{r \in R: r^{n} \in I\right.$ for some $\left.n \geqslant 1\right\}$. We have already shown (in the exercises) that $\sqrt{I}$ is an ideal in $R$.

Theorem 7.6. Let I be an ideal of a ring $R$. Then $\sqrt{I}$ is an ideal of $R$, and moreover is the intersection of all prime ideals in $R$ which contain $I$.

Proof. Consider the quotient homomorphism $\varphi: R \rightarrow R / I$. Then $r \in \sqrt{I}$ if and only if $\varphi(r) \in$ $\operatorname{nil}(R / I)$, thus $\operatorname{rad}(I)=\varphi^{-1}(\operatorname{nil}(R / I))$ and hence is an ideal.
For the second statement we see that

$$
\begin{aligned}
\sqrt{I} & =\varphi^{-1}(\operatorname{nil}(R / I)) \\
& =\varphi^{-1}\left(\bigcap_{\bar{P} \subseteq R / I \text { prime }} \bar{P}\right) \\
& =\bigcap_{\bar{P} \subseteq R / I \text { prime }} \varphi^{-1}(\bar{P}) \\
& =\bigcap_{\substack{P \subseteq R \text { prime } \\
I \subseteq P}} P,
\end{aligned}
$$

where we have again used Proposition 2.10 in the last step.
Example 7.7. (i) Working in $\mathbb{Z}$, we have $\sqrt{4 \mathbb{Z}}=2 \mathbb{Z}$ and $\sqrt{3 \mathbb{Z}}=3 \mathbb{Z}$.
(ii) Again in $\mathbb{Z}$,

$$
\sqrt{12 \mathbb{Z}}=\bigcap_{\substack{\text { Pprime } \\ 12 \mathbb{Z} \subseteq P}} P
$$

The prime ideals in $\mathbb{Z}$ are $p \mathbb{Z}$, and those containing $12 \mathbb{Z}$ are $2 \mathbb{Z}$ and $3 \mathbb{Z}$. Hence $\sqrt{12 \mathbb{Z}}=$ $2 \mathbb{Z} \cap 3 \mathbb{Z}=6 \mathbb{Z}$.
(iii) Let $I=\left\langle x+y, y^{2}\right\rangle \subseteq \mathbb{R}[x, y]$. Then $y \in \sqrt{I}$, and $x^{2}=y^{2}+(x-y)(x+y) \in I$ so also $x \in \sqrt{I}$. Then $\sqrt{I}=\langle x, y\rangle$.

Definition 7.8. Let $R$ be a ring. The Jacobson radical, denoted $J(R)$, is defined to be the set

$$
J(R)=\bigcap_{\mathfrak{m} \subseteq R \text { maximal }} \mathfrak{m} .
$$

Remark. Note that in a local ring $(R, \mathfrak{m})$ (see Definition 4.8), the Jacobson radical is equal to the maximal ideal, i.e. $J(R)=\mathfrak{m}$.

Lemma 7.9. Let $R$ be a ring and $x \in R$. Then $x \in J(R)$ if and only if $1+r x$ is invertible for all $r \in R$.
Proof. See Exercise Sheet 1.
Example 7.10. Let $R=K[[x]]$. Then $R$ is local with maximal ideal $\mathfrak{m}=\langle x\rangle$. Then by definition we have $J(R)=\mathfrak{m}$ but $\operatorname{nil}(R)=\langle 0\rangle$, as $R$ is a domain.

## 8 Modules

Definition 8.1. Let $R$ be a ring. An abelian group $M=(M,+)$ (with identity 0 ) is an $R$-module (or just a module if it is clear from context) if there exists a multiplication map $\cdot: R \times M \rightarrow M$, $(r, m) \mapsto r m$ such that for all $r, s \in R$ and $m, n \in M$ :
(i) $r(s m)=(r s) m$;
(ii) $r(m+n)=r m+r n$;
(iii) $(r+s) m=r m+s m$;
(iv) $1_{R} m=m$.

Example 8.2. (1) If $R$ is a field then an $R$-module is simply a vector space. The axioms for a module are the same as a vector space except $R$ is not necessarily a field.
(2) Ideals in a ring $R$ are also $R$-modules. In general, an ideal is not isomorphic to $R$ as an $R$ module. Take for example $I=\left\langle x^{3}-y z, y^{2}-x z, z^{2}-x^{2} y\right\rangle \subseteq K[x, y, z]$. Then the three generators are not linearly independent over $K[x, y, z]$. One has the relations $y\left(x^{3}-y z\right)+z\left(y^{2}-\right.$ $x z)+x\left(z^{2}-x^{2} y\right)=z\left(x^{3}-y z\right)+x^{2}\left(y^{2}-x z\right)+y\left(z^{2}-x^{2} y\right)=0$. But the three given polynomials are a minimal generating set for $I$. We see that a module does not need to have a basis (different as for vector spaces).
(3) For a ring $R$, the set $R^{n}$ of $n$-tuples of elements of $R$ is an $R$-module.
(4) $R[x]$ is an $R$-module: it is generated by $R \oplus R x \oplus R x^{2} \oplus \cdots$.
(5) $R$ is a module over itself.
(6) Any abelian group is a $\mathbb{Z}$-module (and vice versa!).
(7) If $S \subseteq R$ is a subring then $R$ is an $S$-module.

Modules therefore generalise the idea of vector spaces to rings.
Definition 8.3. A map $\varphi: M \rightarrow N$ between $R$-modules $M$ and $N$ is an $R$-module homomorphism (or $R$-homomorphism) if $\varphi$ is an $R$-linear map, i.e. $\varphi(r m+s n)=r \varphi(m)+s \varphi(n)$ for all $r, s \in R$ and $m, n \in M$. An $R$-module isomorphism (monomorphism, epimorphism) is a (injective, surjective) bijective $R$-homomorphism. The set of all $R$-homomorphisms from $M$ to $N$ is denoted $\operatorname{Hom}_{R}(M, N)$.

Proposition 8.4. The set $\operatorname{Hom}_{R}(M, N)$ is an $R$-module, via the action $(r \varphi)(m)=r \varphi(m)$ for all $r \in R$, $\varphi \in \operatorname{Hom}_{R}(M, N)$ and $m \in M$.

Proof. Exercise.
Example 8.5. If $\varphi: R \rightarrow S$ is a ring homomorphism, then it is also a morphism of $R$-modules. For this define the $R$-module structure on $S$ via $r \cdot s:=\varphi(r) s$. Then it is easy to see that $\varphi$ is $R$-linear.

If $R$ is a field, then $R$-module homomorphisms are simple linear maps between vector spaces.
Definition 8.6. A submodule $U$ of an $R$-module $M$ is a subgroup $(U,+)$ of $(M,+)$, closed under the restricted action of the multiplication, i.e. $r u \in U$ for all $r \in R$ and $u \in U$.

Note that the inclusion map $U \hookrightarrow M$ is an $R$-module homomorphism.
Example 8.7. (i) Let $I \subseteq R$ be an ideal and $M$ an $R$-module. Then

$$
I M=\left\{\sum_{i=1}^{n} a_{i} m_{i}: n \geqslant 1, a_{i} \in I, m_{i} \in M\right\}
$$

is a submodule of $M$.
(ii) If $U, V \subseteq M$ are submodules, then $U \cap V$ is a submodule of $U, V$ and $M$.

The factor group $M / U$ is also an $R$-module, via the action $r(m+U)=(r m)+U$. The quotient $\operatorname{map} \varphi: M \rightarrow M / U$ is an $R$-homomorphism, and this allows us to talk about $I / J$ for ideals $I$ and $J$ of a ring $R$.

Example 8.8. (1) The quotient group $\mathbb{Z} / 6 \mathbb{Z}$ is a $\mathbb{Z}$-module. Note that $2(3+6 \mathbb{Z})=6+6 \mathbb{Z}=0$ in $\mathbb{Z} / 6 \mathbb{Z}$, hence multiplication of non-zero elements of a module by non-zero scalars may result in zero. This is in contrast to the situation in vector spaces.
(2) Let $K$ be a field. Then $K$ is a $K[x]$-module, via $\pi: K[x] \rightarrow K[x] /\langle x\rangle$, which sends $P(x)$ to $P(0)$. Then the multiplication $P(x) \cdot \alpha$ for $P(x) \in K[x]$ and $\alpha \in K$ is simply given by $P(0) \alpha \in K$.

For a general $R$-homomorphism $\varphi: M \rightarrow N$, we can define $\operatorname{Ker} \varphi$ and $\operatorname{Im} \varphi$ in the usual way, and these are submodules of $M$ and $N$ respectively.
Definition 8.9. The cokernel of an $R$-homomorphism $\varphi: M \rightarrow N$ is the set

$$
\text { Coker } \varphi=N / \operatorname{Im} \varphi \text {. }
$$

Let $U, V$ be submodules of an $R$-module $M$. Then the set

$$
U+V=\{u+v: u \in U, v \in V\}
$$

is also a submodule of $M$. This is used in the following theorem.
Theorem 8.10 (Isomorphism theorems). Let $R$ be a ring and $M, N$ be $R$-modules. We have the following:
(i) if $\varphi: M \rightarrow N$ is an $R$-module homomorphism then

$$
M / \operatorname{Ker} \varphi \cong \operatorname{Im} \varphi ;
$$

(ii) if $L \subseteq M \subseteq N$ are submodules then

$$
(N / L) /(M / L) \cong N / M
$$

via the map $(m+L)+M / L \mapsto m+M$;
(iii) if $N$ is a module and $L, M$ are submodules then

$$
M /(M \cap L) \cong(M+L) / L
$$

via the map $m+M \cap L \mapsto m+L$.
These isomorphisms are canonical (i.e. require no choices in their definition).
Proof. Exercise Sheet.
Definition 8.11. Let $R$ be a ring and $M$ an $R$-module. Let $\Gamma$ be a subset of $M$. The submodule of $M$ generated by $\Gamma$, denoted $\langle\Gamma\rangle$ or $\sum_{g \in \Gamma} R g$, is the set

$$
\langle\Gamma\rangle=\left\{\sum_{i=1}^{n} r_{i} g_{i}: n \geqslant 1, r_{i} \in R, g_{i} \in \Gamma\right\} .
$$

The module $M$ is finitely generated if there exists a finite set $\Gamma \subseteq M$ such that $\langle\Gamma\rangle=M$.

Example 8.12. (1) Let $R$ be a ring and $I \subseteq R$ an ideal, then the $R$-module $R / I$ is finitely generated. In fact it is cyclic, i.e. generated by one element, namely $1+I$.
(2) If $R$ is an integral domain and $0 \neq f \in R$, then

$$
R\left[\frac{1}{f}\right]=R+R \frac{1}{f}+R \frac{1}{f^{2}}+\ldots
$$

is usually not finitely generated as an $R$-module.
(3) Let $\Gamma=\left\{x, x^{2}, x^{3}, \ldots,\right\} \subseteq K[x]$. Then $\langle\Gamma\rangle=\langle x\rangle$.

## 9 Nakayama's Lemma

Nakayama's lemma (also known as NAK, where the letters stand for Nakayama-AzumayaKrull) is an important tool in algebraic geometry. In particular it gives a precise definition of what it means for a module to be minimally generated (over a local ring).

Definition 9.1. A minimal generating set for an $R$-module $M$ is a subset $\Gamma \subseteq M$ such that $\Gamma$ generates $M$ but no proper subset of $\Gamma$ generates $M$.
Example 9.2. Consider $\mathbb{Z}_{6}=\mathbb{Z} / 6 \mathbb{Z}$, then $\{1+6 \mathbb{Z}\}$ and $\{2+6 \mathbb{Z}, 3+6 \mathbb{Z}\}$ are both minimal generating sets. Contrast this with vector spaces, where the number of elements in any two minimal generating sets of a given vector space are equal.
Theorem 9.3 (Nakayama's Lemma - NAK). Let $M$ be a finitely generated $R$-module, and $I \subseteq J(R)$ an ideal of $R$. If $M=I M$, then $M=0$.
Proof. Suppose $M \neq 0$. Since $M$ is finitely generated there exists a finite minimal generating set $\Gamma=\left\{g_{1}, \ldots, g_{n}\right\}$ say. Now $M=I M \Longrightarrow g_{1} \in I M$, so there exists $a_{1}, \ldots, a_{n} \in I$ such that

$$
g_{1}=\sum_{i=1}^{n} a_{i} g_{i}
$$

and so

$$
\left(1-a_{1}\right) g_{1}=\sum_{i=2}^{n} a_{i} g_{i}
$$

But $a_{1} \in I \subseteq J(R)$, so by Lemma 7.9,1-a is a unit of $R$. Thus

$$
g_{1}=\left(1-a_{1}\right)^{-1} \sum_{i=2}^{n} a_{i} g_{i}
$$

and $\left\{g_{2}, \ldots, g_{n}\right\}$ is a generating set for $M$ strictly smaller than $\Gamma$, a contradiction.
Corollary 9.4. Let $M$ be a finitely generated $R$-module and $N \subseteq M$ a submodule. Let also $I \subseteq J(R)$ be an ideal of $R$. Then $M=N+I M \Longrightarrow M=N$.

Proof. Take the equality $M=N+I M$ and quotient both sides by the submodule $N$ to obtain $M / N=(N+I M) / N$. By Theorem 8.10, we have $(N+I M) / N \cong I M /(N \cap I M)$. Now the map

$$
\begin{aligned}
I M & \rightarrow I(M / N) \\
\sum_{i=1}^{n} a_{i} m_{i} & \mapsto \sum_{i=1}^{n} a_{i}\left(m_{i}+N\right)
\end{aligned}
$$

is a surjective $R$-module homomorphism, and its kernel is $(I M) \cap N$. Therefore

$$
I(M / N) \cong I M /(I M \cap N) \cong(N+I M) / N
$$

Therefore we have $M / N=I(M / N)$. Since $M$ is finitely generated so too is $M / N$, and hence by Nakayama's Lemma we have $M / N=0$, i.e. $M=N$.

Example 9.5. Consider $K[x, y]$ for some field $K$ and let $\mathfrak{m}=\langle x, y\rangle$. Let $R=K[x, y]_{\mathfrak{m}}$, the localisation at the ideal $\mathfrak{m}$. Then $R$ is a local ring, with maximal ideal $\mathfrak{m}_{\mathfrak{m}}$. We will show that the ideal

$$
I=\left\langle x+x^{2} y+3 y^{2}+x^{4}, y+2 y^{3}+y^{4}+4 x^{7}\right\rangle_{\mathfrak{m}} \subseteq R
$$

is equal to $\mathfrak{m}_{\mathfrak{m}}$. Note first that since $R$ is local it has a unique maximal ideal, hence $J(R)=\mathfrak{m}_{\mathfrak{m}}$. Now

$$
\begin{aligned}
I+\mathfrak{m}_{\mathfrak{m}} \mathfrak{m}_{\mathfrak{m}} & =\left\langle x+x^{2} y+3 y^{2}+x^{4}, y+2 y^{3}+y^{4}+4 x^{7}, x^{2}, x y, y^{2}\right\rangle_{\mathfrak{m}} \\
& =\left\langle x, y, x^{2}, x y, y^{2}\right\rangle_{\mathfrak{m}} \\
& =\langle x, y\rangle_{\mathfrak{m}} \\
& =\mathfrak{m}_{\mathfrak{m}} .
\end{aligned}
$$

So by Nakayama's Lemma, $I=\mathfrak{m}_{\mathfrak{m}}$.
Recall from earlier that we had an issue with minimal generating sets for modules, in that the number of elements in such a set is not well defined. Nakayama's Lemma allows us to fix this in certain cases.

Theorem 9.6. Let $(R, \mathfrak{m})$ be a local ring and $M$ a finitely generated $R$-module. If $\Gamma \subseteq M$ is a set of elements whose images in $M / \mathfrak{m} M$ form a basis of $M / \mathfrak{m} M$ as an $R / \mathfrak{m}$-vector space, then $\Gamma$ is a minimal generating set of $M$ as an $R$-module.

Proof. As $M / \mathfrak{m} M$ is generated by the images of the elements of $\Gamma$, we have $M=\langle\Gamma\rangle+\mathfrak{m} M$. So by Corollary 9.4 to Nakayama's Lemma, we have $M=\langle\Gamma\rangle$. If $\Gamma^{\prime} \subsetneq \Gamma$, then $\left\langle\Gamma^{\prime}\right\rangle+\mathfrak{m} M \neq\langle\Gamma\rangle+\mathfrak{m} M=$ $M$, and so $\Gamma^{\prime}$ is not a generating set.

## 10 Exact sequences

Definition 10.1. A sequence of $R$-modules and $R$-module homomorphisms

$$
\cdots \longrightarrow M_{0} \xrightarrow{f_{1}} M_{1} \xrightarrow{f_{2}} M_{2} \longrightarrow \cdots \xrightarrow{f_{n}} M_{n} \longrightarrow \cdots
$$

is called exact at $M_{i}$ if $\operatorname{Ker} f_{i+1}=\operatorname{Im} f_{i}$. A sequence which is exact at $M_{i}$ for all $i$ is called an exact sequence.

Example 10.2. (i) The sequence $0 \longrightarrow L \stackrel{f}{\longrightarrow} M$ is exact if and only if $f$ is injective.
(ii) The sequence $M \stackrel{g}{\longrightarrow} N \longrightarrow 0$ is exact if and only if $g$ is surjective.
(iii) The sequence $0 \longrightarrow M \xrightarrow{g} N \longrightarrow 0$ is exact if and only if $g$ is an isomorphism.

Definition 10.3. A short exact sequence is an exact sequence of the form

$$
0 \longrightarrow L \xrightarrow{f} M \xrightarrow{g} N \longrightarrow 0 .
$$

Remark. This is equivalent to insisting that $f$ is injective, $g$ is surjective and $\operatorname{Ker} g=\operatorname{Im} f$.
Short exact sequences appear in many different sub-branches of algebra, and are very powerful objects.
Example 10.4. (i) Let $R$ be a ring, $M$ an $R$-module and $N \subseteq M$ a submodule. Then

$$
0 \longrightarrow N \xrightarrow{i} M \xrightarrow{\pi} M / N \longrightarrow 0
$$

where $i$ is the natural inclusion map and $\pi$ is the canonical quotient map, is a short exact sequence.
(ii) Any long exact sequence can be split into short exact sequences. Let

$$
\cdots \xrightarrow{f_{i-1}} M_{i-1} \xrightarrow{f_{i}} M_{i} \xrightarrow{f_{i+1}} M_{i+1} \rightarrow \cdots
$$

be an exact sequence, that is $\operatorname{Im}\left(f_{i}\right)=\operatorname{Ker}\left(f_{i+1}\right)$ for all $i$. Then

$$
0 \rightarrow \operatorname{Ker}\left(f_{i+1}\right) \rightarrow M_{i} \rightarrow M_{i} / \operatorname{Im}\left(f_{i}\right)=\operatorname{Coker}\left(f_{i}\right) \rightarrow 0
$$

is a short exact sequence.
(iii) Let $K$ be a field and

$$
0 \longrightarrow L \xrightarrow{f} M \xrightarrow{g} N \longrightarrow 0
$$

be a short exact sequence of $K$-modules. Then each module is a $K$-vector space, and using facts from linear algebra we have

$$
\begin{aligned}
\operatorname{dim}_{K} M & =\operatorname{dim}_{K} \operatorname{Ker} g+\operatorname{dim}_{K} \operatorname{Im} g \\
& =\operatorname{dim}_{K} \operatorname{Im} f+\operatorname{dim}_{K} N \\
& =\operatorname{dim}_{K} L+\operatorname{dim}_{K} N .
\end{aligned}
$$

More generally, if

$$
0 \longrightarrow M_{0} \xrightarrow{f_{1}} M_{1} \xrightarrow{f_{2}} M_{2} \longrightarrow \cdots \xrightarrow{f_{n}} M_{n} \longrightarrow 0
$$

is an exact sequence of $K$-vector spaces, then $\sum_{i=0}^{n}(-1)^{i} \operatorname{dim}_{K} M_{i}=0$.
Remark 10.5. One can also consider (exact) sequences of other objects, sequences $\cdots \rightarrow A_{0} \xrightarrow{f_{1}}$ $A_{1} \xrightarrow{f_{2}} \cdots$ of abelian groups, where the $f_{i}$ are group homomorphisms.
Definition 10.6. Let $A, B, C, D$ be $R$-modules and let $\alpha, \beta, \gamma, \delta$ be $R$-module homomorphisms. Then the diagram

is commutative (or: the diagram commutes) if $\beta \circ \alpha=\delta \circ \gamma$.
The following lemma is a typical example for statements in homological algebra. We will prove it with diagram chasing.

Theorem 10.7 (Snake Lemma). Suppose the following commutative diagram of $R$-modules and $R$ module homomorphisms

has exact rows. Then there exists a homomorphism $\delta: \operatorname{Ker} \gamma \rightarrow \operatorname{Coker} \alpha$ such that

$$
\operatorname{Ker} \alpha \longrightarrow \operatorname{Ker} \beta \longrightarrow \operatorname{Ker} \gamma \stackrel{\delta}{\longrightarrow} \text { Coker } \alpha \longrightarrow \text { Coker } \beta \longrightarrow \text { Coker } \gamma
$$

is exact.
Furthermore, if $f$ is injective then so too is $\operatorname{Ker} \alpha \rightarrow \operatorname{Ker} \beta$, and if $g^{\prime}$ is surjective then so too is Coker $\beta \rightarrow$ Coker $\gamma$.

The name of this theorem comes from the following diagram:


Proof. We will first define all of the necessary maps, then prove exactness at each site.
The map $\left.f\right|_{\operatorname{Ker} \alpha}: \operatorname{Ker} \alpha \rightarrow \operatorname{Ker} \beta$ is given by the restriction of $f$ to $\operatorname{Ker} \alpha$. Note that if $\ell \in \operatorname{Ker} \alpha$ then $\beta(f(\ell))=f^{\prime}(\alpha(\ell))=0$ by the commutativity of the diagram. Therefore $f(\operatorname{Ker} \alpha) \subseteq \operatorname{Ker} \beta$. That this is a $R$-homomorphism follows from the fact that $f$ itself is. Similarly the map $\left.g\right|_{K e r ~} \beta$ : $\operatorname{Ker} \beta \rightarrow \operatorname{Ker} \gamma$ is given by the restriction of $g$ to $\operatorname{Ker} \beta$.
The map $\bar{f}:$ Coker $\alpha \rightarrow$ Coker $\beta$ is induced from $f^{\prime}$, by setting $\bar{f}\left(\ell^{\prime}+\operatorname{Im} \alpha\right)=f^{\prime}\left(\ell^{\prime}\right)+\operatorname{Im} \beta$. This is well defined, as if $\ell_{1}^{\prime}+\operatorname{Im} \alpha=\ell_{2}^{\prime}+\operatorname{Im} \alpha$ then $\ell_{1}^{\prime}-\ell_{2}^{\prime} \in \operatorname{Im} \alpha$, so $\ell_{1}^{\prime}-\ell_{2}^{\prime}=\alpha(\ell)$ for some $\ell \in L$. Then

$$
\begin{aligned}
f^{\prime}\left(\ell_{1}^{\prime}\right)-f^{\prime}\left(\ell_{2}^{\prime}\right) & =f^{\prime}\left(\ell_{1}^{\prime}-\ell_{2}^{\prime}\right) \\
& =f^{\prime}(\alpha(\ell)) \\
& =\beta(f(\ell)) \\
& \in \operatorname{Im} \beta,
\end{aligned}
$$

so $f^{\prime}\left(\ell_{1}^{\prime}\right)+\operatorname{Im} \beta=f^{\prime}\left(\ell_{2}^{\prime}\right)+\operatorname{Im} \beta$. That $\bar{f}$ is a homomorphism follows from the fact that $f^{\prime}$ is. We similarly define $\bar{g}$ : Coker $\beta \rightarrow$ Coker $\gamma$.
We now construct the connecting homomorphism $\delta: \operatorname{Ker} \gamma \rightarrow$ Coker $\alpha$ by a process known as "diagram chasing". Take $n \in \operatorname{Ker} \gamma \subseteq N$. Since $g$ is surjective, there exists some $m \in M$ such that $n=g(m)$. Then

$$
\begin{aligned}
0 & =\gamma(n) \\
& =\gamma(g(m)) \\
& =g^{\prime}(\beta(m))
\end{aligned}
$$

by the commutativity of the diagram, so $\beta(m) \in \operatorname{Ker} g^{\prime}$. By the exactness of rows, $\operatorname{Ker} g^{\prime}=\operatorname{Im} f^{\prime}$, so $\beta(m)=f^{\prime}\left(\ell^{\prime}\right)$ for some $\ell^{\prime} \in L^{\prime}$. We then define

$$
\delta(n)=\ell^{\prime}+\operatorname{Im} \alpha \in \text { Coker } \alpha
$$

We must show that this is well defined. Since $f^{\prime}$ is injective, the only ambiguity in our process lies in our choice of $m$. Suppose then that $g\left(m_{1}\right)=g\left(m_{2}\right)=n$, and $\ell_{1}^{\prime}, \ell_{2}^{\prime} \in L^{\prime}$ are the unique elements such that $\beta\left(m_{1}\right)=f^{\prime}\left(\ell_{1}^{\prime}\right)$ and $\beta\left(m_{2}\right)=f^{\prime}\left(\ell_{2}^{\prime}\right)$. We must show that $\ell_{1}^{\prime}-\ell_{2}^{\prime} \in \operatorname{Im} \alpha$. Note then that $m_{1}-m_{2} \in \operatorname{Ker} g$, and so by exactness of rows is equal to $f(\ell)$ for some $\ell \in L$. Therefore $\beta\left(m_{1}-m_{2}\right)=\beta(f(\ell))=f^{\prime}(\alpha(\ell))$. By the injectivity of $f^{\prime}$, we then see that $\alpha(\ell)=\ell_{1}^{\prime}-\ell_{2}^{\prime}$. That $\delta$ is a homomorphism is left as an easy exercise.

We now prove exactness at each site.
The composition $\left.\left.g\right|_{\text {Ker } \beta \circ f} f\right|_{\operatorname{Ker} \alpha}=0$ follows from the fact that $\operatorname{Im} f=\operatorname{Ker} g$, therefore $\left.\operatorname{Im} f\right|_{\operatorname{Ker} \alpha} \subseteq$ $\left.\operatorname{Ker} g\right|_{\operatorname{Ker} \beta}$. Suppose now that $m \in \operatorname{Ker} \beta$ with $\left.g\right|_{\operatorname{Ker} \beta}(m)=0$. Then $g(m)=0$ so $m \in \operatorname{Ker} g=\operatorname{Im} f$, say $m=f(\ell)$, and it remains to show that $\ell \in \operatorname{Ker} \alpha$. But

$$
\begin{aligned}
f^{\prime}(\alpha(\ell)) & =\beta(f(\ell)) \\
& =\beta(m) \\
& =0
\end{aligned}
$$

as $m \in \operatorname{Ker} \beta$, and since $f^{\prime}$ is injective we must have $\alpha(\ell)=0$.
For exactness at $\operatorname{Ker} \gamma$, we first calculate $\delta\left(\left.g\right|_{\operatorname{Ker} \beta}(m)\right)$ for $m \in \operatorname{Ker} \beta$. Following our construction of $\delta$ above, we have $g_{\operatorname{Ker} \beta}(m)=g(m)$, and so $\ell^{\prime}$ is chosen so that $\beta(m)=f^{\prime}\left(\ell^{\prime}\right)$. But $\beta(m)=0$, so by the injectivity of $f^{\prime}$ we also have $\delta\left(\left.g\right|_{\operatorname{Ker} \beta}(m)\right)=0$ and hence $\left.\operatorname{Im} g\right|_{\operatorname{Ker} \beta} \subseteq \operatorname{Ker} \delta$. Conversely if $n \in \operatorname{Ker} \gamma$ is such that $\delta(n)=0$, then the corresponding $\ell^{\prime}$ is in $\operatorname{Im} \alpha$, say $\ell^{\prime}=\alpha(\ell)$. Therefore if $m$ is such that $n=g(m)$, we have $\beta(m)=f^{\prime}\left(\alpha\left(\ell^{\prime}\right)\right)=\beta(f(\ell))$, and hence $m-f(\ell) \in \operatorname{Ker} \beta$. Then $\left.g\right|_{\text {Ker } \beta}(m-f(\ell))=g(m)-g(f(\ell))=n$.
For exactness at Coker $\alpha$, note that $\bar{f}(\delta(n))=f^{\prime}\left(\ell^{\prime}\right)+\operatorname{Im} \beta=\beta(m)+\operatorname{Im} \beta=0$ in Coker $\beta$. Therefore $\operatorname{Im} \delta \subseteq \operatorname{Ker} \bar{f}$. Conversely if $l^{\prime}+\operatorname{Im} \alpha \in \operatorname{Coker} \alpha$ is such that $\bar{f}\left(l^{\prime}+\operatorname{Im} \alpha\right)=0$, then $f^{\prime}\left(\ell^{\prime}\right) \in \operatorname{Im} \beta$, say $f^{\prime}\left(\ell^{\prime}\right)=\beta(m)$. But then $\delta(g(m))=\ell^{\prime}+\operatorname{Im} \alpha$.
Finally, for exactness at Coker $\beta$ we see first that $\bar{g}\left(\bar{f}\left(\ell^{\prime}+\operatorname{Im} \alpha\right)\right)=\bar{g}\left(f^{\prime}\left(\ell^{\prime}\right)+\operatorname{Im} \beta\right)=g^{\prime}\left(f^{\prime}\left(\ell^{\prime}\right)\right)+$ $\operatorname{Im} \gamma=0$ since $g^{\prime} \circ f^{\prime}=0$. Therefore $\operatorname{Im} \bar{f} \subseteq \operatorname{Ker} \bar{g}$. Conversely, if $m^{\prime}+\operatorname{Im} \beta \in \operatorname{Coker} \beta$ is such that $\bar{g}\left(m^{\prime}+\operatorname{Im} \beta\right)=0$, then $g^{\prime}\left(m^{\prime}\right) \in \operatorname{Im} \gamma$, say $g^{\prime}\left(m^{\prime}\right)=\gamma(n)$. Since $g$ is surjective, there is some $m \in M$ such that $g(m)=n$, so $g^{\prime}\left(m^{\prime}\right)=\gamma(g(m))$. Commutativity of the diagram then gives $g^{\prime}\left(m^{\prime}\right)=g^{\prime}(\beta(m))$, so $m^{\prime}-\beta(m) \in \operatorname{Ker} g^{\prime}=\operatorname{Im} f^{\prime}$, say $m^{\prime}-\beta(m)=f^{\prime}\left(\ell^{\prime}\right)$. But now $f\left(\ell^{\prime}+\operatorname{Im} \alpha\right)=f^{\prime}\left(\ell^{\prime}\right)+\operatorname{Im} \beta=m^{\prime}-\beta(m)+\operatorname{Im} \beta=m^{\prime}+\operatorname{Im} \beta$.
We leave the last statement as an exercise.
Example 10.8. We reprove part (ii) of Theorem 8.10. Let $L \subseteq M \subseteq N$ be a sequence of submodules and consider the following diagram:


The maps $f, g$ and $f^{\prime}, g^{\prime}$ are pairs of inclusion and quotient maps, so the rows are short exact sequences. We have $\alpha: M \rightarrow M / L$ and $\beta: N \rightarrow N / L$ also quotient homomorphisms, and for all $m \in M$

$$
\begin{aligned}
\beta(f(m)) & =\beta(m) \\
& =m+L \\
& =f^{\prime}(m+L) \text { since } m \in M \\
& =f^{\prime}(\alpha(m))
\end{aligned}
$$

so the first square commutes. Now define $\gamma: N / M \rightarrow(N / L) /(M / L)$ by $\gamma(n+M)=(n+L)+$ $M / L$. This is well defined since if $n+M=n^{\prime}+M$ then $n-n^{\prime} \in M$ so

$$
\begin{aligned}
\gamma(n)-\gamma\left(n^{\prime}\right) & =((n+L)+M / L)-\left(\left(n^{\prime}+L\right)+M / L\right) \\
& =\left(n-n^{\prime}+L\right)+M / L \\
& =M / L=0_{(N / L) /(M / L)} \text { since } n-n^{\prime} \in M
\end{aligned}
$$

It is also a homomorphism (easy check since it is the composition of two quotient maps). Finally we check that the diagram commutes: for all $n \in N$ we have

$$
\begin{aligned}
\gamma(g(n)) & =\gamma(n+M) \\
& =(n+L)+M / L, \text { and } \\
g^{\prime}(\beta(n)) & =g^{\prime}(n+L) \\
& =(n+L)+M / L .
\end{aligned}
$$

By the Snake Lemma, we therefore have an exact sequence

$$
0 \rightarrow \text { Ker } \alpha \rightarrow \text { Ker } \beta \rightarrow \text { Ker } \gamma \rightarrow \text { Coker } \alpha \rightarrow \text { Coker } \beta \rightarrow \text { Coker } \gamma \rightarrow 0
$$

Clearly $\operatorname{Ker} \alpha=\operatorname{Ker} \beta=L$ and Coker $\alpha=\operatorname{Coker} \beta=0$. Therefore our exact sequence is equal to

$$
0 \rightarrow L \rightarrow L \rightarrow \text { Ker } \gamma \rightarrow 0 \rightarrow 0 \rightarrow \text { Coker } \gamma \rightarrow 0
$$

By exactness we immediately see that $\operatorname{Ker} \gamma=\operatorname{Coker} \gamma=0$. Thus $\gamma$ is both injective and surjective, so is an isomorphism between $N / M$ and $(N / L) /(M / L)$.

## 11 Free modules

Let $R$ be a ring, $\Lambda$ a set and $M_{\lambda}$ an $R$-module for each $\lambda \in \Lambda$.
Definition 11.1. The direct product of $\left\{M_{\lambda}\right\}_{\lambda \in \Lambda}$, denoted $\prod_{\lambda \in \Lambda} M_{\lambda}$, consists of all sequences $\left(m_{\lambda}\right)_{\lambda \in \Lambda}$ with $m_{\lambda} \in M_{\lambda}$ for each $\lambda \in \Lambda$. This is a module, with addition

$$
\left(m_{\lambda}\right)_{\lambda \in \Lambda}+\left(n_{\lambda}\right)_{\lambda \in \Lambda}=\left(m_{\lambda}+n_{\lambda}\right)_{\lambda \in \Lambda}
$$

and for any $r \in R$,

$$
r\left(m_{\lambda}\right)_{\lambda \in \Lambda}=\left(r m_{\lambda}\right)_{\lambda \in \Lambda}
$$

The direct sum of $\left\{M_{\lambda}\right\}_{\lambda \in \Lambda}$, denoted $\bigoplus_{\lambda \in \Lambda} M_{\lambda}$, consists of all sequences $\left(m_{\lambda}\right)_{\lambda \in \Lambda}$ with $m_{\lambda} \in M_{\lambda}$ for each $\lambda \in \Lambda$, and all but finitely many of the $m_{\lambda}$ are zero. This is again a module, with addition and scalar multiplication as before.

Note that if $\Lambda$ is finite then $\prod_{\lambda \in \Lambda} M_{\lambda}=\bigoplus_{\lambda \in \Lambda} M_{\lambda}$. For instance, $\mathbb{R} \oplus \mathbb{R} \cong \mathbb{R}^{2}$.
Remark 11.2. The direct sum/product can be defined categorically and are given by universal properties.

Proposition 11.3. If $U, V$ are submodules of $M$, then $M=U \oplus V \Longleftrightarrow M=U+V$ and $U \cap V=\{0\}$.
Proof. Exercise.
Remark. Care needs to be taken when dealing with direct products. For instance, for rings $R$ and $S$ their direct product $R \times S$ has identity (1,1). Then the natural map $\varphi: R \rightarrow R \times S$ given by $\varphi(r)=(r, 0)$ is not a ring homomorphism, since $\varphi(1)=(1,0) \neq(1,1)$.
Definition 11.4. An $R$-module is called free if it is isomorphic to $\bigoplus_{\lambda \in \Lambda} R$ for some set $\Lambda$. We adopt the convention the the zero module is free, with index set $\Lambda=\varnothing$.

Example 11.5. (i) $R^{n}=R \oplus R \oplus \cdots \oplus R$ is clearly free.
(ii) The ring of $m \times n$ matrices over a ring $R$ is free and isomorphic to $R^{m n}$.
(iii) The polynomial ring $R[X]$ is free, as $R[X] \cong R \oplus R X \oplus R X^{2} \oplus \ldots$.

Recall that in contrast to vector spaces, not every module has a basis. However free modules do.

Proposition 11.6. An $R$-module is free if and only if there exists a set of generators $\left\{m_{\lambda}\right\}_{\lambda \in \Lambda}$ of $M$ such that whenever $r_{1} m_{\lambda_{1}}+\ldots r_{n} m_{\lambda_{n}}=0$ with $r_{i} \in R$ and $\lambda_{i} \in \Lambda$ for all $i$, we have $r_{1}=\cdots=r_{n}=0$.
Proof. The "only if" direction is clear.
Conversely, assume we have a set of generators as above and define a map

$$
\begin{aligned}
& \varphi: \bigoplus_{\lambda \in \Lambda} R \rightarrow M \\
&\left(r_{\lambda}\right)_{\lambda \in \Lambda} \mapsto \sum_{\lambda \in \Lambda} r_{\lambda} m_{\lambda} .
\end{aligned}
$$

It is then straightforward to check that this is an isomorphism of $R$-modules.
Definition 11.7. A set of generators as in Proposition 11.6 is called a free basis, or just a basis. The rank of a free module is the cardinality of $\Lambda$, equivalently the number of basis elements.
Example 11.8. (i) $1, X, X^{2}, \ldots$ is a basis of $R[X]$.
(ii) The rank of $R^{n}$ is $n$.
(iii) A $K$-vector space has a basis and so is a free K-module.
(iv) Consider the maximal ideal $\mathfrak{m}=\langle x, y\rangle$ of $R=K[x, y]$. This is generated by two elements but is not free, for instance as $-y x+x y=0$ is a non-trivial dependence relation. However, the module of relations of $\mathfrak{m}$ is freely generated by one element, $(-y, x)$. Thus we get an exact sequence of $R$-modules

$$
0 \longrightarrow R \longrightarrow R^{2} \longrightarrow \mathfrak{m} \longrightarrow 0
$$

This exact sequence can be completed to the Koszul complex of $K$ :

$$
0 \longrightarrow R \longrightarrow R^{2} \longrightarrow R \longrightarrow K \longrightarrow 0
$$

This is what is called a free resolution of the $R$-module $K$. In order to understand the structure of non-free modules $M$, one can study resolutions of $M$ by free modules.
(v) $\mathbb{Z}_{2}$ is not free as a $\mathbb{Z}$-module, since it is generated by $1+2 \mathbb{Z}$ but $2(1+2 \mathbb{Z})=2+2 \mathbb{Z}=0_{\mathbb{Z}_{2}}$, so this is a non-trivial dependence relation.
Proposition 11.9. Let $R$ be a ring and $M$ an $R$-module. Then there exists a free module $F$ and a surjective homomorphism of $R$ modules $\varphi: F \rightarrow M$. Furthermore if $M$ is finitely generated then $F$ can be chosen to have finite rank.

Proof. Any $R$-module can be written as $\langle\Gamma\rangle$ for some $\Gamma \subseteq M$, for instance by setting $\Gamma=M$. Then let $F$ be the free module with basis $\Gamma$. Now define

$$
\begin{gathered}
\varphi: F \rightarrow M \\
\left(r_{g}\right)_{g \in \Gamma} \mapsto \sum_{g \in \Gamma} r_{g} g .
\end{gathered}
$$

Note that this sum is finite since $F$ is a direct sum of copies of $R$. It is an easy exercise to see that this is a surjective $R$-module homomorphism.
If $M$ is finitely generated, say by $\left\{m_{g}\right\}_{g \in \Gamma}$ then we similarly define $F$ to be the free module with finite basis $\Gamma$, and $\varphi: F \rightarrow M$ by $\varphi\left(\left(r_{g}\right)_{g \in \Gamma}\right)=\sum_{g \in \Gamma} r_{g} m_{g}$. It is again easy to check that this is a surjective homomorphism.

Example 11.10. Let $M_{1}, \ldots, M_{n}$ be $R$-modules. Then the sequence

$$
0 \longrightarrow M_{1} \longrightarrow M_{1} \oplus \cdots \oplus M_{n} \longrightarrow M_{2} \oplus M_{3} \oplus \cdots \oplus M_{n} \longrightarrow 0
$$

is exact.

Proposition 11.11. Let $L, M, N$ be $R$-modules and let

$$
0 \rightarrow L \xrightarrow{\alpha} M \xrightarrow{\beta} N \rightarrow 0
$$

be a short exact sequence. Then the following are equivalent:
(i) There exists an isomorphism $M \cong L \oplus N$ under which $\alpha$ is given by $l \mapsto(l, 0)$ and $\beta$ as $(l, n) \mapsto n$.
(ii) There exists a section of $\beta$, that is, a map $s: N \rightarrow M$ such that $\beta s=\operatorname{Id}_{N}$.
(iii) There exists a retraction for $\alpha$, that is, a map $r: M \rightarrow L$ such that $r \alpha=\operatorname{Id}_{L}$.

Definition 11.12. If any of the three equivalent condition of the above proposition is satisfied, then the short exact sequence

$$
0 \rightarrow L \xrightarrow{\alpha} M \xrightarrow{\beta} N \rightarrow 0
$$

is called a split exact sequence.

## Proof. Exercise.

Example 11.13. (1) For finite dimensional $K$-vector spaces, every short exact sequence is split.
(2) The short exact sequence

$$
0 \rightarrow\langle x\rangle \xrightarrow{\text { incl }} K[x] \xrightarrow{\pi} K \rightarrow 0
$$

is nonsplit as a sequence of $K[x]$-modules. (See this by trying to constract a section $K \rightarrow K[x]$ !)

## 12 Noetherian rings and modules

Being finitely generated is obviously a good property for a module to have. But if $M$ is a finitely generated $R$-module then there is no guarantee that its submodules will be.

Example 12.1. Let $R=K\left[x_{1}, x_{2}, x_{3}, \ldots\right]$. Then $R$ is an $R$-module and is finitely generated by $\{1\}$. However the submodule $\left\langle x_{1}, x_{2}, x_{3} \ldots\right\rangle$ is not.

This motivates the following:
Definition 12.2. A module $M$ is called a Noetherian ${ }^{1}$ module if every submodule of $M$ is finitely generated. A ring $R$ is called a Noetherian ring if it is a Noetherian module over itself (i.e. all ideals are finitely generated).

Examples are hard to give without a bit of extra theory, so we present this first.
Theorem 12.3. Let $M$ be an $R$-module. Then the following are equivalent:
(i) all submodules of $M$ are finitely generated;
(ii) $M$ satisfies the ascending chain condition (ACC), i.e. every chain of submodules

$$
M_{1} \subseteq M_{2} \subseteq M_{3} \subseteq \ldots
$$

of $M$ is stationary, that is there exists some $N$ with $M_{n}=M_{N}$ for all $n \geqslant N$;
(iii) every non-empty set of submodules of $M$ has a maximal element.

[^0]Proof. (i) $\Longrightarrow$ (ii) : The union $\bigcup_{i} M_{i}$ is a submodule of $M$, so is finitely generated by assumption. Each of these generators must lie in some $M_{j}$, and taking $N$ to be the maximum of these $j$ we have $\bigcup_{i} M_{i}=M_{N}$. Hence $M_{n}=M_{N}$ for all $n \geqslant N$.
(ii) $\Longrightarrow$ (iii) : Let $S$ be a non-empty set of submodules of $M$ and suppose $S$ has no maximal element. Since $S$ is non-empty we can take some $M_{1} \in S$. Since $M_{1}$ is not maximal we can find some $M_{2} \in S$ with $M_{1} \subsetneq M_{2}$. Repeating this argument we can construct inductively a nonstationary ascending chain of submodules of $M$, contradicting (ii).
(iii) $\Longrightarrow(i)$ : Let $U$ be a submodule of $M$ and $S$ the set of finitely generated submodules of $U$. This is non-empty as it contains the zero module, so has a maximal element $U^{\prime}=\left\langle u_{1}, \ldots, u_{n}\right\rangle$. Now take any $v \in U$, then $U^{\prime}+\langle v\rangle=\left\langle u_{1}, \ldots, u_{n}, v\right\rangle$ is a finitely generated submodule of $U$, so by maximality must equal $U^{\prime}$. Hence $U=U^{\prime}$ is finitely generated.

We can now give some examples of Noetherian rings and modules.
Example 12.4. (i) Let $R$ be a field, then the only ideals of $R$ are $R$ and $\{0\}$ which are finitely generated. Therefore $R$ is a Noetherian ring.
(ii) Modules and rings with a finite number of elements are Noetherian.
(iii) Any principal ideal domain is a Noetherian ring. Therefore $\mathbb{Z}, \mathbb{Z}[i]$ and $K[x]$ ( $K$ a field) are Noetherian rings (as they are Euclidean domains).
(iv) Finite dimensional $K$-vector spaces are Noetherian $K$-modules, since any subspace (submodule) has a finite basis.

Theorem 12.5. Let $0 \longrightarrow L \longrightarrow M \longrightarrow N \longrightarrow 0$ be an exact sequence of $R$-modules. Then $M$ is Noetherian if and only if both $L$ and $N$ are Noetherian.

Proof. Note that the property of being Noetherian is preserved by isomorphisms, thus it is sufficient to prove the theorem in the case $L \subseteq M$ and $N=M / L$. [One can prove this using the snake lemma. Look at the diagram of short exact sequences:

where $\gamma: N \rightarrow M / \alpha(L)$ is defined via: since $\beta$ is surjective, for any $n \in N$ there exists an $m \in M$ such that $\beta(m)=n$. Then set $\gamma(n)=m+\alpha(L)$. This is well-defined, since for any $m^{\prime} \in M$ with $\beta\left(m^{\prime}\right)=n$, one has that $m-m^{\prime} \in \operatorname{Ker}(\beta)$, which is equal to $\operatorname{Im}(\alpha)$, since the top sequence is exact. But this means that $m-m^{\prime} \in \alpha(L)$ and thus the cosets $m+\alpha(L)=m^{\prime}+\alpha(L)$ in $M / \alpha(L)$. For the bottom row note that $\alpha(L) \cong L$, since $\alpha$ is injective. The bottom row is exact by construction. It is easy to see that the diagram commutes, and then an application of the snake lemma yields the result.]
Suppose first that $M$ is Noetherian and let $L^{\prime}$ be a submodule of $L$. Then $L^{\prime}$ is a submodule of $M$ so is finitely generated, and hence $L$ is Noetherian. Next, any submodule $N^{\prime}$ of $M / L$ is of the form $M^{\prime}$ / $L$ for some submodule $M^{\prime}$ of $M$. Therefore $M^{\prime}$ is finitely generated, and reduction of these generators modulo $L$ shows that $N^{\prime}$ is also finitely generated.
Conversely suppose that both $L$ and $N$ are Noetherian and consider a submodule $M^{\prime} \subseteq M$. Then the submodules $M^{\prime} \cap L \subseteq L$ and $M^{\prime} / L \subseteq N$ are both finitely generated, say by $x_{1}, \ldots, x_{n}$ and $y_{1}+L, \ldots, y_{m}+L$ respectively. Now for any $m \in M^{\prime}$ we have $m+L=\left(b_{1} y_{1}+\cdots+b_{m} y_{m}\right)+L$ for some $b_{i} \in R$, thus $m-\left(b_{1} y_{1}+\cdots+b_{m} y_{m}\right) \in L$. But also $m, y_{1}, \ldots, y_{m} \in M^{\prime}$, so $m-\left(b_{1} y_{1}+\cdots+\right.$ $\left.b_{m} y_{m}\right)=a_{1} x_{1}+\cdots+a_{n} x_{n}$ for some $a_{i} \in R$. Hence $m=a_{1} x_{1}+\cdots+a_{n} x_{n}+b_{1} y_{1}+\cdots+b_{m} y_{m}$, and so $M^{\prime}$ is finitely generated. Therefore $M$ is Noetherian.

Proposition 12.6. Let $R$ be a Noetherian ring and $M$ an $R$-module. Then $M$ is Noetherian if and only if $M$ is finitely generated.

Proof. The "only if" direction is by definition.
Suppose $M$ is finitely generated, then there is a surjection $\varphi: R^{n} \rightarrow M$ for some $n \geqslant 0$. The sequence $0 \longrightarrow \operatorname{Ker} \varphi \longrightarrow R^{n} \longrightarrow M \longrightarrow 0$ is then exact, and since $R^{n}$ is Noetherian then so too is $M$ by Theorem 12.5.

Proposition 12.7. Let $R$ be a Noetherian ring.
(i) Let $I \subseteq R$ be an ideal. Then $R / I$ is a Noetherian ring.
(ii) Let $A \subseteq R$ be a multiplicatively closed subset. Then $A^{-1} R$ is a Noetherian ring.

Proof. (i) Let $J$ be an ideal or $R / I$. Its preimage under the canonical quotient map is finitely generated, therefore so too is $J$.
(ii) Similarly for an ideal $J$ of $A^{-1} R$, its preimage under the natural map $R \rightarrow A^{-1} R$ is finitely generated. Therefore so too is $J$.

Remark 12.8. One can also define Noetherian spaces: Let $X$ be a topological space. Then $X$ is called noetherian if every descending chain of closed subsets becomes stationary. In particular $X=$ $\mathbb{A}_{K}^{n}$ is a noetherian space, where one takes the closed subsets to be $V(I)$, where $I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ is an ideal. This topology is called Zariski topology. Since for ideal $I \subseteq J$ in $K\left[x_{1}, \ldots, x_{n}\right]$, one has $V(J) \subseteq V(I)$ (see part about algebraic geometry), one can show that a descending chain of closed subsets in $X$ corresponds to an ascending chain of ideals in $K\left[x_{1}, \ldots, x_{n}\right]$.

Remark 12.9. If an $R$-module $M$ satisfies the descending chain condition, that is, every descending chain of submodules $M_{1} \supseteq M_{2} \supseteq \cdots$ becomes stationary, then $M$ is called Artinian module. A ring $R$ is called Artinian if it is Artinian as a module over itself. This condition is much rarer than noetherian: if $R$ is Artinian, then it is also Noetherian. An example of an Artinian ring is $R=K[x] /\left\langle x^{n}\right\rangle$ for $n \geqslant 1$.
But on the other hand, take for example the polynomial ring $K[x]$ : here $\langle x\rangle \supsetneq\left\langle x^{2}\right\rangle \supsetneq\left\langle x^{3}\right\rangle \supsetneq \cdots$ is a strictly decreasing chain of ideals that never becomes stationary.

## 13 Hilbert's Basis Theorem

This theorem was proved by David Hilbert in 1890. It is fundamental for algebraic geometry and also important for practical computations, in particular, Gröbner basis calculations.

Theorem 13.1. If $R$ is Noetherian, then the polynomial ring $R[x]$ is Noetherian.
Remark 13.2. In the lecture I did a different proof, following Atiyah-Macdonald [1, p.81f]. The idea of both proofs is the same: take an ideal $I$ in $R[x]$ and look at the ideal generated by all the leading coefficients of polynomials in $I$. The leading coefficients are in $R$, so this ideal $l c(I)$ has to be finitely generated. Then look at the corresponding ideal $I^{\prime} \subseteq R[x]$ generated by all the polynomials, whose leading coefficient generate $l c(I)$. Show with a "division algorithm" that any element in I belongs to a finitely generated module (namely $I^{\prime}$ and the "remainders").

Proof. Suppose there exists an ideal $I \subseteq R[x]$ which is not finitely generated. Choose a sequence $f_{1}, f_{2}, f_{3}, \ldots$ of polynomials in $R[x]$ such that

$$
\begin{aligned}
& f_{1} \in I \\
& f_{2} \in I \backslash\left\langle f_{1}\right\rangle \\
& f_{3} \in I \backslash\left\langle f_{1}, f_{2}\right\rangle, \ldots
\end{aligned}
$$

of minimal possible degree. If $d_{i}=\operatorname{deg}\left(f_{i}\right)$, say $f_{i}=a_{i} x^{d_{i}}+$ lower terms, then $d_{1} \leqslant d_{2} \leqslant d_{3} \leqslant \ldots$ and

$$
\left\langle a_{1}\right\rangle \subseteq\left\langle a_{1}, a_{2}\right\rangle \subseteq\left\langle a_{1}, a_{2}, a_{3}\right\rangle \subseteq \ldots
$$

is an ascending chain of ideals in $R$. Since $R$ is Noetherian this chain is stationary, i.e. there is some $N$ such that $\left\langle a_{1}, \ldots, a_{N}\right\rangle=\left\langle a_{1}, \ldots, a_{N+1}\right\rangle$. Hence $a_{N+1}=\sum_{i=1}^{N} b_{i} a_{i}$ for some suitable $b_{i} \in R$. Now consider

$$
\begin{aligned}
g & =f_{N+1}-\sum_{i=1}^{N} b_{i} x^{d_{N+1}-d_{i}} f_{i} \\
& =a_{N+1} x^{d_{N+1}}-\left(\sum_{i=1}^{N} b_{i} a_{i}\right) x^{d_{N+1}}+\text { lower terms. }
\end{aligned}
$$

Since $f_{N+1} \in I \backslash\left\langle f_{1}, \ldots, f_{N}\right\rangle$, it follows that $g \in I \backslash\left\langle f_{1}, \ldots, f_{N}\right\rangle$ is a polynomial of degree smaller than $d_{N+1}$, a contradiction to the choice of $f_{N+1}$.

Corollary 13.3. If $R$ is Noetherian, then $R\left[x_{1}, \ldots, x_{n}\right]$ is Noetherian. In particular, if $K$ is a field then $K\left[x_{1}, \ldots, x_{n}\right]$ is Noetherian.

Proof. Exercise (easy induction).
Corollary 13.4. If $R$ is Noetherian and $\varphi: R \rightarrow B$ is a ring homomorphism, such that $B$ is a finitely generated extension ring of $\operatorname{Im}(\varphi)$ (i.e., $B \cong R\left[x_{1}, \ldots, x_{n}\right] / I$ ), then $B$ is noetherian.
Proof. See p. 55 of [6].
Example 13.5. Similarly one can show that $K[[x]]$, the power series ring over $K$, is Noetherian.

## 14 Primary decomposition

This is sometimes also called Lasker-Noether decomposition and an analogue of decomposition of an integer into prime factors for more general rings. It also has a geometric content: we will see that the (isolated) components of a minimal primary decomposition of an ideal $I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ correspond to the irreducible components of the algebraic set $V(I) \subseteq \mathbb{A}_{K}^{n}$.

Motivation: Consider $R=\mathbb{Z}$. Then every $z \in \mathbb{Z}$ may be written as $z=p_{1}^{k_{1}} \cdots p_{n}^{k_{n}}$. One can express this in ideal notation:

$$
\langle z\rangle=\left\langle p_{1}^{k_{1}}\right\rangle \cap \cdots\left\langle p_{n}^{k_{n}}\right\rangle .
$$

Here one sees that the ideals on the right hand side are just powers of prime ideals. It is not so clear how to generalize this to Noetherian rings.

Example 14.1. Let $I=\left\langle x^{2} y, x^{2} z, x y^{2}, x z^{2}, x y z, y^{2} z, y z^{2}\right\rangle \subseteq K[x, y, z]$. Then $I$ may be written as intersection of ideals

$$
I=\langle x, y\rangle \cap\langle x, z\rangle \cap\langle y, z\rangle \cap\left\langle x, y^{2}, z^{2}\right\rangle \cap\left\langle x^{2}, y, z^{2}\right\rangle \cap\left\langle x^{2}, y^{2}, z\right\rangle .
$$

Not all of the ideals on the right hand side are powers of primes! For example, set $\mathfrak{m}=\langle x, y, z\rangle$. Then $\mathfrak{m} \supsetneq\left\langle x, y^{2}, z^{2}\right\rangle \supsetneq \mathfrak{m}^{3}$. Taking the radicals of all three ideals and noting that if $I \subseteq J$, then $\sqrt{I} \subseteq \sqrt{J}$, it follows that $\sqrt{\left\langle x, y^{2}, z^{2}\right\rangle}=\mathfrak{m}$. Since $\left\langle x, y^{2}, z^{2}\right\rangle$ is not equal to $\mathfrak{m}^{2}$, it cannot be a power of a prime ideal.

To get a bit more flexibility one makes the following
Definition 14.2. A proper ideal $\mathfrak{q} \subseteq R$ is called primary if $x y \in \mathfrak{q} \Longrightarrow$ either $x \in \mathfrak{q}$ or $y^{n} \in \mathfrak{q}$ for some $n \geqslant 1$. Equivalently, $\mathfrak{q}$ is primary if and only if $R / \mathfrak{q} \neq 0$ and every zero-divisor in $R / \mathfrak{q}$ is nilpotent.
Remark 14.3. A prime ideal is a generalisation of a prime number. In turn, a primary ideal is a generalisation of the power of a prime number. This will allow us to talk about "unique factorisation" of ideals in much the same way we do for integers or polynomials say.

Example 14.4. (i) If $I$ is prime, then $I$ is primary.
(ii) The ideal $I=\left\langle x, y^{2}, z^{2}\right\rangle$ is primary in $R=K[x, y, z]$. To see this, look at the quotient $R / I \neq \cong$ $K[y, z] /\left\langle y^{2}, z^{2}\right\rangle \neq 0$. If $\bar{f} \neq \overline{0}$ in $R / I$ is a zero-divisor, then it is easy to see that $\bar{f} \in\langle\bar{y}, \bar{z}\rangle$ and that $\bar{f}^{3}=\overline{0}$ in $R / I$.
(iii) On the other hand, if $\mathfrak{p}$ is prime, then $\mathfrak{p}^{n}$ is not necessarily primary: let $R=K[x, y, z] /\langle x y-$ $\left.z^{2}\right\rangle$. Then $I=\langle\bar{x}, \bar{z}\rangle$ is prime (since $R / I \cong K[y]$ is an integral domain). Calculate $I^{2}=$ $\left\langle\bar{x}^{2}, \overline{x z}, \bar{z}^{2}\right\rangle$. Here $\bar{z}^{2}=\overline{x y} \in I^{2}$. But neither $\bar{x}$, nor $\bar{y}$ are contained in $I=\sqrt{I}$ (direct calculation), so no power of them is in $I$. But this means that $I^{2}$ violates the condition of being a primary ideal.
(iv) $\{0\}$ and $\left\langle p^{n}\right\rangle$ for $p$ a prime, $n \geqslant 1$ are the primary ideals in $\mathbb{Z}$. These are the only ideals with prime radical, and it is then clear that they are primary.
Proposition 14.5. (1) Let $I \subseteq R$ be a primary ideal, then $\sqrt{I}$ is a prime ideal.
(2) If $\sqrt{I}=\mathfrak{m}$ is maximal, then $I$ is primary.

Proof. (1) Exercise.
(2) We show that every zero divisor in $R / I$ is nilpotent. We begin by noting from Theorem 7.6 that $\sqrt{I}$ is the intersection of all prime ideals of $R$ containing $I$. Since $\sqrt{I}$ is maximal, there is precisely one prime ideal containing $I$, namely $\mathfrak{m}$. Now by Remark 3.5 the prime ideals of $R / I$ are in correspondence with the prime ideals of $R$ containing $I$, in particular there is only one prime ideal in $R / I$ which we can write as $\operatorname{nil}(R / I)$.
Now assume that $x+I \in R / I$ is a zero divisor. Then there is some $y \notin I$ such that $x y+I=$ $0_{R / I} \in \operatorname{nil}(R / I)$. Since $\operatorname{nil}(R / I)$ is prime, we have either $x+I \in \operatorname{nil}(R / I)$ or $y+I \in \operatorname{nil}(R / I)$. Now since $\operatorname{nil}(R / I)$ is the unique prime ideal in $R / I$, and maximal ideals are also prime, we see that $R / I$ is local. So by Homework Sheet 2, Q5, we can write $\operatorname{nil}(R / I)$ as the set of non-units in $R / I$. Since $x+I$ is a zero divisor, it is not invertible and thus $x+I \in \operatorname{nil}(R / I)$ as required.

Definition 14.6. Let $R$ be a ring and let $\mathfrak{p} \subseteq R$ be a prime ideal. We say that an ideal $I \subseteq R$ is $\mathfrak{p}$-primary if $I$ is primary and $\sqrt{I}=\mathfrak{p}$. If $I$ is primary, then $\mathfrak{p}$ is called the associated prime ideal.

Theorem 14.7. Let $\mathfrak{q}_{1}, \ldots, \mathfrak{q}_{n}$ be $\mathfrak{p}$-primary ideals in $R$. Then $\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}$ is $\mathfrak{p}$-primary.
Proof. As $\sqrt{\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}}=\sqrt{\mathfrak{q}_{1}} \cap \cdots \cap \sqrt{\mathfrak{q}_{n}}=\mathfrak{p}$, we need only check that $\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}$ is primary. Assume $x, y \in R$ are such that $x y \in \mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}$. If $x \notin \mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}$ then $x \notin \mathfrak{q}_{j}$ for some $1 \leqslant j \leqslant n$. Now $x y \in \mathfrak{q}_{j}$ and since $\mathfrak{q}_{j}$ is primary we have $y^{m} \in \mathfrak{q}_{j}$ for some $m \geqslant 1$, i.e. $y \in \sqrt{\mathfrak{q}_{j}}=$ $P=\sqrt{\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}}$, and the result follows.
Definition 14.8. A primary decomposition of an ideal $I$ in a ring $R$ is an expression of $I$ as a finite intersection of primary ideals

$$
I=\bigcap_{i=1}^{n} \mathfrak{q}_{i}
$$

The decomposition is minimal (sometimes: irredundant or reduced) if:
(i) $\sqrt{\mathfrak{q}_{i}}$ are distinct for all $i$;
(ii) $\bigcap_{\substack{1 \leqslant j \leqslant n \\ j \neq i}} \mathfrak{q}_{j} \nsubseteq \mathfrak{q}_{i}$ for all $1 \leqslant i \leqslant n$.

Remark 14.9. One can always obtain a minimal primary decomposition from a given one: if $I=$ $\bigcap_{i=1}^{n} \mathfrak{q}_{i}$ is an intersection of primary ideals, then if $\mathfrak{q}_{i_{1}}, \ldots, \mathfrak{q}_{i_{k}}$ have the same associated prime $\mathfrak{p}_{i}$, we collect them together as $\mathfrak{q}_{i}^{\prime}:=\mathfrak{q}_{i_{1}} \cap \ldots \cap \mathfrak{q}_{i_{k}}$ (which is $\mathfrak{p}_{i}$-primary by Thm. 14.7). If $\bigcap_{\substack{1 \leqslant j \leqslant n \\ j \neq i}} \mathfrak{q}_{j} \subseteq \mathfrak{q}_{i}$, then omit $\mathfrak{q}_{i}$.

Theorem 14.10 (Lasker-Noether). Let $R$ be a Noetherian ring, $I \subseteq R$ an ideal. Then I has a minimal primary decomposition

$$
I=\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}
$$

Moreover, for any two minimal primary decompositions

$$
I=\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}=\mathfrak{q}_{1}^{\prime} \cap \cdots \cap \mathfrak{q}_{m}^{\prime}
$$

we have $n=m$ and (possibly after reordering) $\sqrt{\mathfrak{q}_{i}}=\sqrt{\mathfrak{q}_{i}^{\prime}}$ for all $1 \leqslant i \leqslant n$. The set $\left\{\sqrt{\mathfrak{q}_{1}}, \ldots, \sqrt{\mathfrak{q}_{n}}\right\}$ is equal to the set of prime ideals of $R$ of the form $\sqrt{(I:\langle x\rangle)}$ for some $x \in R$.
In particular, if $I=\sqrt{I}=\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}$ then the primary decomposition is unique and all $\mathfrak{q}_{i}$ are prime.
Example 14.11. (i) Let $I$ be the ideal from example 14.1: $I=\langle x, y\rangle \cap\langle x, z\rangle \cap\langle y, z\rangle \cap\left\langle x, y^{2}, z^{2}\right\rangle \cap$ $\left\langle x^{2}, y, z^{2}\right\rangle \cap\left\langle x^{2}, y^{2}, z\right\rangle$. Then we have seen this is a primary decomposition of $I$. However, this decomposistion is not minimal, since $\sqrt{\left\langle x, y^{2}, z^{2}\right\rangle}=\sqrt{\left\langle x^{2}, y, z^{2}\right\rangle}=\sqrt{\left\langle x^{2}, y^{2}, z\right\rangle}=$ $\langle x, y, z\rangle$. Use the remark above and set

$$
\mathfrak{q}^{\prime}=\left\langle x, y^{2}, z^{2}\right\rangle \cap\left\langle x^{2}, y, z^{2}\right\rangle \cap\left\langle x^{2}, y^{2}, z\right\rangle=\left\langle x^{2}, y^{2}, z^{2}, x y z\right\rangle .
$$

It is now easy to see that replacing the three ideals with $\mathfrak{q}^{\prime}$ yields a minimal primary decomposition of $I$.
(ii) Suppose $I=\langle f\rangle \subseteq K\left[x_{1}, \ldots, x_{n}\right]$, and $f=f_{1}^{n_{1}} \ldots f_{r}^{n_{r}}$ is the factorisation into irreducibles over $K$. Then $I=\left\langle f_{1}^{n_{1}}\right\rangle \cap \cdots \cap\left\langle f_{r}^{n_{r}}\right\rangle$ is a minimal primary decomposition, with associated primes $\left\{\left\langle f_{1}\right\rangle, \ldots,\left\langle f_{r}\right\rangle\right\}$.

Now we come to the proof of the primary decomposition theorem: it mainly consists of two parts - existence and uniqueness. For the existence one introduces the notion of irreducible ideals, and first shows that any ideal in a Noetherian ring can be written as an intersection of irreducible ideals, and finally that any irreducible ideal is primary.

Definition 14.12. We call an ideal $I \subseteq R$ irreducible if it cannot be written as $I_{1} \cap I_{2}$, where $I_{1}$ and $I_{2}$ are proper ideals of $R$ which strictly contain $I$.

Example 14.13. (i) $\left\langle x^{2}+1\right\rangle \subseteq \mathbb{R}[x]$ is irreducible.
(ii) $\left\langle\left(y-x^{2}\right)\left(y^{2}-x^{3}\right)\right\rangle=\left\langle y-x^{2}\right\rangle \cap\left\langle y^{2}-x^{3}\right\rangle \subseteq R[x, y]$ is reducible.

Proposition 14.14. Every proper ideal of a Noetherian ring $R$ is the intersection of finitely many irreducible ideals.

Proof. Let $S$ be the set of all ideals which are not the intersection of finitely many irreducible ideals. If $S \neq \varnothing$ then by Theorem 12.3(iii) it has a maximal element, $J$ say. Now $J$ is not irreducible, so $J=J_{1} \cap J_{2}$ for some ideals $J_{1}, J_{2} \supsetneq J$. By the maximality of $J$, it must be possible to write $J_{1}$ and $J_{2}$ as the intersection of finitely many irreducible ideals, and therefore we can also write $J$ as such. This is a contradiction, so $S=\varnothing$ and the result follows.

For the next proposition we need to recall the quotient ideal

$$
(I: J)=\{r \in R: r J \subseteq I\}
$$

for ideals $I, J \subseteq R$ from Proposition 2.5. It is an easy exercise to show that $\left(I: J_{1}+J_{2}\right)=(I:$ $\left.J_{1}\right) \cap\left(I: J_{2}\right)$ and $\left(I_{1} \cap I_{2}: J\right)=\left(I_{1}: J\right) \cap\left(I_{2}: J\right)$, which allows us to prove:

Proposition 14.15. Irreducible ideals in Noetherian rings are primary.

Proof. Let $R$ be Noetherian. We first show that if the zero ideal is irreducible then it is primary. Let $x y=0$ with $y \neq 0$ and consider the chain

$$
(0:\langle x\rangle) \subseteq(0:\langle x\rangle) \subseteq(0:\langle x\rangle) \subseteq \ldots
$$

By ACC this is stationary, i.e. $\left(0:\left\langle x^{n}\right\rangle\right)=\left(0:\left\langle x^{n+1}\right\rangle\right)=\ldots$ for some $n \geqslant 1$. It follows that $\left\langle x^{n}\right\rangle \cap\langle y\rangle=\{0\}$, for if $a \in\langle y\rangle$ then $a x=0$ so if also $a \in\left\langle x^{n}\right\rangle$ then $a=b x^{n}$ and $a x=b x^{n+1}=0$. Hence $b \in\left(0:\left\langle x^{n+1}\right\rangle\right)=\left(0:\left\langle x^{n}\right\rangle\right)$, so $b x^{n}=a=0$. Since $\{0\}$ is irreducible and $\langle y\rangle \neq 0$ we must therefore have $x^{n}=0$, i.e. $\{0\}$ is primary.
Now let $I \subseteq R$ be irreducible. Then $R / I$ is Noetherian by Theorem 12.5 and the zero ideal $\{0+$ $I\} \subseteq R / I$ is irreducible by Proposition 2.10. Therefore $\{0+I\}$ is primary, so for any $x, y \in R$ we have $x y \in I$ implies that $(x+I)(y+I) \in\{0+I\}$, thus either $x+I=0+I$ or $y^{n}+I=0+I$ for some $n$. But this is equivalent to having either $x \in I$ or $y^{n} \in I$, hence $I$ is primary.

Corollary 14.16. Every proper ideal of a Noetherian ring can be written as an intersection of finitely many primary ideals.
Proof. Exercise, use Propositions 14.14 and 14.15.
For the proof of uniqueness in the Lasker-Noether theorem and also for practical computations, one needs the following

Lemma 14.17. Let $\mathfrak{q}$ be a primary ideal in $R$. Then for any $x \in R$

$$
\sqrt{(\mathfrak{q}:\langle x\rangle)}= \begin{cases}R & \text { if } x \in \mathfrak{q} \\ \sqrt{\mathfrak{q}} & \text { if } x \notin \mathfrak{q}\end{cases}
$$

Proof. Exercise.
Proof of Thm. 14.10. Corollary 14.16 tells us that primary decompositions always exist, and now Theorem 14.7 allows us to reduce this to a minimal decomposition.
Suppose first that $\sqrt{(I:\langle x\rangle)}$ is prime for some $x \in R$. Then we have

$$
\begin{aligned}
\sqrt{(I:\langle x\rangle)} & =\sqrt{\left(\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}:\langle x\rangle\right)} \\
& =\sqrt{\left(\mathfrak{q}_{1}:\langle x\rangle\right)} \cap \cdots \cap \sqrt{\left(\mathfrak{q}_{n}:\langle x\rangle\right)}
\end{aligned}
$$

Recall from Theorem 3.9 that $I_{1} \cap \cdots \cap I_{n} \subseteq P \Longleftrightarrow I_{j} \subseteq P$ for some $j$, where $I_{i}$ are ideals and $P$ is prime. It is an easy exercise to show that in the "only if" direction, the subsets can be replaced by equalities, and hence $\sqrt{(I:\langle x\rangle)}=\sqrt{\left(\mathfrak{q}_{j}:\langle x\rangle\right)}$ for some $j$. Since $\sqrt{(I:\langle x\rangle)} \neq R$ we must have $\sqrt{(I:\langle x\rangle)}=\sqrt{\left(\mathfrak{q}_{j}:\langle x\rangle\right)}=\sqrt{\mathfrak{q}_{j}}$ by Lemma 14.17. Therefore the set of prime ideals of the form $\sqrt{(I:\langle x\rangle)}$ is a subset of $\left\{\sqrt{\mathfrak{q}_{1}}, \ldots, \sqrt{\mathfrak{q}_{n}}\right\}$.
Now consider $\sqrt{\mathfrak{q}_{i}}$. By minimality of the primary decomposition we can choose $x \in \mathfrak{q}_{j}$ for all $j \neq i$ but $x \notin \mathfrak{q}_{i}$. But then we have

$$
\begin{aligned}
\sqrt{(I:\langle x\rangle)} & =\sqrt{\left(\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}:\langle x\rangle\right)} \\
& =\sqrt{\left(\mathfrak{q}_{1}:\langle x\rangle\right)} \cap \cdots \cap \sqrt{\left(\mathfrak{q}_{n}:\langle x\rangle\right)} \\
& =\sqrt{\mathfrak{q}_{i}} .
\end{aligned}
$$

Thus $\left\{\sqrt{\mathfrak{q}_{1}}, \ldots, \sqrt{\mathfrak{q}_{n}}\right\}$ is a subset of the set of prime ideals of the form $\sqrt{(I:\langle x\rangle)}$, and the equality is established. The final statement follows immediately, since the set of primes of the form $\sqrt{(I:\langle x\rangle)}$ is independent of any choice of primary decomposition.

Definition 14.18. For any ideal $I$ of a Noetherian ring $R$, the associated primes of $I$ is the set

$$
\operatorname{Ass}(I)=\left\{\sqrt{\mathfrak{q}_{i}}: 1 \leqslant i \leqslant n, I=\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n} \text { is a minimal primary decomposition }\right\} .
$$

A minimal element in $\operatorname{Ass}(I)$ (w.r.t. inclusion) is called an isolated or minimal prime ideal. A non-isolated prime ideal is called embedded. The $\mathfrak{q}_{i}$ are called the (isolated or embedded) primary components of $I$.

If $\sqrt{I}=I=\mathfrak{q}_{1} \cap \cdots \cap \mathfrak{q}_{n}$, then the primary components are the $\sqrt{\mathfrak{q}_{i}}=\mathfrak{q}_{i}=\mathfrak{p}_{i}$ and all $\mathfrak{p}_{i}$ are isolated.

Example 14.19. An ideal $I$ is primary if and only if $\operatorname{Ass}(I)$ consists of one element.
An ideal $I$ is prime if and only if $\operatorname{Ass}(I)=I$.
Proposition 14.20. For any ideal I of a Noetherian ring $R$, the set

$$
\{x+I: x \in P \text { for some } P \in \operatorname{Ass}(I)\}
$$

is precisely the set of zero divisors of $R / I$.

## Proof. Exercise.

Example 14.21. (i) $R=\mathbb{Z}, I=\langle 12\rangle=\langle 3\rangle \cap\langle 4\rangle$. Then $\mathfrak{q}_{1}=\langle 4\rangle, \mathfrak{q}_{2}=\langle 3\rangle$ which have radicals $\langle 2\rangle$ and $\langle 3\rangle$ respectively. Therefore $\operatorname{Ass}(\langle 12\rangle)=\{\langle 2\rangle,\langle 3\rangle\}$.
(ii) Consider $I=\left\langle x, y^{2}\right\rangle \cap\langle y\rangle \subseteq K[x, y]$. Then $\mathfrak{q}_{1}=\left\langle x, y^{2}\right\rangle, \mathfrak{q}_{2}=\langle y\rangle$ have radicals $\langle x, y\rangle$ and $\langle y\rangle$ respectively, so $\operatorname{Ass}(I)=\{\langle x, y\rangle,\langle y\rangle\}$. Here $\langle y\rangle$ is an embedded component and $\langle x, y\rangle$ is an isolated component.
But $I$ also has the minimal primary decomposition $I=\langle y\rangle \cap\left\langle x^{2}, x y, y^{2}\right\rangle$ which have the same radicals as $q_{1}$ and $q_{2}$.

## 15 Noether normalisation and Hilbert's Nullstellensatz

Both of these classical theorems have a geometric background. We will only sketch this in the case of Noether normalisation, the geometric meaning of the Nullstellensatz is part of the next chapter. We will also provide proofs of both results in the next chapter (in Section 17).

For the Noether normalisation let $X=V(I) \subseteq \mathbb{A}_{K}^{n}$ be an algebraic set, where $I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ is an ideal. The normalisation theorem says that there exists a (linear) surjective and finite morphism $\pi: X \rightarrow \mathbb{A}_{K}^{d}$ onto the linear space $\mathbb{A}_{K}^{d}$. Finite is an algebraic condition and means that $K\left[x_{1}, \ldots, x_{n}\right] / I$ is a finitely generated $K\left[x_{1}, \ldots, x_{d}\right]$-module under the map $\pi^{*}: K\left[x_{1}, \ldots, x_{d}\right] \rightarrow$ $K\left[x_{1}, \ldots, x_{n}\right] / I, f \mapsto \pi^{*}(f)=f \circ \pi$. In particular, if $\pi$ is finite, then it has finite fibers, that is, for any $b \in \mathbb{A}_{K}^{d}$ the set $\pi^{-1}(b)$ consists of a finite number of points.

Example 15.1. (i) Let $X=V\left(y-x^{2}\right) \subseteq \mathbb{A}_{\mathbb{R}}^{2}$. We can project $X$ onto each of the two coordinate axes: $\pi_{x}: X \rightarrow \mathbb{A}_{\mathbb{R}}^{1}:(x, y) \mapsto x$ and $\pi_{y}: X \rightarrow \mathbb{A}_{\mathbb{R}}^{1}:(x, y) \mapsto y$. The first projection $\pi_{x}$ is even bijective, for $\pi_{y}$ the fibers $\pi_{y}^{-1}(b), b \in \mathbb{A}_{\mathbb{R}}^{1}$, consist of either 1 or 2 points.
Algebraically for $\pi_{x}^{*}$ we have $\pi_{x}^{*}: \mathbb{R}[x] \rightarrow \mathbb{R}[x, y] /\left\langle y-x^{2}\right\rangle \cong \mathbb{R}\left[x, x^{2}\right]$. Clearly, $\mathbb{R}\left[x, x^{2}\right]=$ $\mathbb{R}[x]$ is finitely generated as an $\mathbb{R}[x]$-module here!
(ii) Consider the cross $V(x y) \subseteq \mathbb{A}_{\mathbb{R}}^{2}$ and take again the projections $\pi_{x}$ and $\pi_{y}$ onto the two coordinate axes. Here neither of the two projections is finite, since $\pi_{x}^{-1}(0)$ is the whole $y$ axis, and $\pi_{y}^{-1}(0)$ is the $x$-axis. Algebraically, one sees for example that for $\pi_{x}^{*}: \mathbb{R}[x] \rightarrow$ $\mathbb{R}[x, y] /\langle x y\rangle$ the module $\mathbb{R}[x, y] /\langle x y\rangle$ is not finitely generated over $\mathbb{R}[x]$ : it is the infinite direct sum $\mathbb{R}[x] \oplus y \mathbb{R}[x] \oplus y^{2} \mathbb{R}[x] \oplus \cdots$.
In the second example above, the (proof of the) Noether normalisation theorem will tell us how to modify $X$ to obtain a finite projection onto a linear space. For this first recall the following

Definition 15.2. Let $R$ be a ring. An $R$-algebra is a ring $S$ with a ring homomorphism $\varphi: R \rightarrow S$. We say $S$ is a finite $R$-algebra if it is finitely generated as an $R$-module, i.e. there exist $x_{1}, \ldots, x_{n} \in S$ such that

$$
S=R x_{1}+\cdots+R x_{n}
$$

If also $R$ is a field then we say $S$ is a finite dimensional $R$-algebra.
We say $S$ is a finitely generated $R$-algebra if there exist $x_{1}, \ldots, x_{n} \in S$ such that $S=R\left[x_{1}, \ldots, x_{n}\right]$.
Example 15.3. (i) $R[x]$ is an $R$-algebra via the natural inclusion map. It is not finite, but it is finitely generated.
(ii) $\mathbb{Q}[\sqrt{2}]$ is finitely generated over $\mathbb{Q}$ and also finite, since $\mathbb{Q}[\sqrt{2}]=\mathbb{Q} \oplus \mathbb{Q} \sqrt{2}$ as $\mathbb{Q}$-vector space.
(iii) $K[t]$ is a finitely generated $R=K\left[t^{2}, t^{3}\right]$-algebra: $K[t]=R[t]$ as algebras and $K[t]=R+R t$ as $R$-module.
(iv) Any finitely generated $K$-algebra is of the form $K\left[x_{1}, \ldots, x_{n}\right] / I$, where $I$ is an ideal in $K\left[x_{1}, \ldots, x_{n}\right]$ : Let $S=K\left[a_{1}, \ldots, a_{n}\right]$ be a finitely generated $K$-algebra, with $a_{i} \in S$. We have an algebra homomorphism (this is a ring homomorphism that is also a $K$-module homomorphism) $\varphi: K\left[x_{1}, \ldots, x_{n}\right] \rightarrow S, x_{i} \mapsto a_{i}$. Then by construction $\varphi$ is surjective, and by the homomorphism theorem $S \cong K\left[x_{1}, \ldots, x_{n}\right] / \operatorname{Ker}(\varphi)$.

The homomorphism $\varphi$ turns $S$ into an $R$-module, where multiplication is defined by $r \cdot s=\varphi(r) s$ for all $r \in R, s \in S$.
When $R \subseteq S$, we call $S$ an extension ring of $R$. If in addition $R$ and $S$ are fields, then we call $S$ an extension field of $R$.

Definition 15.4. Let $S$ be an $R$-algebra. An element $s \in S$ is integral over $R$ if there exists a monic polynomial

$$
f(x)=x^{n}+a_{n-1} x^{n-1}+a_{n-2} x^{n-2}+\cdots+a_{0} \in R[x]
$$

such that $f(s)=0$.
We say $S$ is integral over $R$ if every $s \in S$ is integral over $R$. If also $R \subseteq S$, then we call $S$ an integral extension.

Example 15.5. (i) The integral elements of $\mathbb{Q}$ over $\mathbb{Z}$ are the integers.
(ii) $K\left[x^{2}\right] \subseteq K[x]$ is an integral extension.

The following result will be crucial in the proof of Hilbert's Nullstellensatz. For a proof see Section 17.

Theorem 15.6 (Noether Normalisation). Let $K$ be an infinite field and $S$ a finitely generated $K$-algebra. Then there exist $z_{1}, \ldots, z_{m} \in S$ such that:
(i) $z_{1}, \ldots, z_{m}$ are algebraically independent over $K$, i.e. there is no non-zero polynomial $f \in K\left[x_{1}, \ldots, x_{m}\right]$ such that $f\left(z_{1}, \ldots, z_{m}\right)=0$;
(ii) $S$ is finite over $R=K\left[z_{1}, \ldots, z_{m}\right]$.

Remark. (i) In fact Theorem 15.6 does hold for finite fields, but an alternative proof is needed (for instance, see [6] or [1]). In the following we will assume the normalisation theorem for any field.
(ii) Theorem 15.6 shows that any finitely generated extension $K \subseteq S$ can be written as a composite

$$
K \subseteq K\left[z_{1}, \ldots, z_{m}\right] \subseteq S
$$

where the first extension is polynomial and the second is finite.

Theorem 15.7 (Weak Nullstellensatz). Let $K$ be a field and $S$ a finitely generated $K$-algebra. If $S$ is also a field, then $S$ is finitely generated as a K-module.
In particular, if $K$ is algebraically closed then every maximal ideal of $K\left[x_{1}, \ldots, x_{n}\right]$ is of the form $\left\langle x_{1}-\right.$ $\left.a_{1}, \ldots, x_{n}-a_{n}\right\rangle$ for some $a_{1}, \ldots, a_{n} \in K$.

The proof of this theorem is also deferred to Section 17.

## Part II

## Algebraic Geometry

## 16 The algebra-geometry dictionary

Let $K$ be a field (we will usually assume it to be algebraically closed) and consider the polynomial ring $K\left[x_{1}, \ldots, x_{n}\right]$. Normally we deal with polynomials simply as elements in the ring, but we will now consider them as maps from $K^{n}$ to $K$ by substituting the variables $x_{1}, \ldots, x_{n}$ with elements of $K$.

Definition 16.1. Let $S$ be a subset of $K\left[x_{1}, \ldots, x_{n}\right]$. The vanishing locus of $S$ is the set

$$
\mathbb{V}(S)=\left\{\left(a_{1}, \ldots, a_{n}\right) \in K^{n}: f\left(a_{1}, \ldots, a_{n}\right)=0 \text { for all } f \in S\right\} .
$$

A set $X \subseteq K^{n}$ is called an algebraic set or algebraic variety if $X=\mathbb{V}(S)$ for some such $S$. The set $K^{n}$ is often denoted $\mathbb{A}_{K}^{n}$ and is called affine $n$-space (this is done to avoid giving 0 special status). If $I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ is and ideal, then $\mathbb{V}(I)$ is called the vanishing set of $I$.

Remark 16.2. If $I=\left\langle f_{1}, \ldots, f_{m}\right\rangle$, then $\mathbb{V}\left(f_{1}, \ldots, f_{m}\right)=\mathbb{V}(I)$ and every algebraic set is of the form $\mathbb{V}(I)$ for some ideal $I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ (see this with Hilbert's basis theorem!).

Example 16.3. (i) $\mathbb{V}\left(x^{2}+y^{2}-1\right) \subseteq \mathbb{A}_{\mathbb{R}}^{2}$ is a circle.
(ii) $\mathbb{V}(x y z) \subseteq \mathbb{A}_{\mathbb{R}}^{3}$ is the union of the three planes $\{x=0\},\{y=0\}$ and $\{z=0\}$, see Fig. II.1.


Figure II.1: Union of the three coordinate planes $\mathbb{V}(x y z)$
(iii) $\mathbb{V}\left(\left(x^{2}+y^{2}\right)^{3}-4 x^{2} y^{2}\right)$ is the four leaf clover:

(iv) $\mathbb{V}\left(z-x^{3}, y-x^{2}\right)$ is a curve in $\mathbb{A}_{\mathbb{R}}^{3}$ and is a twisted cubic $t \mapsto\left(t, t^{2}, t^{3}\right)$ :

(v) $\mathbb{V}\left(y^{2}-x^{3}-a x-b\right) \subseteq \mathbb{A}_{\mathbb{C}}^{2}$ gives an elliptic curve. These are very important in many branches of mathematics.
(vi) The surface $\mathbb{V}\left(z^{2}+x\left(y^{2}-x^{2}\right)\right) \subseteq \mathbb{A}_{\mathbb{R}}^{3}$ looks like three cones meeting at a point, see Fig. II.2. This surface is a so-called $D_{4}$-singularity and example of an ADE-surface singularity. For more visualizations of these surfaces see http://www1.maths.leeds.ac.uk/~pmtemf/web/ gallery-ADE.html.
(vii) $\mathbb{V}\left(16 x^{4} z-4 x^{3} y^{2}-128 x^{2} z^{2}+144 x y^{2} z-27 y^{4}+256 z^{3}\right) \subseteq \mathbb{A}_{\mathbb{R}}^{3}$ is the so-called swallowtail, see Fig. II.3. This surface appears in many contexts, e.g. as the discriminant of a quartic polynomial, see also https://imaginary.org/sites/default/files/snapshots/snapshot-2014-007. pdf
(viii) $\mathbb{V}\left(x z-y^{2}, x^{3}-y z, z^{2}-x^{2} y\right) \subseteq \mathbb{A}_{\mathbb{R}}^{3}$ gives a singular twisted cubic $t \mapsto\left(t^{3}, t^{4}, t^{5}\right)$. Note that this has codimension 2 , but has 3 generators. In fact this this set cannot be 2-generated.
(ix) If $a_{1}, \ldots, a_{n} \in K$, then $\mathbb{V}\left(x_{1}-a_{1}, \ldots, x_{n}-a_{n}\right) \subseteq \mathbb{A}_{K}^{n}$ is the point $\left(a_{1}, \ldots, a_{n}\right)$.


Figure II.2: The surface $\mathbb{V}\left(z^{2}+x\left(y^{2}-x^{2}\right)\right)$ in $\mathbb{R}^{3}$ (the highlighted curve is the intersection of the surface with the plane $\mathbb{V}(z))$.


Figure II.3: The swallowtail in $\mathbb{R}^{3}$.
(x) Spirals $r=\cos \theta$ are not algebraic sets, as they give a polynomial with an infinite number of zeros.

Remark 16.4. More visualizations of algebraic surfaces can be found on the Imaginary portal https://imaginary.org/galleries.

Some properties of algebraic sets:

$$
\mathbb{V}\left(f_{1}, \ldots, f_{r}\right)=\bigcap_{i=1}^{r} \mathbb{V}\left(f_{i}\right)
$$

so every algebraic set in the intersection of a finite number of hypersurfaces, algebraic sets generated by a single non-zero polynomial. In particular, the algebraic subsets of $A_{K}^{1}$ are just the finite subsets plus all of $K$ (as $\mathbb{V}(\{0\})=K)$.
Now we get a functor $\mathbb{V}$ :

$$
\mathbb{V}: \text { Ideals in } K\left[x_{1}, \ldots, x_{n}\right] \longrightarrow \text { Algebraic sets in } \mathbb{A}_{K}^{n}
$$

Proposition 16.5. Let $R=K\left[x_{1}, \ldots, x_{n}\right]$. Then:
(i) $\mathbb{V}(\{0\})=\mathbb{A}_{K}^{n}$ and $\mathbb{V}(R)=\varnothing$;
(ii) $I \subseteq J \Longrightarrow \mathbb{V}(I) \supseteq \mathbb{V}(J)$ for ideals $I$, $J$ of $R$;
(iii) $\mathbb{V}(I J)=\mathbb{V}(I \cap J)=\mathbb{V}(I) \cup \mathbb{V}(J)$ for ideals $I$, $J$ of $R$;
(iv) for any set $\left\{I_{\lambda}\right\}_{\lambda \in \Lambda}$ of ideals of $R$,

$$
\mathbb{V}\left(\sum_{\lambda \in \Lambda} I_{\lambda}\right)=\bigcap_{\lambda \in \Lambda} \mathbb{V}\left(I_{\lambda}\right)
$$

Proof. (i) Exercise.
(ii) Exercise.
(iii) Since $I$ and $J$ both contain $I \cap J$, which in turn contains $I J$, we see from (ii) that

$$
\mathbb{V}(I) \cup \mathbb{V}(J) \subseteq \mathbb{V}(I \cap J) \subseteq \mathbb{V}(I J)
$$

Now if $x \notin \mathbb{V}(I) \cup \mathbb{V}(J)$ then there exists $f \in I$ and $g \in J$ such that $f(x) \neq 0$ and $g(x) \neq 0$. Hence $(f g)(x) \neq 0$ so $x \notin \mathbb{V}(I J)$. Thus $\mathbb{V}(I J) \subseteq \mathbb{V}(I) \cup \mathbb{V}(J)$.
(iv) Since $I_{\mu} \subseteq \sum_{\lambda \in \Lambda} I_{\lambda}$ for all $\mu \in \Lambda$, we have from (ii) that

$$
\mathbb{V}\left(\sum_{\lambda \in \Lambda} I_{\lambda}\right) \subseteq \bigcap_{\lambda \in \Lambda} \mathbb{V}\left(I_{\lambda}\right)
$$

Now if $x \in \bigcap_{\lambda \in \Lambda} \mathbb{V}\left(I_{\lambda}\right)$ and $f \in \sum_{\lambda \in \Lambda}$ then $f=\sum_{i=1}^{m} f_{\lambda_{i}}$ for some $m \in \mathbb{N}, \lambda_{i} \in \Lambda$ and $f_{\lambda_{i}} \in I_{\lambda_{i}}$. Then we have $f(x)=\sum_{i=1}^{m} f_{\lambda_{i}}(x)=0$.

Example 16.6. Consider $\mathbb{A}_{\mathbb{R}}^{3}$. Then

$$
\mathbb{V}(x z, y z)=\mathbb{V}(\langle z\rangle \cap\langle x, y\rangle)=\mathbb{V}(z) \cup \mathbb{V}(x, y)
$$

is the union of the $(x, y)$-plane and the $z$-axis.
Remark 16.7. Proposition 16.5 can be used to show that the sets $\mathbb{V}(S)$ for $S \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ define the closed sets for a topology on $\mathbb{A}_{K}^{n}$. We call this topology the Zariski topology. Facts from commutative algebra, e.g. Hilbert's Basis Theorem, properties of Noetherian rings etc., can be used to prove statements about the Zariski topology, for instance any closed subset of $\mathbb{A}_{K}^{n}$ is compact. More generally one can also define the Zariski topology for any commutative ring $R$ : the closed sets are then of the form $\mathbb{V}(I)$ for any ideal $I \subseteq R$ and are defined as

$$
\mathbb{V}(I)=\{\mathfrak{p} \in \operatorname{Spec}(R): I \subseteq \mathfrak{p}\}
$$

We now introduce an "inverse" to $\mathbb{V}$ :

$$
\left.\mathbb{I}: \text { Subsets of } \mathbb{A}_{K}^{n}\right\} \longrightarrow \text { Ideals in } K\left[x_{1}, \ldots, x_{n}\right] .
$$

Definition 16.8. For a subset $X \subseteq \mathbb{A}_{K}^{n}$ let the vanishing ideal of $X$ be the set

$$
\mathbb{I}(X)=\left\{f \in K\left[x_{1}, \ldots, x_{n}\right]: f(x)=0 \text { for all } x \in X\right\} .
$$

That this is an ideal is clear.
However, in general one does not have $\mathbb{I}(\mathbb{V}(J))=J$ for an ideal $J \subseteq K\left[x_{1}, \ldots, x_{n}\right]$.
Example 16.9. (i) Let $J=\langle x, y\rangle^{2}$ in $K[x, y]$. Clearly, $\mathbb{V}(J)=\{(0,0)\}$, but $\mathbb{I}(\{(0,0)\})=\langle x, y\rangle$.
(ii) Let $X$ be the cusp $\mathbb{V}\left(x^{3}-y^{2}\right)$ in $\mathbb{A}_{K}^{2}$. In this case we have $\mathbb{I}(X)=\left\langle x^{3}-y^{2}\right\rangle$.
(iii) Let $X=\left\{n \in \mathbb{Z} \subseteq \mathbb{A}_{\mathbb{R}}^{1}\right\}$. This set is not algebraic! But we still can find $\mathbb{I}(X)$ :

$$
\mathbb{I}(X)=\{f \in \mathbb{R}[x]: f(x)=0 \text { for all } x \in \mathbb{N}\}=\langle 0\rangle
$$

Proposition 16.10. (i) $\mathbb{I}(\varnothing)=K\left[x_{1}, \ldots, x_{n}\right]$. If $K$ is infinite then $\mathbb{I}\left(\mathbb{A}_{K}^{n}\right)=\{0\}$.
(ii) $X \subseteq Y \subseteq \mathbb{A}_{K}^{n} \Longrightarrow \mathbb{I}(X) \supseteq \mathbb{I}(Y)$.
(iii) $X, Y \subseteq \mathbb{A}_{K}^{n} \Longrightarrow \mathbb{I}(X \cup Y)=\mathbb{I}(X) \cap \mathbb{I}(Y)$.

Proof. (i) The first part is clear. The second follows from Lemma 17.3.
(ii) Straightforward: if $X \subseteq Y$ then we need more functions to define it.
(iii) We have

$$
\begin{aligned}
f \in \mathbb{I}(X \cup Y) & \Longleftrightarrow f(x)=0 \text { for all } x \in X \cup Y \\
& \Longleftrightarrow f(x)=0 \text { for all } x \in X \text { and for all } x \in Y \\
& \Longleftrightarrow f \in \mathbb{I}(X) \cap \mathbb{I}(Y) .
\end{aligned}
$$

Remark. Note that in (i) the assumption that $K$ is infinite is necessary. For instance if $p \in \mathbb{Z}$ is prime, $K=\mathbb{Z}_{p}$ and $f(x)=x^{p}-x \in K[x]$, then $f \in \mathbb{I}\left(\mathbb{A}_{K}^{1}\right)$.
Proposition 16.11. Let I be an ideal of $K\left[x_{1}, \ldots, x_{n}\right]$ and $X$ a subset of $\mathbb{A}_{K}^{n}$. Then:
(i) $X \subseteq \mathbb{V}(\mathbb{I}(X))$, with equality if and only if $X$ is an algebraic set;
(ii) $I \subseteq \mathbb{I}(\mathbb{V}(I))$.

Proof. The two inclusions are mostly tautological, for instance if $\mathbb{I}(X)$ is defined to be the set of functions vanishing on $X$ then for any point $x \in X$ all functions in $\mathbb{I}(X)$ vanish on it.
If $X=\mathbb{V}(\mathbb{I}(X))$ then $X$ is algebraic as it is of the form $X=\mathbb{V}(J)$ for some ideal $J$. Conversely if $X$ is algebraic then $X=\mathbb{V}(J)$ for some ideal $J$. But $J \subseteq \mathbb{I}(X)$ so $\mathbb{V}(\mathbb{I}(X)) \subseteq \mathbb{V}(J)=X$.

We would like a condition to ensure equality in Proposition 16.11(ii). This is not so easy, as two types of problems can occur:
(i) $\left\langle x^{n}\right\rangle \subsetneq \mathbb{I}\left(\mathbb{V}\left(x^{n}\right)\right)=\langle x\rangle$ for all $n \geqslant 2$, so non-reduced elements present a challenge, and
(ii) in $\mathbb{R}[x],\left\langle x^{2}+1\right\rangle \subsetneq \mathbb{I}\left(\mathbb{V}\left(x^{2}+1\right)\right)=\mathbb{I}(\varnothing)=\mathbb{R}[x]$, so the ideal may product no zeroes.

We can attempt to solve (i) by using the radical $\sqrt{I}$, but even this is not enough. In fact, (ii) gives a strict inclusion here too as $\sqrt{\left\langle x^{2}+1\right\rangle}=\left\langle x^{2}+1\right\rangle$. The correct way to fix this is using Hilbert's Nullstellensatz.
Recall that a field $K$ is called algebraically closed if every non-constant polynomial in $K[x]$ has a root in K.
Theorem 16.12 (Nullstellensatz). Let $K$ be an algebraically closed field and $I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ an ideal. Then:
(Weak) $I \neq K\left[x_{1}, \ldots, x_{n}\right] \Longrightarrow \mathbb{V}(I) \neq \varnothing$;
(Strong) $\mathbb{I}(\mathbb{V}(I))=\sqrt{I}$.
We will prove this Theorem in Section 17. This theorem says that we have correspondences


So it is not clear yet to which algebraic subsets the prime ideals correspond. This will be tightly connected with the geometric interpretation of primary decomposition.

Example 16.13. (i) Let $J=\left\langle x^{2}-y^{2}\right\rangle$ in $K[x, y]$. $J$ is not prime since $(x+y)(x-y) \in J$ but none of the two factors is an element of $J$. We have already seen that $\mathbb{V}(J)=\mathbb{V}(x+y) \cup \mathbb{V}(x-y)$ is a union of two hyperplanes.
(ii) Let $J=\left\langle x^{2} y, x^{2} z, y^{2} x, y^{2} z, z^{2} x, z^{2} y, x y z\right\rangle$ be an ideal in $K[x, y, z]$. A minimal primary decomposition of $J$ is $J=\langle x, y\rangle \cap\langle x, z\rangle \cap\langle y, z\rangle \cap\left\langle x^{2}, y^{2}, z^{2}, x y z\right\rangle$. Here we see that $\sqrt{J}=$ $\langle x, y\rangle \cap\langle x, z\rangle \cap\langle y, z\rangle$ and

$$
\mathbb{V}(J)=\mathbb{V}(\sqrt{J})=\mathbb{V}(x, y) \cup \mathbb{V}(x, z) \cup \mathbb{V}(y, z)
$$

is the union of the three coordinate axes.
Lemma 16.14. Every non-empty set of algebraic subsets of $\mathbb{A}_{K}^{n}$ has a minimal element.
Proof. (This was an exercise in the lecture!) Suppose that $\Sigma$ is a non-empty set of algebraic subsets of $\mathbb{A}_{K}^{n}$ with no minimal element. Then we can find a strictly descending chain $X_{1} \supsetneq X_{2} \supsetneq X_{3} \supsetneq$ $\ldots$. Recall that $X_{1} \supseteq X_{2} \Longrightarrow \mathbb{I}\left(X_{1}\right) \subseteq \mathbb{I}\left(X_{2}\right)$, and note that if the left subset is strict then so too is the right. Therefore we have a strictly ascending chain of ideals

$$
\mathbb{I}\left(X_{1}\right) \subsetneq \mathbb{I}\left(X_{2}\right) \subsetneq \mathbb{I}\left(X_{3}\right) \subsetneq \ldots,
$$

in $K\left[x_{1}, \ldots, x_{n}\right]$. But $K\left[x_{1}, \ldots, x_{n}\right]$ is Noetherian, so this is a contradiction.
Definition 16.15. An algebraic set $X \subseteq \mathbb{A}_{K}^{n}$ is called irreducible if for all decompositions $X=$ $X_{1} \cup X_{2}$ with $X_{1}, X_{2} \subseteq X$ algebraic sets, we have either $X=X_{1}$ or $X=X_{2}$. Sometimes in the literature an Irreducible algebraic set is called algebraic variety. (However, we use the term algebraic variety for any algebraic set here!)

Example 16.16. $\mathbb{V}(x y) \subseteq \mathbb{A}_{\mathbb{R}}^{2}$ is the two coordinate axes which can be written as the union $\mathbb{V}(x) \cup$ $\mathbb{V}(y)$, so is reducible.

Proposition 16.17. (i) Let $X \subseteq \mathbb{A}_{K}^{n}$ be an algebraic set and $\mathbb{I}(X)$ the vanishing ideal of $X$. Then $X$ is irreducible if and only if $\mathbb{I}(X)$ is prime.
(ii) Any algebraic set has an expression

$$
X=X_{1} \cup \cdots \cup X_{r}
$$

unique up to reordering of the $X_{i}$, with $X_{i}$ irreducible and $X_{i} \nsubseteq X_{j}$ for $i \neq j$. The $X_{i}$ are called the irreducible components of $X$.

Proof. (i) We prove that $X$ is reducible if and only if $\mathbb{I}(X)$ is not prime. Indeed, suppose $X=$ $X_{1} \cup X_{2}$ is a non-trivial decomposition of $X$ into algebraic sets. Then $X_{1}, X_{2} \subsetneq X$ means that there is some $f_{1} \in \mathbb{I}\left(X_{1}\right) \backslash \mathbb{I}(X)$ and some $f_{2} \in \mathbb{I}\left(X_{2}\right) \backslash \mathbb{I}(X)$. The product $f_{1} f_{2}$ vanishes at all points of $X$, so $f_{1} f_{2} \in \mathbb{I}(X)$. Therefore $\mathbb{I}(X)$ is not prime.
Conversely, suppose that $\mathbb{I}(X)$ is not prime. Then there exists $f_{1}, f_{2} \notin \mathbb{I}(X)$ such that $f_{1} f_{2} \in$ $\mathbb{I}(X)$. Let $X_{1}=\mathbb{V}\left(\mathbb{I}(X)+\left\langle f_{1}\right\rangle\right)$ and $X_{2}=\mathbb{V}\left(\mathbb{I}(X)+\left\langle f_{2}\right\rangle\right)$. Then by Proposition 16.5

$$
\begin{aligned}
X_{1} & =\mathbb{V}(\mathbb{I}(X)) \cap \mathbb{V}\left(f_{1}\right) \\
& =X \cap \mathbb{V}\left(f_{1}\right) \text { since } X \text { is an algebraic set } \\
& \subsetneq X \text { since } f_{1} \notin \mathbb{I}(X),
\end{aligned}
$$

similarly $X_{2} \subsetneq X$, and both are algebraic sets. So $X_{1} \cup X_{2} \subseteq X$, and moreover

$$
\begin{aligned}
\left(\mathbb{I}(X)+\left\langle f_{1}\right\rangle\right)\left(\mathbb{I}(X)+\left\langle f_{2}\right\rangle\right) & =\mathbb{I}(X)^{2}+\left\langle f_{1}\right\rangle \mathbb{I}(X)+\left\langle f_{2}\right\rangle \mathbb{I}(X)+\left\langle f_{1} f_{2}\right\rangle \\
& \subseteq \mathbb{I}(X),
\end{aligned}
$$

so by Propositions 16.5 and 16.11 we have $X=\mathbb{V}(\mathbb{I}(X)) \subseteq \mathbb{V}\left(\left(\mathbb{I}(X)+\left\langle f_{1}\right\rangle\right)\left(\mathbb{I}(X)+\left\langle f_{2}\right\rangle\right)\right)=$ $X_{1} \cup X_{2}$. Thus $X=X_{1} \cup X_{2}$, but neither component is equal to $X$, so $X$ is reducible.
(ii) Let $\Sigma$ be the set of algebraic subsets of $\mathbb{A}_{K}^{n}$ which do not have such a decomposition. If $\Sigma=\varnothing$ then we are done, otherwise by Lemma 16.14 there is a minimal element $X \in \Sigma$. If $X$ is irreducible, then $X \notin \Sigma$, a contradiction. Otherwise $X$ has a non-trivial decomposition $X=X_{1} \cup X_{2}$, and the minimality of $X$ shows that $X_{1}, X_{2} \notin \Sigma$ and so have a decomposition into irreducibles. But then putting these decompositions together gives a decomposition of $X$, so $X \notin \Sigma$, another contradiction. Therefore $\Sigma=\varnothing$ and the existence is proved.
Uniqueness is left as an exercise.

Remark. The decomposition of $X$ into irreducibles $X_{i}$ corresponds to a minimal primary decomposition of $\mathbb{I}(X)$. The associated primes in the latter case are the prime ideals $\mathbb{I}\left(X_{i}\right)$.

Example 16.18. We will decompose the algebraic set $X=\mathbb{V}\left(x^{2}-y z, x z-x\right) \subseteq \mathbb{A}_{K}^{3}$ into its irreducible components, assuming that the field $K$ is infinite. We begin by considering ( $p_{1}, p_{2}, p_{3}$ ) $\in$ $X$, and note that if $p_{1}=0$ then we must also have $p_{2} p_{3}=0$ (so either $p_{2}=0$ or $p_{3}=0$ ). This part corresponds to the algebraic set $\mathbb{V}(x, y z)=\mathbb{V}(x, y) \cup \mathbb{V}(x, z)$.
If now $p_{1} \neq 0$ then $p_{1} p_{3}-p_{1}=0 \Longrightarrow p_{3}=1$, and thus $p_{1}^{2}=p_{2}$. Therefore this part corresponds to the algebraic set $\mathbb{V}\left(x^{2}-y, z-1\right)$, and we can decompose

$$
X=\mathbb{V}(x, y) \cup \mathbb{V}(x, z) \cup \mathbb{V}\left(x^{2}-y, z-1\right)
$$

We will prove using Proposition 16.17 that each of the three components is irreducible. By the strong Nullstellensatz we have $\mathbb{I}(\mathbb{V}(x, y))=\sqrt{\langle x, y\rangle}$, but $\langle x, y\rangle$ is prime as $K[x, y, z] /\langle x, y\rangle \cong K[z]$ so $\sqrt{\langle x, y\rangle}=\langle x, y\rangle$. Thus by Proposition 16.17(i) we see that $\mathbb{V}(x, y)$ is irreducible. Similarly $\mathbb{V}(x, z)$ is irreducible. Finally, $K[x, y, z] /\left\langle x^{2}-y, z-1\right\rangle \cong K[x]$ so $\left\langle x^{2}-y, z-1\right\rangle$ is also prime so this component is also irreducible.
To sum up, we obtain the following dictionary between algebra and geometry, cf. [2, Ch. 4,§8]:

| Algebra |  | Geometry |
| :---: | :---: | :---: |
| radical ideal |  | algebraic variety |
| $J$ | $\longrightarrow$ | $\mathbb{V}(J)$ |
| $\mathbb{I}(X)$ | $\longleftarrow$ | X |
| sum of ideals |  | intersection of varieties |
| $I+J$ | $\longrightarrow$ | $\mathbb{V}(I) \cap \mathbb{V}(J)$ |
| $\sqrt{\mathbb{I}(X)+\mathbb{I}(Y)}$ | $\longleftarrow$ | $X \cap Y$ |
| intersection (multiplication) of ideals |  | union of varieties |
| $I \cap J(I \cdot J)$ | $\longrightarrow$ | $\mathbb{V}(I) \cup \mathbb{V}(J)$ |
| $\mathbb{I}(X) \cap \mathbb{I}(Y)(\sqrt{\mathbb{I}(X) \cdot \mathbb{I}(Y)})$ | $\stackrel{ }{ }$ | $X \cup Y$ |
| minimal primary decomposition $I=\sqrt{I}=\mathfrak{p}_{1} \cap \cdots \cap \mathfrak{p}_{m}$ | $\longrightarrow$ | decomposition into irreducible components $\mathbb{V}(I)=\mathbb{V}\left(\mathfrak{p}_{1}\right) \cup \cdots \cup \mathbb{V}\left(\mathfrak{p}_{m}\right)$ |
| $\mathbb{I}(X)=\bigcap_{i=1}^{m} \mathbb{I}\left(X_{i}\right)$ <br> prime ideal | $\longleftarrow$ | $X=\bigcup_{i=1}^{m} X_{i}$ <br> irreducible variety |
| maximal ideal |  | point in $\mathbb{A}_{K}^{n}$ (where $K$ alg. closed) |
| ascending chain condition |  | descending chain condition |

From here some natural further questions about the geometry of algebraic varieties arise:
(i) Most basic here is the question whether $X \subseteq \mathbb{A}_{K}^{n} \neq \varnothing$. If $X=\mathbb{V}\left(f_{1}, \ldots, f_{m}\right) \subseteq \mathbb{A}_{K}^{n}(K$ algebraically closed), then by Hilbert's Nullstellensatz $X=\varnothing$ if and only if $1 \in\left\langle f_{1}, \ldots, f_{m}\right\rangle$. In order to solve the geometric problem, we thus have to solve the ideal membership problem, see Section 18 on Gröbner bases!
(ii) Determine the irreducible components of $X$ : this can be done, as soon as we can compute a minimal primary decomposition of $\mathbb{I}(X)$. Therefor one also uses Gröbner bases, but the discussion of the algorithms is beyond the scope of this course. See [4, Chapter 4] for a general discussion and [5] for the special case of monomial ideals.
(iii) Determine the intersection behavior of the "smooth" irreducible components of an algebraic variety $X$. This leads to the field of intersection theory.
(iv) Study of singular points: these are the points on an algebraic variety $X$, where $X$ is not "smooth". This leads to classsification problems and the problem of resolution of singularities, both active research areas.

## 17 The proofs of the Noether Normalisation lemma and Hilbert's Nullstellensatz

In this section we prove the Normalisation lemma (Theorem 15.6) and the algebraic version of Hilbert's Nullstellensatz (Theorem 15.7). From this we will obtain a proof of the weak and strong
geometric version of the Nullstellensatz (Theorem 16.12).
First we need a few facts about finite algebras and integral elements.
Proposition 17.1. (i) Let $R \subseteq S \subseteq T$ be rings. If $S$ is a finite $R$-algebra and $T$ is a finite $S$-algebra, then $T$ is a finite $R$-algebra.
(ii) If $R \subseteq S$ is a finite $R$-algebra and $t \in S$, then $t$ satisfies a monic polynomial over $R$.
(iii) If $S$ is an $R$-algebra and $t \in S$ is integral over $R$, then $R[t]$ is a finite $R$-algebra.

## Proof. (i) Exercise.

(ii) Suppose $S=\sum_{i=1}^{n} R s_{i}$. Then for each $i, t s_{i} \in S$ so there exist $r_{i j} \in R$ such that

$$
t s_{i}=\sum_{j=1}^{n} r_{i j} s_{j} \Longrightarrow \sum_{j=1}^{n}\left(t \delta_{i j}-r_{i j}\right) s_{j}=0
$$

where $\delta_{i j}$ is the Kronecker Delta, taking value 1 if $i=j$ and 0 otherwise. Now let $A$ be the matrix with $A_{i j}=t \delta_{i j}-r_{i j}$, and set $\Delta=\operatorname{det} A$ and $\underline{s}=\left(s_{1}, \ldots, s_{n}\right)^{\top}$. Then $A \underline{s}=0$, hence $0=\left(A^{\text {adj }}\right) A \underline{s}=\Delta \underline{s}$ where $A^{\text {adj }}$ is the adjoint matrix. Therefore $\Delta s_{i}=0$ for all $i$. But $1 \in S$ is a linear combination of the $s_{i}$, so in particular we have $\Delta=\Delta \cdot 1=$. Therefore the monic polynomial $\operatorname{det}\left(x \delta_{i j}-r_{i j}\right)$ over $R$ is satisfied by $t$.
(iii) Exercise.

Corollary 17.2. Let $S$ be a field and $R$ a subring of $S$ such that $S$ is a finite $R$-algebra. Then $R$ is a field.
Proof. For any $0 \neq r \in R$, the inverse $r^{-1}$ exists in $S$, so we must show $r^{-1} \in R$. Now by Proposition 17.1(ii), $r^{-1}$ satisfies a monic polynomial over $R$, say

$$
r^{-n}+a_{n-1} r^{-n+1}+\cdots+a_{1} r^{-1}+a_{0}=0
$$

for some $a_{i} \in R$. Then multiply by $r^{n-1}$ to get

$$
r^{-1}=-\left(a_{n-1}+a_{n-2} r+\cdots+a_{0} r^{n-1}\right) \in R,
$$

so $R$ is a field.
We will prove the normalisation theorem for infinite fields $K$, and for this the following lemma is crucial:

Lemma 17.3. Let $K$ be an infinite field and $f \in K\left[x_{1}, \ldots, x_{n}\right]$ be a non-zero polynomial. Then there exist $\alpha_{1}, \ldots, \alpha_{n} \in K$ such that $f\left(\alpha_{1}, \ldots, \alpha_{n}\right) \neq 0$.

Proof. We prove this by induction on $n$, with the case $n=0$ being trivial. If now $n=1$ then any non-zero $f \in K\left[x_{1}\right]$ has at most $\operatorname{deg}(f)$ roots. Since $K$ is infinite, we can choose $\alpha_{1}$ not equal to any of these roots and thus $f\left(\alpha_{1}\right) \neq 0$.
Assume now that $n>1$ and the result holds for $n-1$. Let $f \in K\left[x_{1}, \ldots, x_{n}\right]$ be non-zero. If $f \in K\left[x_{1}, \ldots, x_{n-1}\right]$ then we are done, so assume this is not the case. Then we can write

$$
f=g_{r} x_{n}^{r}+\cdots+g_{1} x_{n}+g_{0}
$$

for some $g_{i} \in K\left[x_{1}, \ldots, x_{n-1}\right]$ with $g_{r} \neq 0$. Now by induction, there exist $\alpha_{1}, \ldots, \alpha_{n-1} \in K$ such that $g_{r}\left(\alpha_{1}, \ldots, \alpha_{n-1}\right) \neq 0$. Therefore $f\left(\alpha_{1}, \ldots, \alpha_{n-1}, x_{n}\right) \in K\left[x_{n}\right]$ is a non-zero polynomial, so by the $n=1$ case above we see that there exists $\alpha_{n} \in K$ with $f\left(\alpha_{1}, \ldots, \alpha_{n}\right) \neq 0$.

Now we have all necessary preliminaries to prove the Normalisation Lemma (Thm. 15.6):

Proof of Thm. 15.6. Suppose $S=K\left[y_{1}, \ldots, y_{n}\right]$ and $f \in K\left[x_{1}, \ldots, x_{n}\right]$ is such that $f\left(y_{1}, \ldots, y_{n}\right)=0$, i.e. $y_{1}, \ldots, y_{n}$ are algebraically dependent over $K$. Then choose $\alpha_{1}, \ldots, \alpha_{n-1} \in K$ and set $z_{i}=$ $y_{i}-\alpha_{i} y_{n}$ for $1 \leqslant i \leqslant n-1$. Now let $g \in K\left[x_{1}, \ldots, x_{n}\right]$ be such that

$$
g\left(z_{1}, \ldots, z_{n-1}, y_{n}\right)=f\left(z_{1}+\alpha_{1} y_{n}, \ldots, z_{n-1}+\alpha_{n-1} y_{n}, y_{n}\right)=0
$$

If $f$ has degree $d$ then let $f_{d}$ be the sum of all monomials of $f$ of degree $d$ (the homogeneous piece of $f$ of degree $d$ ). Then

$$
\begin{aligned}
f_{d}\left(z_{1}+\alpha_{1} y_{n}, \ldots, z_{n-1}+\alpha_{n-1} y_{n}, y_{n}\right) & =f_{d}\left(\alpha_{1} y_{n}, \ldots, \alpha_{n-1} y_{n}, y_{n}\right)+\text { lower order terms in } y_{n} \\
& =f_{d}\left(\alpha_{1}, \ldots, \alpha_{n-1}, 1\right) y_{n}^{d}+\text { lower order terms in } y_{n} .
\end{aligned}
$$

Therefore considering $g$ as a polynomial in $y_{n}$ over $K\left[z_{1}, \ldots, z_{n-1}\right]$ we have

$$
g\left(z_{1}, \ldots, z_{n-1}, y_{n}\right)=f_{d}\left(\alpha_{1}, \ldots, \alpha_{n-1}, 1\right) y_{n}^{d}+\text { lower order terms in } y_{n}
$$

Since $f_{d} \neq 0($ as $\operatorname{deg}(f)=d)$, we have by Lemma 17.3 that there exist $\alpha_{1}, \ldots, \alpha_{n-1}$ such that $f_{d}\left(\alpha_{1}, \ldots, \alpha_{n-1}, 1\right) \neq 0$. For this choice we have

$$
f_{d}\left(\alpha_{1}, \ldots, \alpha_{n-1}, 1\right)^{-1} g\left(z_{1}, \ldots, z_{n-1}, y_{n}\right)=0
$$

a monic polynomial over $K\left[z_{1}, \ldots, z_{n-1}\right]$ satisfied by $y_{n}$. Therefore $y_{n}$ is integral over $K\left[z_{1}, \ldots, z_{n-1}\right]$. The proof of the theorem is now by induction on the number $n$ of generators of $S$. Suppose $S=$ $k\left[y_{1}, \ldots, y_{n}\right]$ is such that $y_{1}, \ldots, y_{n}$ are algebraically independent, then we are done. Otherwise there exists some $f \in K\left[x_{1}, \ldots, x_{n}\right]$ such that $f\left(y_{1}, \ldots, y_{n}\right)=0$. Then by the above we can choose $z_{1}, \ldots, z_{n-1} \in S$ such that $y_{n}$ is integral over $S^{*}=K\left[z_{1}, \ldots, z_{n-1}\right]$ and $S=S^{*}\left[y_{n}\right]$. By the induction hypothesis applied to $S^{*}$ there exist elements $w_{1}, \ldots, w_{m} \in S^{*}$ that are algebraically independent over $K$ with $S^{*}$ finite dimensional over $R=K\left[w_{1}, \ldots, w_{m}\right]$. Now since $y_{n}$ is integral over $S^{*}$ it follows by Proposition 17.1(iii) that $S^{*}\left[y_{n}\right]$ is a finite $S^{*}$-algebra. Since both extensions $R \subseteq S^{*}$ and $S^{*} \subseteq S$ are finite, it follows by Proposition 17.1(i) that the extension $R \subseteq S$ is finite as required.

Example 17.4. Let again $S=K[x, y] /\langle x y\rangle=K[\bar{x}, \bar{y}]$. We want to show that $S$ is finite over some $K[z]$. As in the proof of the theorem, $f(\bar{x}, \bar{y})=\bar{x} \cdot \bar{y}=\overline{0}$ in S. Thus we have $d=\operatorname{deg} f=2$. Now we find an $\alpha_{1} \in K$ such that $f\left(\alpha_{1}, 1\right) \neq \overline{0}$, e.g., $\alpha_{1}=1$. Then set $z:=\bar{x}-1 \cdot \bar{y}$ and get $g(z, \bar{y})=f(z+\bar{y}, \bar{y})=(z+\bar{y}) \bar{y}=z \bar{y}+\bar{y}^{2}$. One has $g(z, \bar{y})=\overline{0}$ and thus $S=K[z, y] /\left\langle y z+y^{2}\right\rangle$ is finite over $R=K[z]$.

Now the algebraic version of the Weak Nullstellensatz (Theorem 15.7) can be proven using the results of this Section:

Proof of Theorem 15.7. Using Theorem 15.6 (Noether Normalisation) there exists a polynomial subalgebra $R=K\left[x_{1}, \ldots, x_{r}\right]$ of $S$, over which $S$ is a finite algebra. If $S$ is a field then so is $R$ by Corollary 17.2. If $r \geqslant 1$ then $\left\langle x_{1}\right\rangle$ is a proper ideal in $R$, a contradiction. Therefore $S$ is finitely generated as an $R$-module.
For the second part, suppose $R=K\left[x_{1}, \ldots, x_{n}\right]$ and $\mathfrak{m} \subseteq R$ is a maximal ideal. Then by the first part of the theorem we have that $R / \mathfrak{m}$ is a finite dimensional $K$-algebra. So given $\alpha \in R / \mathfrak{m}$ we have $m(\alpha)=0$ for some monic polynomial $m \in K[t]$ of degree $r$ by Proposition 17.1(ii). Since $K$ is algebraically closed, we can write $m=\left(t-\alpha_{r}\right) \ldots\left(t-\alpha_{r}\right)$ for some $\alpha_{1}, \ldots, \alpha_{r} \in K$. As $R / \mathfrak{m}$ is a field and $m(\alpha)=0$ we have $\alpha=\alpha_{i}$ for some $i$. Therefore $\alpha \in K$ and so $R / \mathfrak{m}=K$. Thus $x_{i}+\mathfrak{m}=a_{i}+\mathfrak{m}$ for some $a_{i} \in K$, and so $\left\langle x_{1}-a_{1}, \ldots, x_{n}-a_{n}\right\rangle \subseteq \mathfrak{m}$. Since both sides are maximal ideals, this is an equality.

Finally, we can prove the weak and strong form of the geometric version of Hilbert's Nullstellensatz (sometimes just denoted by "the Nullstellensatz").

Proof of Theorem 16.12. (Weak): If $I$ is a proper ideal of $K\left[x_{1}, \ldots, x_{n}\right]$ then $I$ is contained in some maximal ideal $\mathfrak{m}$ by Proposition 4.6. But we know from the weak algebraic form of the Nullstellensatz (Theorem 15.7) that $\mathfrak{m}=\left\langle x_{1}-a_{1}, \ldots, x_{n}-a_{n}\right\rangle$ for some $a_{1}, \ldots, a_{n} \in K$. Now $I \subseteq \mathfrak{m} \Longrightarrow$ $\mathbb{V}(\mathfrak{m}) \subseteq \mathbb{V}(I)$ and $\mathbb{V}(\mathfrak{m})=\left\{\left(a_{1}, \ldots, a_{n}\right)\right\} \neq \varnothing$.
(Strong): Note first that if $f \in \sqrt{I}$ then $f^{m} \in I$ for some $m \geqslant 1$. But since $K\left[x_{1}, \ldots, x_{n}\right]$ is an integral domain, the set of zeros of $f^{m}$ is the same as the set of zeros of $f$ (counted without multiplicity). Thus $f \in \mathbb{I}(\mathbb{V}(I))$.
We now show that for all $f \in \mathbb{I}(\mathbb{V}(I))$ we have $f \in \sqrt{I}$. This is obvious for $f=0$ so assume $f \neq 0$. Let $f_{1}, \ldots, f_{m}$ generate $I$ and set

$$
J=\left\langle f_{1}, \ldots, f_{m}, y f-1\right\rangle \subseteq K\left[x_{1}, \ldots, x_{n}, y\right]
$$

for a new variable $y$. Then

$$
\begin{aligned}
\mathbb{V}(J) & =\mathbb{V}\left(\left\langle f_{1}, \ldots, f_{m}, y f-1\right\rangle\right) \\
& =\mathbb{V}\left(\left\langle f_{1}, \ldots, f_{m}\right\rangle\right) \cap \mathbb{V}(\langle y f-1\rangle) \\
& =\mathbb{V}(I) \cap \mathbb{V}(\langle y f-1\rangle)
\end{aligned}
$$

by Proposition 16.5. But since $f \in \mathbb{I}(\mathbb{V}(I))$, any point in $\mathbb{V}(I)$ will not be in $\mathbb{V}(\langle y f-1\rangle)$. Therefore $\mathbb{V}(J)=\varnothing$ and by the weak Nullstellensatz we have $1 \in J$. Hence

$$
1=\sum_{i=1}^{m} g_{i}\left(x_{1}, \ldots, x_{n}, y\right) f_{i}+h\left(x_{1}, \ldots, x_{n}, y\right)(y f-1)
$$

for some $g_{i}, h \in K\left[x_{1}, \ldots, x_{n}, y\right]$. Let $z=\frac{1}{y}$ and choose $N \geqslant \max \left\{\operatorname{deg}\left(g_{1}\right), \ldots, \operatorname{deg}\left(g_{m}\right), \operatorname{deg}(h)+\right.$ $1\}$. Then

$$
\begin{aligned}
z^{N} & =\sum_{i=1}^{m} z^{N} g_{i}\left(x_{1}, \ldots, x_{n}, y\right) f_{i}+z^{N-1} h\left(x_{1}, \ldots, x_{n}, y\right) z(y f-1) \\
& =\sum_{i=1}^{m} \widetilde{g}_{i}\left(x_{1}, \ldots, x_{n}, z\right) f_{i}+\widetilde{h}\left(x_{1}, \ldots, x_{n}, z\right)(f-z)
\end{aligned}
$$

in $K\left[x_{1}, \ldots, x_{n}, z\right]$. Substituting $f$ for $z$ then gives

$$
f^{N}=\sum_{i=1}^{m} \widetilde{g}_{i}\left(x_{1}, \ldots, x_{n}, f\right) f_{i} \in I
$$

so $f \in \sqrt{I}$.

## 18 Gröbner bases

Gröbner bases allow to generalize the Euclidean division algorithm for polynomials in $K[x]$ to several variables. First recall the Euclidean algorithm in one variable:

Let $P(x) \in K[x]$ with $\operatorname{deg}(P)=d$ and let $Q \in K[x]$ be any polynomial. Then there exist unique polynomials $A, B \in K[x]$ such that $\operatorname{deg} B<d$ and

$$
Q=A \cdot P+B
$$

Moreover, $A$ and $B$ may be calculated by a finite algorithm.
One may interpret this using monomial orders (here $<_{\varepsilon}$ is the usual order on $K[x]$ by degree):
$\operatorname{deg} P=d$ means that $\operatorname{lm}_{\varepsilon}(P)=x^{d}$ and $\operatorname{deg} B<d$ means that $\operatorname{lm}_{\varepsilon}(P)$ does not divide any of the
monomials appearing in $B$. If we let $K[x]_{<d}:=\{B \in K[x]: \operatorname{deg} B<d\}$, then we get a $K$-vector space decomposition

$$
K[x]=\langle P\rangle \oplus K[x]_{<d},
$$

or said differently

$$
K[x] /\langle P\rangle \cong K[x]_{<d} \cong K \oplus K x \oplus \cdots \oplus K x^{d-1}
$$

Remark 18.1. This decomposition is particularly easy to find if $P=\operatorname{lm}(P)=x^{d}$ is a monomial. Then for a $Q=\sum_{i=1}^{m} c_{i} x^{i}$ write

$$
Q=\underbrace{\sum_{i=d}^{m} c_{i} x^{i-d}}_{A}+\underbrace{x^{d} \sum_{i=0}^{d-1} c_{i} x^{i}}_{P \cdot B}
$$

So we first consider the case of monomial ideals in $K\left[x_{1}, \ldots, x_{n}\right]$. A monomial ideal is an ideal $I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ such that there exists a (possibly infinite!) set $A \subseteq \mathbb{N}^{n}$ such that $I=\left\langle\underline{x}^{\alpha}: \alpha \in A\right\rangle$. If $P=\sum_{\alpha} a_{\alpha} \underline{\alpha}^{\alpha} \subseteq K\left[x_{1}, \ldots, x_{n}\right]$, then the support of $P$ is $\operatorname{Supp}(P)=\left\{\alpha \in \mathbb{N}^{n}: a_{\alpha} \neq 0\right\}$. Note that $\operatorname{Supp}(P) \subseteq \mathbb{N}^{n}$ is always a finite set.
Further, for any set $A \subseteq \mathbb{N}^{n}$, denote $K[\underline{x}]^{A}:=\{B \in K[\underline{x}]: \operatorname{Supp}(B) \subseteq A\}$.
Lemma 18.2. Let $P_{\alpha}=\underline{x}^{\alpha}, \alpha \in V \subseteq \mathbb{N}^{n}$ and let $E=\bigcup_{\alpha \in V}\left(\alpha+\mathbb{N}^{n}\right)$ and $F=\mathbb{N}^{n} \backslash E$. Then any $Q \in K\left[x_{1}, \ldots, x_{n}\right]$ has a decomposition into

$$
Q=\sum_{\alpha \in V} A_{\alpha} P_{\alpha}+B
$$

where $\sum_{\alpha \in V} A_{\alpha} P_{\alpha}$ is a finite sum and $B$ is a unique polynomial with $\operatorname{Supp}(B) \subseteq F$.
Proof. Since clearly $\mathbb{N}^{n}=E \dot{\cup} F$, it follows that $K[\underline{x}]=K[\underline{x}]^{E} \oplus K[\underline{x}]^{F}$. This means that $K[\underline{x}]=$ $\left\langle P_{\alpha}, \alpha \in V\right\rangle \oplus K[\underline{x}]^{F}$.
Example 18.3. Let $n=2$ and $P_{1}=x^{2}, P_{2}=x y^{2}$, and $P_{3}=y^{4}$. Here $E=\left((2,0)+\mathbb{N}^{2}\right) \cup$ $\left((1,2)+\mathbb{N}^{2}\right) \cup\left((0,4)+\mathbb{N}^{2}\right)$. Then e.g.,

$$
P=x^{5}+x^{3} y^{3}-y=\left(x^{3} P_{1}+x^{2} y P_{2}\right)+(-y)
$$

Here $B=-y$ has $\operatorname{Supp}(B)=\{(0,1)\} \subseteq F$.
Lemma 18.4 (Dickson's lemma). Let $I \subseteq K[\underline{x}]$ be a monomial ideal. Then $I$ is already generated by finitely many monomials. Equivalently, if $E \subseteq \mathbb{N}^{n}$ is an ideal, that is, if $E+\mathbb{N}^{n}=E$, then $E$ is finitely generated, that is, there exists a finite set $V \subseteq \mathbb{N}^{n}$ such that $E=\bigcup_{\alpha \in V}\left(\alpha+\mathbb{N}^{n}\right)$.

This is a special case of Hilbert's basis theorem, so we omit a proof. There exist many direct proofs without using the basis theorem, see e.g. [2].

If $I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ is an arbitrary ideal, then the idea is to "approximate" $I$ by monomial ideals: Choose a monomial order $<_{\varepsilon}$ on $\mathbb{N}^{n}$ and let

$$
\operatorname{lm}_{\varepsilon}(I)=\left\langle\operatorname{lm}_{\varepsilon}(f): f \in I\right\rangle
$$

be the leading ideal of $I$. Clearly, $\operatorname{lm}_{\varepsilon}(I)$ is a monomial ideal, and moreover, one gets a $K$-vector space decomposition

$$
K[\underline{x}]=\operatorname{lm}_{\varepsilon}(I) \oplus K[\underline{x}]^{F},
$$

where $F=\mathbb{N}^{n} \backslash \operatorname{Supp}\left(\operatorname{lm}_{\varepsilon}(I)\right)$. The division theorem will then prove that actually one has $K[\underline{x}]=$ $I \oplus K[\underline{x}]^{F}$.

Example 18.5. Let $P=x^{2}-y$ be in $K[x, y]$ and let $Q=x^{2} y$. Then we can easily find two different ways to express $Q$ as a multiple of $P$ plus a remainder:

$$
Q=y P+y^{2}=x^{2} P+x^{4}
$$

It is not clear which one of the two is preferable!
Theorem 18.6 (Division through 1 polynomial). Let $<_{\varepsilon}$ be a chosen monomial order on $K\left[x_{1}, \ldots, x_{n}\right]$ and let $P \in K\left[x_{1}, \ldots, x_{n}\right]$ with $\operatorname{lm}_{\varepsilon}(P)=\underline{x}^{\alpha}$ for some $\alpha \in \mathbb{N}^{n}$ and denote by $E=\alpha+\mathbb{N}^{n}$ and $F=$ $\mathbb{N}^{n} \backslash E$. Then for any $Q \in K\left[x_{1}, \ldots, x_{n}\right]$ the exist unique polynomials $A, B$ with $B \in K\left[x_{1}, \ldots, x_{n}\right]^{F}$ such that

$$
Q=A \cdot P+B
$$

Moreover, $A$ and $B$ can be calculated with an algorithm.
Remark 18.7. $A$ and $B$ both depend on the monomial order $<_{\varepsilon}$ !
Proof. First we prove the existence (constructively): Without loss of generality assume that the leading coefficient of $P=1$. Let $Q \in K\left[x_{1}, \ldots, x_{n}\right]$, then write

$$
Q=A_{1} \underline{x}^{\alpha}+B_{1}
$$

where $\operatorname{lm}_{\varepsilon}(P)=\underline{x}^{\alpha}$ and $\operatorname{Supp}\left(B_{1}\right) \subseteq F$. Grouped differently

$$
Q=\underbrace{A_{1} P}_{\in\langle P\rangle}+\underbrace{A_{1}\left(\underline{x}^{\alpha}-P\right)}_{=: Q_{1}}+B_{1}
$$

Now write $Q_{1}$ as $Q_{1}=A_{2} \underline{x}^{\alpha}+B_{2}$ with $B_{2} \in K\left[x_{1}, \ldots, x_{n}\right]^{F}$.
Claim: $\operatorname{lm}_{\varepsilon}\left(A_{2}\right)<_{\varepsilon} \operatorname{lm}_{\varepsilon}\left(A_{1}\right)$.
It is enough to show $\operatorname{lm}_{\varepsilon}\left(A_{2} \underline{x}^{\alpha}\right)<_{\varepsilon} \operatorname{lm}_{\varepsilon}\left(A_{1} \underline{x}^{\alpha}\right)$ (properties of monomial orders!). But this holds since $\operatorname{lm}_{\varepsilon}\left(\underline{x}^{\alpha}-P\right)<_{\varepsilon} \operatorname{lm}_{\varepsilon}(P)$. This proves the claim.

Now use induction on $\operatorname{lm}_{\varepsilon}\left(A_{i}\right)$ for $Q_{i}=A_{i} \underline{x}^{\alpha}+B_{i}$. Thus we may assume that $Q_{1}=\tilde{A} \cdot P+\tilde{B}$ with $\tilde{B} \in K\left[x_{1}, \ldots, x_{n}\right]^{F}$. Then

$$
Q=A_{1} P+\tilde{A} P+B+\tilde{B}=\left(A_{1}+\tilde{A}\right) P+(B+\tilde{B})
$$

with $\left(A_{1}+\tilde{A}\right) P \in\langle P\rangle$ and $(B+\tilde{B}) \in K\left[x_{1}, \ldots, x_{n}\right]^{F}$.
For uniqueness assume that $Q=A P+B=A^{\prime} P+B^{\prime}$ with $B, B^{\prime} \in K\left[x_{1}, \ldots, x_{n}\right]^{F}$. Then

$$
0=\left(A-A^{\prime}\right) P+\left(B-B^{\prime}\right)
$$

that is $\left(A-A^{\prime}\right) P=B^{\prime}-B$. Looking at the leading monomials, we see that $\operatorname{lm}_{\varepsilon}\left(B^{\prime}-B\right) \in$ $K\left[x_{1}, \ldots, x_{n}\right]^{F}$ and $\operatorname{lm}_{\varepsilon}\left(\left(A-A^{\prime}\right) P\right) \in K\left[x_{1}, \ldots, x_{n}\right]^{E}$. Since $E \cap F=\varnothing$, also $K\left[x_{1}, \ldots, x_{n}\right]^{E} \cap$ $K\left[x_{1}, \ldots, x_{n}\right]^{F}=0$ and thus $B=B^{\prime}$ and $A=A^{\prime}$.
Example 18.8. Let $Q=x^{2} y$ and $P=x^{2}-y$ in $K[x, y]$. If we choose $<_{\varepsilon}=<_{l e x}$ with $x>y$, then $\operatorname{lm}_{\varepsilon}(P)=x^{2}$ and $Q=y P+y^{2}$. If, on the other hand, we choose $<_{\varepsilon}=<_{l e x}$ with $y>x$, then $\operatorname{lm}_{\varepsilon}(P)=y$ and $Q=\left(-x^{2}\right) P+x^{4}$. In both cases the remainder lies in $K[x, y]^{F}$.
We have proven so far that we have a unique division for principal ideals, but if $I=\left\langle P_{1}, \ldots, P_{m}\right\rangle$, then the remainder depends on the order of the divisions.

Theorem 18.9. Let $<_{\varepsilon}$ be a monomial order on $K\left[x_{1}, \ldots, x_{n}\right]$ and $P_{1}, \ldots, P_{k} \in K\left[x_{1}, \ldots, x_{n}\right]$ with $\operatorname{lm}_{\varepsilon}\left(P_{i}\right)=\underline{x}^{\alpha_{i}}, \alpha_{i} \in \mathbb{N}^{n}$. Then for each $A \in K\left[x_{1}, \ldots, x_{n}\right]$ there exist polynomials $A_{1}, \ldots, A_{k}$ and $B$ such that

$$
Q=\sum_{i=1}^{k} A_{i} P_{i}+B
$$

where $B \in K\left[x_{1}, \ldots, x_{n}\right]^{F}$. (E and $F$ defined as above). Again, there is an algorithm to compute $A_{i}$ and $B$ but they are not unique in general.

Proof. See [2, Chapter 2, $\S 3$, Theorem 3].
Example 18.10. (i) Let $P_{1}=x^{2} y$ and $P_{2}=x y^{2}$ and $Q=x^{3} y^{3}+x y$ in $K[x, y]$. Then

$$
Q=x y^{2} P_{1}+x y=x^{2} y P_{2}+x y=\frac{1}{2} x y^{2} P_{1}+\frac{1}{2} x^{2} y P_{2}+x y .
$$

This shows that the $A_{i}$ are not unique.
(ii) Let $P_{1}=x^{2}-y^{2}$ and $P_{2}=x y-y^{3}$ and $Q=x^{3}$ and choose $<_{\varepsilon}=<_{l e x}$ with $x>y$. Then

$$
Q=x P_{1}+x y^{2}=x P_{1}+y P_{2}+y^{4}
$$

which shows that the remainder is not unique.
In general we can at least make the remainder unique: For this consider an ideal $I=\left\langle P_{1}, \ldots, P_{k}\right\rangle \subseteq$ $K\left[x_{1}, \ldots, x_{n}\right]$. It always holds that

$$
\left\langle\operatorname{lm}_{\varepsilon}\left(P_{1}\right), \ldots, \operatorname{lm}_{\varepsilon}\left(P_{k}\right)\right\rangle \subseteq \operatorname{lm}_{\varepsilon}(I)=\left\langle\operatorname{lm}_{\varepsilon}(P): P \in I\right\rangle
$$

Definition 18.11. A collection of polynomials $P_{1}, \ldots, P_{k} \in K\left[x_{1}, \ldots, x_{n}\right]$ is called a Gröbner basis with respect to a chosen monomial order $<_{\varepsilon}$ if

$$
\left\langle\operatorname{lm}_{\varepsilon}\left(P_{1}\right), \ldots, \operatorname{lm}_{\varepsilon}\left(P_{k}\right)\right\rangle=\operatorname{lm}_{\varepsilon}\left(\left\langle P_{1}, \ldots, P_{k}\right\rangle\right) .
$$

Theorem 18.12. Let $<_{\varepsilon}$ be monomial order on $K\left[x_{1}, \ldots, x_{n}\right]$ and $P_{1}, \ldots, P_{k}$ be a Gröbner basis. Then for each $Q \in K\left[x_{1}, \ldots, x_{n}\right]$ there exist $A_{1}, \ldots, A_{k} \in K\left[x_{1}, \ldots, x_{n}\right]$ and a unique $B \in K\left[x_{1}, \ldots, x_{n}\right]^{F}$ such that

$$
Q=\sum_{i=1}^{k} A_{i} P_{i}+B
$$

Here $B$ is unique but depends on $<_{\varepsilon}$.
Proof. See [2, Chapter 2, §6, Prop. 6].
Remark 18.13. (a) If $I=\langle P\rangle$ is a principal ideal, then $P$ is a Gröbner basis with respect to any monomial order, since $\operatorname{lm}(I)=\left\langle\operatorname{lm}(A P): A \in K\left[x_{1}, \ldots, x_{n}\right]\right\rangle=\langle\operatorname{lm}(A) \cdot \operatorname{lm}(P)\rangle=\langle\operatorname{lm}(P)\rangle$.
(b) Not every set of generators of $I$ is a Gröbner basis, e.g., take $P_{1}=x^{2}-y^{3}$ and $P_{2}=x y-y^{4}$ and the monomial order $<_{\text {lex }}$ with $X>y$. Then $\operatorname{lm}\left(P_{1}\right)=x^{2}$ and $\operatorname{lm}\left(P_{2}\right)=x y$ but $\left\langle x^{2}, x y\right\rangle \subsetneq$ $\operatorname{lm}\left(\left\langle P_{1}, P_{2}\right\rangle\right)$. In order to see this, consider $P_{3}=y P_{1}-x P_{2}=-y^{4}+y^{4} x$ and $P_{4}=P_{3}-y^{2} P_{2}=$ $-y^{4}+y^{7}$. Clearly $\operatorname{lm}\left(P_{4}\right)=y^{7}$ is not contained in $\left\langle x^{2}, x y\right\rangle$.
(c) Let $P_{1}, \ldots, P_{k} \in K\left[x_{1}, \ldots, x_{n}\right], I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ is an ideal, and $<_{\varepsilon}$ a monomial order. If the $P_{i}$ are all contained in $I$ and satisfy $\left\langle\operatorname{lm}_{\varepsilon}\left(P_{1}\right), \ldots, \operatorname{lm}_{\varepsilon}\left(P_{k}\right)\right\rangle=\operatorname{lm}(I)$, then $P_{1}, \ldots, P_{k}$ generate $I$ (see Homework 5!).
(d) Every ideal $I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ has a Gröbner basis: since $\operatorname{lm}(I) \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ is a monomial ideal, we can find finitely many generators, i.e., $\operatorname{lm}(I)=\left\langle\underline{\alpha}^{\alpha_{1}}, \ldots, \underline{x}^{\alpha_{k}}\right\rangle$. By definition of $\operatorname{lm}(I)$ there exist polynomials $P_{1}, \ldots, P_{k}$ in $I$ such that $\operatorname{lm}\left(P_{i}\right)=\underline{x}^{\alpha_{i}}$, so $\left\langle\operatorname{lm}\left(P_{1}\right), \ldots, \operatorname{lm}\left(P_{k}\right)\right\rangle=\operatorname{lm}(I)$. This implies that the $P_{i}$ are a Gröbner basis of $I$.

The next problem is to decide whether a given set of polynomials forms a Gröbner basis with respect to a given monomial order. Furthermore, one wants to construct a Gröbner basis from a given set of polynomials. Both problems will be solved with the following two theorems.

Definition 18.14. Let $P_{1}, \ldots, P_{k}$ and $Q \in K\left[x_{1}, \ldots, x_{n}\right]$ and define $\bar{Q}^{\left(P_{1}, \ldots, P_{k}\right)}$ as the rest of the divisions of $Q$ by $P_{1}, \ldots, P_{k}$ (in this order). If $P_{1}, \ldots, P_{k}$ are a Gröbner basis, then $\bar{Q}^{\left(P_{1}, \ldots, P_{k}\right)}=\bar{Q}^{I}$ is independent of the order of divisions (here $I=\left\langle P_{1}, \ldots, P_{k}\right\rangle$ ).

Lemma 18.15. Let $\mathcal{P}:=\left(P_{1}, \ldots, P_{k}\right)$.
(a) If $Q_{1}, Q_{2} \in K\left[x_{1}, \ldots, x_{n}\right]$, then $\overline{Q_{1}+Q_{2}}{ }^{\mathcal{P}}={\overline{Q_{1}}}^{\mathcal{P}}+{\overline{Q_{2}}}^{\mathcal{P}}$.
(b) If ${\overline{Q_{1}}}^{\mathcal{P}}=0$ and ${\overline{Q_{2}}}^{\mathcal{P}}=0$ and $A_{1}, A_{2}$ are any polynomials, then

$$
{\overline{A_{1} Q_{1}+A_{2} Q_{2}}}^{\mathcal{P}}=0
$$

Proof. Exercise (see Homework 5).
Definition 18.16. Let $<_{\varepsilon}$ be a monomial order on $\mathbb{N}^{n}$. Let $P_{1}, \ldots, P_{k} \in K\left[x_{1}, \ldots, x_{n}\right]$ and define

$$
\operatorname{Rel}\left(P_{1}, \ldots, P_{k}\right):=\left\{R=\left(R_{1}, \ldots, R_{k}\right) \in K\left[x_{1}, \ldots, x_{n}\right]^{k}: \sum_{i=1}^{k} R_{i} P_{i}\right\} .
$$

Then $\operatorname{Rel}\left(P_{1}, \ldots, P_{k}\right) \subseteq K\left[x_{1}, \ldots, x_{n}\right]^{k}$ is a $K\left[x_{1}, \ldots, x_{n}\right]$-module, the module of relations of the $P_{i}$.
Since $K\left[x_{1}, \ldots, x_{n}\right]$ is noetherian, $\operatorname{Rel}\left(P_{1}, \ldots, P_{k}\right)$ is finitely generated, say by $S_{1}, \ldots, S_{m}$ with $S_{j}=$ $\left(S_{j 1}, \ldots, S_{j k}\right)$ for $j=1, \ldots, m$. Written differently, $S_{j} \cdot \mathcal{P}=\sum_{i=1}^{k} S_{j i} P_{i}=0$.

Example 18.17. (a) Let $\mathcal{P}=\left(P_{1}, P_{2}, P_{3}\right)=(x, y, z)$ in $K[x, y, z]^{3}$. Here $\operatorname{Rel}(\mathcal{P})$ is generated by $S_{1}=(y,-x, 0), S_{2}=(z, 0,-x), S_{3}=(0, z,-y)$.
(b) Let $\mathcal{P}=(y z, x z, x y) \in K[x, y, z]^{3}$. Then $\operatorname{Rel}(\mathcal{P})$ is generated by $S_{1}=(x,-y, 0)$ and $S_{2}=$ $(x, 0,-z)$.

Theorem 18.18 (Buchberger's criterion). Let $<_{\varepsilon}$ be monomial order on $\mathbb{N}^{n}$ and let $\mathcal{P}=\left(P_{1}, \ldots, P_{k}\right)$ with $P_{i} \in K\left[x_{1}, \ldots, x_{n}\right]$. Then $P_{1}, \ldots, P_{k}$ are a Gröbner basis with respect to $<_{\varepsilon}$ if and only if for any relation $S \in K\left[x_{1}, \ldots, x_{n}\right]^{k}$ of $\operatorname{lm}_{\varepsilon}\left(P_{1}\right), \ldots, \operatorname{lm}_{\varepsilon}\left(P_{k}\right)$ one has

$$
\overline{S \cdot \overline{\mathcal{P}}^{\mathcal{P}}}=\overline{\sum_{i=1}^{k} S_{i} P_{i}}=0
$$

Equivalently: If $S_{1}, \ldots, S_{m}$ generate $\operatorname{Rel}\left(\operatorname{lm}_{\varepsilon}\left(P_{1}\right), \ldots, \operatorname{lm}_{\varepsilon}\left(P_{k}\right)\right)$, one has

Proof. See [2, Chapter 2, §6, Theorem 6].
Remark 18.19. Relations between the $\operatorname{lm}_{\varepsilon}\left(P_{i}\right)$ can be easily determined: let $\operatorname{lm}\left(P_{i}\right)=\underline{x}^{\alpha_{i}}$. The relations between $\operatorname{lm}\left(P_{1}\right)$ and $\operatorname{lm}\left(P_{2}\right)$ are for example of the form $\underline{x}^{\alpha_{1}} \underline{x}^{\gamma}-\underline{x}^{\alpha_{2}} \underline{x}^{\delta}=0$. Here first set $\omega_{i}=\max \left(\alpha_{1 i}, \alpha_{2 i}\right)$, or equivalently, $\underline{x}^{\omega}=\operatorname{lcm}\left(\underline{x}^{\alpha_{1}}, \underline{x}^{\alpha_{2}}\right)$. Then $\gamma, \delta$ can be determined from $\omega=\alpha_{1}+\gamma=\alpha_{2}+\delta$. Then the relations between $\operatorname{lm}\left(P_{1}\right)$ and $\operatorname{lm}\left(P_{2}\right)$ are generated by the vector $\left(\underline{x}^{\gamma},-\underline{x}^{\delta}, 0, \ldots, 0\right)$. Similarly for the other $\operatorname{lm}\left(P_{i}\right)$ and $\operatorname{lm}\left(P_{j}\right)$.
Note that in general, one also has to take into account the leading coefficients of the $P_{i}$ !
Example 18.20. Let $P_{1}=x y+1$ and $P_{2}=y^{2}-1$ with any monomial order. Then $\operatorname{lm}\left(P_{1}\right)=x y$ and $\operatorname{lm}\left(P_{2}\right)=y_{2}$, and consequently $\operatorname{Rel}\left(x y, y^{2}\right)=\langle(y,-x)\rangle$. We get

$$
(y,-x)\left(P_{1}, P_{2}\right)^{T}=x y^{2}+y-x y^{2}+x=x+y .
$$

Moreover $\overline{x+y}\left(P_{1}, P_{2}\right)=\overline{x+y}^{\left(P_{2}, P_{1}\right)}=x+y \neq 0$. Thus $P_{1}, P_{2}$ are not a Gröbner basis.

Definition 18.21. Let $<_{\varepsilon}$ be a monomial order on $K\left[x_{1}, \ldots, x_{n}\right]$, let $P_{1}, \ldots, P_{k} \in K\left[x_{1}, \ldots, x_{n}\right]$ and let $S_{1}, \ldots, S_{m} \in K\left[x_{1}, \ldots, x_{n}\right]^{k}$ be a generating set of $\operatorname{Rel}\left(\operatorname{lm}_{\varepsilon}\left(P_{1}\right), \ldots, \operatorname{lm}_{\varepsilon}\left(P_{k}\right)\right)$. Then the polynomials $\sum_{i=1}^{k} S_{j i} P_{i}$ are called S-polynomials of $P_{1}, \ldots, P_{k}$ with respect to $<_{\varepsilon}$.
Explicitly, for $P_{i}, P_{j}$ with $\operatorname{lcm}\left(\operatorname{lm}_{\varepsilon}\left(P_{i}\right), \operatorname{lm}_{\varepsilon}\left(P_{j}\right)\right)=\underline{x}^{\omega}$, the $S$-polynomial is

$$
S\left(P_{i}, P_{j}\right)=\frac{\underline{x}^{\omega}}{\operatorname{lt}\left(P_{i}\right)} \cdot P_{i}-\frac{\underline{x}^{\omega}}{\operatorname{lt}\left(P_{j}\right)} \cdot P_{j} .
$$

(Note that here $\operatorname{lt}(f)$ stands for the leading term of the polynomial $f$, so we are also inverting the leading coefficients here!)

Remark 18.22. The name "S-polynomial" comes from the word syzygy, and this words stands for the relations between polynomials: the relations between polynomials $P_{1}, \ldots, P_{k}$ are called first syzygies, the relations between the first syzygies are the second syzygies, and so on.
By Buchberger's criterion, $P_{1}, \ldots, P_{k}$ are a Gröbner basis with respect to $<_{\varepsilon}$ if and only if all $S$ polynomials reduce to 0 after division through $P_{1}, \ldots, P_{k}$.

Example 18.23. Let $P_{1}=y-x^{2}$ and $P_{2}=z-x^{3}$ in $K[x, y, z]$ and choose $<_{l e x}$ with $y>z>x$. Then $\operatorname{lm}\left(P_{1}\right)=y$ and $\operatorname{lm}\left(P_{2}\right)=z$. The relations between these two monomials are generated by $S_{1}=(z,-y)$. Then

$$
S_{12}:=S\left(P_{1}, P_{2}\right)=z P_{1}-y P_{2}=x^{3} y-x^{2} z
$$

The leading monomial $\operatorname{lm}\left(S_{12}\right)=x^{3} y$ is divisible by $\operatorname{lm}\left(P_{1}\right)$, so we get ${\overline{S_{12}}}^{P_{1}}=x^{5}-x^{2} z=x^{2} P_{2}$. Thus $\overline{S_{12}}{ }^{\left(P_{1}, P_{2}\right)}=0$ and it follows that $P_{1}$ and $P_{2}$ are a Gröbner basis.

If, on the other hand, we choose $<_{l e x}$ with $x>y>z$, then $\operatorname{lm}\left(P_{1}\right)=x^{2}$ and $\operatorname{lm}\left(P_{2}\right)=x^{3}$, and in this case $S_{12}=x y-z$. No monomial of $S_{12}$ is divisible by $\operatorname{lm}\left(P_{1}\right)$ or $\operatorname{lm}\left(P_{2}\right)$, so it follows that $P_{1}$ and $P_{2}$ are not a Gröbner basis with respect to this order.

Remark 18.24. The right choice of a monomial order can sometimes simplify computations significantly! In particular useful here are the linear orders, that were defined in Section 5.

Theorem 18.25 (Buchberger's algorithm). Let $P_{1}, \ldots, P_{k} \in K\left[x_{1}, \ldots, x_{n}\right]$ and choose a monomial order $<_{\varepsilon}$ on $\mathbb{N}^{n}$. Define for $m \in \mathbb{N}$ the following vectors: $F^{0}:=\left(P_{1}, l\right.$ dots,$\left.P_{k}\right), F^{1}:=\left(P_{1}, \ldots, P_{k}, S_{i j}\right.$ for $1 \leqslant i<j \leqslant k)$, and

$$
F^{m+1}:=\left(F^{m}, \text { all S-polynomials of components of } F^{m}\right) .
$$

(Here we mean S-polynomials after reduction by $F^{m}$ !). Then there exists an $m_{0}$ such that $F^{m_{0}}$ is a Gröbner basis with respect to $<_{\varepsilon}$.

Proof. See [2, Chapter 2, §7, Theorem 2].
This algorithm yields a Gröbner basis but can be computationally complex.

## Applications

We list here a few applications of Gröbner bases - many more can be found in e.g. [2, 3, 4].

## Ideal membership

Let $I=\left\langle P_{1}, \ldots, P_{k}\right\rangle$ be an ideal in $K\left[x_{1}, \ldots, x_{n}\right]$ and $Q$ any polynomial. How can one determine whether $Q \in I$ ?
To answer this question, first use Buchberger's algorithm (Theorem 18.25) to complete $P_{1}, \ldots, P_{k}$ to a Gröbner basis $P_{1}^{\prime}, \ldots, P_{m}^{\prime}$ of $I$ (with respect to a suitably chosen monomial order $<_{\varepsilon}$ ). Then set $\mathcal{P}^{\prime}=\left(P_{1}^{\prime}, \ldots, P_{m}^{\prime}\right)$ and calculate $\bar{Q}^{\mathcal{P}^{\prime}}=: B$. If $B \neq 0$, then $Q \notin I$. If $B=0$, then $Q \in I$.
Remark 18.26. If $\operatorname{lm}_{\varepsilon}(Q) \notin \operatorname{lm}_{\varepsilon}(I)$, then $Q \notin I$.

## Solving polynomial systems of equations

Let $I=\left\langle P_{1}, \ldots, P_{k}\right\rangle \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ be an ideal (here we assume that $K=\bar{K}$ is algebraically closed). Then by Hilbert's Nullstellensatz, $\mathbb{V}(I)=\varnothing$ if and only if $I=K\left[x_{1}, \ldots, x_{n}\right]$ if and only if $1 \in I$. So in order to determine whether the system of polynomial equations $\left\{P_{1}=\cdots=P_{k}=0\right\}$ has a solution, we just need to check whether $1 \in I$. This can be done with the method from above.

Explicitly determining solutions for the system requires some more work:

## Elimination

The idea of elimination is that for a system of polynomial equations in $n$ variables, one first tries to eliminate some of the variables.

Theorem 18.27. Let $I \subseteq K\left[x_{1}, \ldots, x_{n}\right]$ be an ideal, let $<_{\varepsilon}=<_{\text {lex }}$ with $x_{1}>x_{2}>\ldots>x_{n}$, and let $P_{1}, \ldots, P_{k}$ be a Gröbner basis of I with respect to this order. Set $F:=\left\{P_{1}, \ldots, P_{k}\right\}$ and for $l \leqslant n$ let $I_{l}:=I \cap K\left[x_{l+1}, \ldots, x_{n}\right]$.
Then $F_{l}:=F \cap K\left[x_{l+1}, \ldots, x_{n}\right]$ is a Gröbner basis of the ideal $I_{l}$ with respect to $<_{\text {lex }}$ on $\mathbb{N}^{n-l}$ with $x_{l+1}>\ldots>x_{n}$.
Proof. See [2, Chapter 3, $\S 1$, Theorem 2].
Example 18.28. Solve the system of equations in $\mathbb{C}^{3}$

$$
\begin{aligned}
& P_{1}:=x^{2}+y^{2}+z^{2}-1=0 \\
& P_{2}:=x^{2}+y^{2}-z=0 \\
& P_{3}:=x-z=0
\end{aligned}
$$

Consider $I=\left\langle P_{1}, P_{2}, P_{3}\right\rangle$ and compute a Gröbner basis of $I$ with respect to $<_{l e x}$ with $x>y>z$. The Gröbner basis is given by the three polynomials $P_{1}^{\prime}=x-z, P_{2}^{\prime}=y^{2}+z^{2}-z, P_{3}^{\prime}=z^{2}+z-1$. The third elimination ideal is then $I_{3}=\left\langle z^{2}+z-1\right\rangle$, which yields two possibilities $z_{+/-}=\frac{-1 \pm \sqrt{5}}{2}$. The second elimination ideal is $I_{2}=\left\langle P_{3}^{\prime}, P_{2}^{\prime}\right\rangle$. Plugging both values for $z$ into $P_{3}^{\prime}=P_{2}^{\prime} \stackrel{2}{=} 0$, we obtain two solutions for $y$ in each case (in total 4 pairs of solutions $(y, z)$, two of them with imaginary $y$-values coming from $z_{-}$). Now $I_{1}=\left\langle P_{1}^{\prime}, P_{2}^{\prime}, P_{3}^{\prime}\right\rangle$ and $P_{1}^{\prime}=P_{2}^{\prime}=P_{3}^{\prime}=0$ has reduced to one equation in one variable. In total one gets 4 different triples of solutions $(x, y, z) \in \mathbb{C}^{3}$ :

$$
\begin{array}{r}
\quad\left(\frac{\sqrt{5}-1}{2}, \sqrt{-2+\sqrt{5}}, \frac{\sqrt{5}-1}{2}\right),\left(\frac{\sqrt{5}-1}{2},-\sqrt{-2+\sqrt{5}}, \frac{\sqrt{5}-1}{2}\right), \\
\left(\frac{-\sqrt{5}-1}{2}, i \sqrt{2+\sqrt{5}}, \frac{-\sqrt{5}-1}{2}\right),\left(\frac{-\sqrt{5}-1}{2},-i \sqrt{2+\sqrt{5}}, \frac{-\sqrt{5}-1}{2}\right) .
\end{array}
$$

Interpreted geometrically, this means that the three surfaces defined by $P_{1}, P_{2}$ and $P_{3}$ intersect in 4 different points in $\mathbb{C}^{3}$, and only two of them are real.

Other application of Gröbner bases include: computation of radical of an ideal, intersection of ideals, ideal quotient, Gauss algorithm, ...

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[^0]:    ${ }^{1}$ Named after Emmy Noether (1882-1935),

