# An Introduction to Gröbner Bases 

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#### Abstract

The aim of this paper is to provide an undergraduate friendly introduction to the concept of Gröbner bases. After motivating the subject, we build a solid foundation before formally defining, proving the existence of, and providing the original Buchberger's algorithm to compute, said Gröbner basis. We conclude by briefly summarizing a few applications.


## 1 Prerequisites

It is expected that the reader has completed an undergraduate course in ring theory and/or is familiar with multi-variable polynomial ring concepts and results. We highlight a few of these which are used directly in definitions and proofs.

Definition 1. Let $K$ be a field (in more generality, $K$ can be any commutative ring). A monomial in the polynomial ring $K\left[x_{1}, \ldots, x_{n}\right]$ is represented by $c x^{\alpha}$ where $c$ is an element of $K, x=\left(x_{1}, \ldots, x_{n}\right)$ and $\alpha=\left(a_{1}, \ldots, a_{n}\right)$ in $\mathbb{Z}_{\geq 0}^{n}$. For example, $x^{2} y^{3}$.
Lemma 1. ([1]) Let $I=\left\langle x^{\alpha} \mid \alpha \in S \subseteq \mathbb{Z}_{>0}^{n}\right\rangle$ be a monomial ideal (ideal generated by monomials) in $K\left[x_{1}, \ldots, x_{n}\right]$. Then, a monomial $x^{\beta}$ is an element of $I$ if and only if $x^{\alpha}$ divides $x^{\beta}$ for some $\alpha$ in $S$.

Definition 2. A commutative ring with unity $K$ is called Noetherian if and only if for every increasing chain of ideals $I_{1} \subseteq I_{2} \subseteq I_{3} \subseteq \ldots$ of $K$, there exists $N$ in $\mathbb{Z}_{\geq 0}$ such that $I_{n}=I_{N}$ for all $n \geq N$.
We have an equivalent definition of Noetherian:
Definition 2. A commutative ring with unity $K$ is called Noetherian if and only if every ideal of $K$ is finitely generated.

Theorem 1 (Hilbert Basis Theorem). ([2]) If $K$ is a Noetherian ring then $K[x]$ is a Noetherian ring. Inductively, $K\left[x_{1}, \ldots, x_{n}\right]$ is a Noetherian ring.

Proposition 1. ([2]) Every field is Noetherian.

## 2 Introduction

In his 1965 Ph.D. thesis, Bruno Buchberger introduced the new concept of a Gröbner basis, named after his advisor Wolfgang Gröbner. He provided the Buchberger Algorithm to compute them (the one we will state and prove in this paper) at the same time ([3]).

We acknowledge that many years earlier in 1913, Nikolai Günther of Russia made a similar discovery; his work was published in Russian journals but ignored internationally until it was recognized in the 1980s ([3]).

The Buchberger algorithm is motivated by what is known as the Ideal Membership Problem (6) which asks: given a field $K$ and an ideal $I$ of $K\left[x_{1}, \ldots x_{n}\right]$, how can we determine if $f$ in $K\left[x_{1}, \ldots, x_{n}\right]$ is an element of $I$ ? We know that for a commutative ring, we can write an element of $I$ as a linear combination of the generators. Thus, as one would do in the single variable case, it would be intuitive to divide our polynomial by the generators of $I$ simultaneously using the division algorithm for multi-variable polynomial rings (3), and if the algorithm terminates with zero remainder, we have an element of $I$. Alas, we will show in the following section that we run into several issues using this method. Namely, unlike the division algorithm in one variable, the multi-variable division algorithm does not output a unique remainder unless the order of divisors is fixed. Consequently, an element of $I$ can have a non-zero remainder. This dilemma says that zero remainder is sufficient for ideal membership but not necessary. We are left to wonder: can we find a basis such that the unique remainder is independent of divisor order, thus providing a sufficient and necessary way to determine elements of $I$ ? Indeed! The Gröbner basis of $I$ !

This relatively modern idea has led to many fascinating applications and discoveries in mathematics and science.

## 3 Ordering and Division Algorithm in $K\left[x_{1}, \ldots, x_{n}\right]$

In one variable, we are familiar with the terms degree of a monomial, and degree of a polynomial. For example, $4 x^{2}$ is a degree 2 monomial and $4 x^{2}+x+1$ is a degree 2 polynomial in $\mathbb{Q}[x]$. How do these definitions change for multi-variable polynomials?
Let us consider the two variable case: $\mathbf{x}=(x, y)$. Suppose you have the monomials $x^{2} y$ and $x y^{2}$. Both have the same total degree $(1+2=2+1=3)$. Which one would you define as greater than the other? In fact, it depends on the definition of monomial ordering you choose. We provide a few common examples of monomial orderings before outlining the formal criteria.

Definition 3. (Lexicographic Order) Let $\alpha=\left(a_{1}, \ldots, a_{n}\right), \beta=\left(b_{1}, \ldots, b_{n}\right) \in \mathbb{Z}_{\geq 0}^{n}$. If $\alpha-\beta$ has a positive leftmost non-zero entry, then we say the degree $\alpha$ is greater than the degree $\beta$ with respect to lexicographic ordering, abbreviated as $\alpha>_{\text {lex }} \beta$. We can equivalently say that $x^{\alpha}$ is greater than $x^{\beta}$ as monomials.
Example. $x>y>z$ since $x=x^{1} y^{0} z^{0}, y=x^{0} y^{1} z^{0}, z=x^{0} y^{0} z^{1}$ and
$(1,0,0)-(0,1,0)=(1,-1,0),(0,1,0)-(0,0,1)=(0,1,-1)$, i.e. we have alphabetical
ordering.
$\underline{\text { Definition 4. (Graded Lex Order) }}$ Let $\alpha, \beta$ as above. If

$$
|\alpha|=\sum_{i=1}^{n} a_{i}>|\beta|=\sum_{i=1}^{n} b_{i}
$$

or

$$
|\alpha|=|\beta| \text { and } \alpha>_{\text {lex }} \beta,
$$

then we say $\alpha>{ }_{\text {grlex }} \beta$.
$|\alpha|$ is denoted the total degree.
Example. We first notice that again we have alphabetical ordering of the variables. Another example is $x^{4} y^{7} z>_{\text {grlex }} x^{4} y^{2} z^{3}$ since $|(4,7,1)|=12>9=|(4,2,3)|$.
$\underline{\text { Definition 5. (Graded Reverse Lex Order) }}$ Let $\alpha, \beta$ as above. If

$$
|\alpha|=\sum_{i=1}^{n} a_{i}>|\beta|=\sum_{i=1}^{n} b_{i}
$$

or

$$
|\alpha|=|\beta| \text { and the rightmost nonzero entry of } \alpha-\beta \text { is negative, }
$$

then we say $\alpha>_{\text {grevlex }} \beta$.
Example. We have the same example as above: $x^{4} y^{7} z>_{\text {grevlex }} x^{4} y^{2} z^{3}$ because grevlex and grlex both order by total degree first, but break ties in different ways.
Another example is $(2,3,2)>_{\text {grevlex }}(0,0,7)$ since $(2,3,2)-(0,0,7)=(2,3,-5)$.
Definition 6. We call $>$ a monomial ordering on the set of monomials $\left\{x^{\alpha} \mid x=\left(x_{1}, \ldots, x_{n}\right), \alpha \in \mathbb{Z}_{\geq 0}^{n}\right\}$ in $K\left[x_{1}, \ldots, x_{n}\right]$ if it is a relation satisfying the following conditions:
(i) $>$ is a total (linear) ordering, i.e. exactly one of

$$
x^{\alpha}>x^{\beta}, x^{\alpha}=x^{\beta}, x^{\beta}>x^{\alpha}
$$

is true and the ordering in transitive.
(ii) If $x^{\alpha}>x^{\beta}, \gamma \in \mathbb{Z}_{\geq 0}^{n}$, then $x^{\alpha+\gamma}>x^{\beta+\gamma}$
(iii) $>$ is a well-ordering, i.e. there exists a (not necessarily unique) minimal monomial.

From now on, > represents an arbitrary fixed monomial ordering. You may be wondering... Do we have to impose an ordering? If we do not, then we also run into uniqueness of remainder problems as mentioned in the introduction.

We need just a few more pieces of terminology before we have the foundation to discuss Gröbner bases.

Definition 7. Let $f=\sum_{i=1}^{n} c_{i} x^{\alpha_{i}}$ in $K\left[x_{1}, \ldots, x_{n}\right]$ (without loss of generality, none of the $\alpha_{i}$ are equal, otherwise combine them). The leading term of $f, L T(f)=c_{k} x^{\alpha_{k}}$, is such that $\alpha_{k}>\alpha_{i}$ for all $i . c_{k}$ is denoted the leading coefficient.

Definition 8. As above, let $f=\sum_{i=1}^{n} c_{i} x^{\alpha_{i}}$ be a polynomial in $K\left[x_{1}, \ldots, x_{n}\right]$. The leading monomial of $f$ is $L M(f)=x^{\alpha_{k}}$ such that $\alpha_{k}>\alpha_{i}$ for all $i$.

Remark Note that we distinguish the leading monomial as the leading term without its coefficient. This is different than the general definition of a monomial.

For the definition and proof of the division algorithm in $K\left[x_{1}, \ldots, x_{n}\right]$, we reference [1] (chapter 2.3). The computation is easy to grasp by examples, which follow. Like normal division, we have a divisor, dividend, quotient and remainder; what's different is that we can have multiple divisors and quotients. Take a peak at the example below to familiarize yourself with the setup.

Essentially, we start by fixing a monomial ordering. This determines leading terms. Then, we apply the division algorithm as we normally would in one variable to the top most divisor (which we will call $f_{1}$ ) and the dividend. That is, we ask, what can we multiply the leading term of $f_{1}$ by to get the leading term of the dividend? For example, $x y$ can be multiplied by $x$ to get $x^{2} y$ below.

Dividing by $f_{1}$ builds quotient one. If in any particular step we cannot divide by divisor one (for example, $x y$ cannot be multiplied by anything to get $y^{2}$ below), we leave quotient one empty, move to $f_{2}$ and build quotient two, and so on. At each step we return to $f_{1}$ and work down the list.

If at any step we cannot divide the leading term by any of the divisors, that leading term becomes the remainder of that step (bold below) and we continue the division by carrying down the rest of the terms. At the end, we add the remainders of each step for a total remainder.

As we will see in the coming examples, the choice of ordering can simplify or complicate polynomial division in $K\left[x_{1}, \ldots, x_{n}\right]$ (sometimes greatly). We also see that, unless order of divisors is fixed, the remainder is not unique - even with a fixed ordering.

The following problems and the worked example in 4 are taken from [1], where additional examples and exercises can be found, but worked out independently.

Example. Order of divisors matters. We use lex ordering in the following two divisions and work in $\mathbb{Q}[x, y]$.

To begin, the polynomials are ordered (greatest monomial to the left). Step one is multiplying $x y-1$ by $x$ to get the leading term of the dividend. We write the product below the dividend and subtract, and add $x$ to quotient one.

$$
\begin{array}{r}
\left.\begin{array}{r}
q_{1}=x \\
f_{1}=x y-1 \\
f_{2}=y^{2}-1
\end{array}\right) \\
\frac{-\left(x^{2} y-x\right)}{x y^{2}+x+y^{2}}
\end{array}
$$

Step two is multiplying $x y-1$ by $y$ to get the leading term of the difference. We write the
new product below the first difference and subtract, and add $y$ to quotient one. We encourage the reader to follow through the remaining steps.

$$
\begin{array}{r}
\left.\begin{array}{r}
q_{1}=x+y \\
q_{2}=1 \\
f_{1}=x y-1 \\
f_{2}=y^{2}-1
\end{array}\right) \\
\frac{x^{2} y+x y^{2}+y^{2}}{-}\left(x^{2} y-x\right) \\
x y^{2}+x+y^{2} \\
\frac{-\left(x y^{2}-y\right)}{\mathbf{x}+y^{2}+y} \\
y^{2}+y \\
\\
-\quad\left(y^{2}-1\right) \\
\frac{\mathbf{y}+1}{\mathbf{1}}
\end{array}
$$

We finish with remainder $\mathbf{x}+\mathbf{y}+\mathbf{1}$.
We conclude that $x^{2} y+x y^{2}+y^{2}=(x+y) \cdot(x y-1)+1 \cdot\left(y^{2}-1\right)+x+y+1$.

On the other hand, switching the order of divisors gives

$$
\begin{aligned}
& \left.\begin{array}{c}
f_{1}=y^{2}-1 \\
f_{2}=x y-1
\end{array}\right) \frac{\begin{array}{c}
q_{1}=x+1 \\
q_{2}=x
\end{array}}{x^{2} y+x y^{2}+y^{2}} \\
& \frac{-\left(x^{2} y-x\right)}{x y^{2}+x+y^{2}} \\
& \frac{-\quad\left(x y^{2}-x\right)}{\frac{2 \mathbf{x}+y^{2}}{y^{2}}} \\
& -\quad\left(y^{2}-1\right)
\end{aligned}
$$

with remainder $\mathbf{2 x}+\mathbf{1}$ - a different remainder!
We conclude that $x^{2} y+x y^{2}+y^{2}$ can also be written as $(x+1) \cdot\left(y^{2}-1\right)+x \cdot(x y-1)+2 x+1$.
We now compute the second division again, this time using grlex ordering:

$$
\begin{aligned}
& \left.\begin{array}{l}
f_{1}=y^{2}-1 \\
f_{2}=x y-1
\end{array}\right) \frac{\begin{array}{c}
q_{1}=x+1 \\
q_{2}=x
\end{array}}{x^{2} y+x y^{2}+y^{2}} \\
& \frac{-\left(x^{2} y-x\right)}{x y^{2}+y^{2}+x} \\
& \frac{-\quad\left(x y^{2}-x\right)}{y^{2}+2 x} \\
& \frac{-\quad\left(y^{2}-1\right)}{2 \mathbf{x}+\mathbf{1}}
\end{aligned}
$$

We highlight that we were able to continue dividing by $d_{1}$ at row 5 , unlike before, because $y^{2}>_{\text {grlex }} 2 x$ whereas $2 x>_{\text {lex }} y^{2}$.

## 4 Algorithm to Construct Gröbner Bases and Proof of Existence

The preceding section brings to light the obstacle that the Ideal Membership Problem faces. With a fixed ordering, we now construct a new basis from our original basis of $I$ such that division in any order of this new basis results in a unique remainder.

Definition 9. Let $K$ be a field. Let $I$ be an ideal of $K\left[x_{1}, \ldots, x_{n}\right]$. We say the subset $G=\left\{g_{1}, \ldots, g_{d}\right\}$ of $I$ is a Gröbner basis (with respect to fixed ordering) for $I$ if $\left\langle L T\left(g_{1}\right), \ldots, L T\left(g_{d}\right)\right\rangle=\langle L T(I)\rangle$, where $\langle L T(I)\rangle$ is the ideal generated by the set of leading terms of elements of $I$.

Remark 1. $G$ generates $I$; that is, $\left\langle g_{1}, \ldots, g_{d}\right\rangle=I$.
Proof. $\left\langle g_{1}, \ldots, g_{d}\right\rangle$ is contained in $I$ since $G$ is contained in $I$, by the properties of ideals. Conversely, let $f$ be an element of $I$. We can write $f=p_{1} g_{1}+\ldots+p_{d} g_{d}+r$ for some remainder $r$ using the division algorithm. Since $f$ is in $I, r=f-p_{1} g_{1}+\ldots+p_{d} g_{d}$ is an element of $I$. Assume $r$ is non-zero. By the definition of Gröbner basis then, $L T(r)$ is in $\left\langle L T\left(g_{1}\right), \ldots, L T\left(g_{d}\right)\right\rangle$, and by Lemma 1 some $L T\left(g_{i}\right)$ divides $L T(r)$. But this contradicts that $r$ is the smallest remainder, thus $r$ must be zero, and $f$ is in $\left\langle g_{1}, \ldots, g_{d}\right\rangle$.

Remark 2. Note that by 1, since $K$ is a field, $K$ is Noetherian, and thus $K\left[x_{1}, \ldots, x_{n}\right]$ is Noetherian. Hence, every ideal in $K\left[x_{1}, \ldots, x_{n}\right]$ is finitely generated. The existence of this generating set is what we will eventually build a Gröbner basis from. A Gröbner basis is a specific kind of basis, and in fact is more of a spanning set unless minimal, and non-unique unless reduced 5 .

Definition 10. Let $x^{\alpha}, x^{\beta}$ be monomials. Let $\operatorname{lcm}\left(x^{\alpha}, x^{\beta}\right)=x^{\gamma}, \gamma=\left(c_{1}, \ldots c_{n}\right)$ and $c_{i}=\max \left\{a_{i}, b_{i}\right\}$. We define the $S$-polynomial: $S(f, g)=\frac{x^{\gamma}}{L T(f)} \cdot f-\frac{x^{\gamma}}{L T(g)} \cdot g$ where $x^{\gamma}=\operatorname{lcm}(L M(f), L M(g))$.
Theorem 2 (Buchberger's Criterion). Let $I$ be an ideal of $K\left[x_{1}, \ldots, x_{n}\right]$ such that $I=\left\langle g_{1}, \ldots g_{d}\right\rangle$. Then $G=\left\{g_{1}, \ldots, g_{d}\right\}$ is a Gröbner basis if and only if ${\overline{S\left(g_{i}, g_{j}\right)}}^{G}=0$ for all $i, j$ in $\{1, \ldots, d\}$, where ${\overline{S\left(g_{i}, g_{j}\right)}}^{G}$ denotes the remainder of $S\left(g_{i}, g_{j}\right)$ divided by $G$.

Proof. ( $\Leftarrow$ :) We need to show that $\left\langle L T\left(g_{1}\right), \ldots L T\left(g_{d}\right)\right\rangle=\langle L T(I)\rangle$.
$\left\langle L T\left(g_{1}\right), \ldots L T\left(g_{d}\right)\right\rangle \subseteq\langle L T(I)\rangle$ is clear since $g_{i}$ in $I$ implies $L T\left(g_{i}\right)$ is in $\langle L T(I)\rangle$ for all $i$ in $\{1, \ldots, d\}$. Thus $\left\langle L T\left(g_{1}\right), \ldots L T\left(g_{d}\right)\right\rangle \subseteq\langle L T(I)\rangle$ since $\langle L T(I)\rangle$ is closed under finite linear combinations and an arbitrary element of $\left\langle L T\left(g_{1}\right), \ldots L T\left(g_{d}\right)\right\rangle$ is written $f_{1} L T\left(g_{1}\right)+\ldots+f_{d} L T\left(g_{d}\right)$ for $f_{i}$ in $K\left[x_{1}, \ldots, x_{n}\right]$.

It remains to show that $\langle L T(I)\rangle \subseteq\left\langle L T\left(g_{1}\right), \ldots L T\left(g_{d}\right)\right\rangle$, and by the same reasoning as above, it is sufficient to show that $L T(I) \subseteq\left\langle L T\left(g_{1}\right), \ldots L T\left(g_{d}\right)\right\rangle$. Let $f$ be arbitrary in $I$ and $L T(f)=a_{n} x^{\alpha_{n}}$ in $L T(I)$. Write:

$$
f=a_{n} x^{\alpha_{n}}+a_{n-1} x^{\alpha_{n-1}}+\ldots+a_{0}=\sum_{i=1}^{d} h_{i} g_{i}, h_{i} \in K\left[x_{1}, \ldots, x_{n}\right]
$$

We can write $f$ this way since $I=\left\langle g_{1}, \ldots, g_{d}\right\rangle$.
By Lemma 1 , we need to show that $a_{n} x^{\alpha_{n}}$ is divisible by one of $\operatorname{LT}\left(g_{i}\right), \ldots, L T\left(g_{d}\right)$.
Let $\delta=\max \left(\operatorname{deg}\left(h_{i} g_{i}\right)\right)$. There are two cases:

1. $\delta=\operatorname{deg}(f)=\alpha_{n}$, the case with no cancellation in the sum of $f$.
2. $\delta>\operatorname{deg}(f)$, the case with cancellation in the sum of $f$.

Let's deal with case 1 first, then we will produce $f=\sum_{i=1}^{d} h_{i}^{\prime} g_{i}$ such that $\max \left(\operatorname{deg}\left(h_{i}^{\prime} g_{i}\right)\right)=\delta^{\prime}$ is strictly less than $\delta$, so we can inductively reduce to case 1 .

Suppose $f=h_{1} g_{1}+\ldots+h_{d} g_{d}$ with $\operatorname{deg}(f)=\delta$ and $\operatorname{deg}\left(h_{i} g_{i}\right) \leq \delta$ (at least one equality). Then for some $i$ in $\{1, \ldots, d\}, L M(f)=L M\left(h_{i} g_{i}\right)=L M\left(g_{i}\right) L M\left(h_{i}\right)$. So in this case, $L M\left(g_{i}\right)$ divides $L M(f)$. This gives that $L T\left(g_{i}\right)$ divides $L T(f)$ since $L M(f)=x^{\delta}=L M\left(h_{i}\right) L M\left(g_{i}\right)$ implies

$$
\begin{aligned}
a_{n} x^{\delta} & =a_{n} L M\left(h_{i}\right) L M\left(g_{i}\right) \\
& =\frac{a_{n}}{b_{i}} L M\left(h_{i}\right) \cdot b_{i} L M\left(g_{i}\right) \\
& :=\frac{a_{n}}{b_{i}} L M\left(h_{i}\right) L T\left(g_{i}\right)
\end{aligned}
$$

where $b_{i} \neq 0$ since $L T\left(g_{i}\right) \neq 0$. This gives case 1 .
We now look at case 2. We will need a lemma which we state and prove now.
Lemma 2. If $\operatorname{deg}\left(p_{i}\right)=\delta$ for all $i$ and $\operatorname{deg}\left(\sum_{i=1}^{n} p_{i}\right)<\delta$, then there exists $a_{i j}$ in $K$ such that $\sum_{i=1}^{n} p_{i}=\sum_{i, j}^{n} a_{i j} S\left(p_{i}, p_{j}\right)$.

Proof. Write $S\left(p_{i}, p_{j}\right)=\frac{x^{\delta}}{L T\left(p_{i}\right)} \cdot p_{i}-\frac{x^{\delta}}{L T\left(p_{j}\right)} \cdot p_{j}$ where $L T\left(p_{i}\right)=b_{i} x^{\delta}$ and $L T\left(p_{j}\right)=b_{j} x^{\delta}$.
Fix $j$. Then,

$$
\begin{aligned}
\sum_{\substack{i=1 \\
i \neq j}}^{n} b_{i}\left(\frac{p_{i}}{b_{i}}-\frac{p_{j}}{b_{j}}\right) & =\sum_{\substack{i=1 \\
i \neq j}}^{n}\left(p_{i}-\frac{b_{i}}{b_{j}} p_{j}\right) \\
& =\sum_{\substack{i=1 \\
i \neq j}}^{n} p_{i}-\frac{p_{j}}{b_{j}} \sum_{\substack{i=1 \\
i \neq j}}^{n} b_{i} \\
& =\sum_{i=1}^{n} p_{i}
\end{aligned}
$$

where $\left(\frac{p_{i}}{b_{i}}-\frac{p_{j}}{b_{j}}\right)=S\left(p_{i}, p_{j}\right)$ and $\sum_{\substack{i=1 \\ i \neq j}}^{n} b_{i}=0$ since the sum has degree less than $\delta$.

Continuing, we now write

$$
f=\sum_{\operatorname{deg}\left(h_{i} g_{i}\right)>\operatorname{deg}(f)} h_{i} g_{i}+\sum_{\operatorname{deg}\left(h_{i} g_{i}\right) \leq \operatorname{deg}(f)} h_{i} g_{i}
$$

and since $\operatorname{deg}\left(h_{i} g_{i}\right)=\delta$ implies $\operatorname{deg}\left(L T\left(h_{i}\right) g_{i}\right)=\delta$ by well-ordering, we have

$$
\begin{aligned}
& =\sum_{\operatorname{deg}\left(h_{i} g_{i}\right)=\delta} L T\left(h_{i}\right) g_{i}+\sum_{\operatorname{deg}\left(h_{i} g_{i}\right)=\delta}\left(h_{i}-L T\left(h_{i}\right)\right) g_{i} \\
& +\sum_{\operatorname{deg}(f)<\operatorname{deg}\left(h_{i} g_{i}\right)<\delta} h_{i} g_{i}+\sum_{\operatorname{deg}\left(h_{i} g_{i}\right) \leq \operatorname{deg}(f)} h_{i} g_{i}
\end{aligned}
$$

where the first sum must have degree less than $\delta$ and the second sum subtracts $L T\left(h_{i}\right)$, thus has degree less than $\delta$.

Using Lemma 2 above, there exists $a_{i j}$ such that

$$
\begin{aligned}
\sum_{\operatorname{deg}\left(h_{i} g_{i}\right)=\delta} L T\left(h_{i}\right) g_{i} & =\sum_{i, j}^{d} a_{i j} S\left(L T\left(h_{i}\right) g_{i}, L T\left(h_{j}\right) g_{j}\right) \\
& =\sum_{i, j}^{d} a_{i j}\left(\frac{x^{\delta}}{L T\left(h_{i} g_{i}\right)} L T\left(h_{i}\right) g_{i}-\frac{x^{\delta}}{L T\left(h_{j} g_{j}\right)} L T\left(h_{j}\right) g_{j}\right) \\
& =\sum_{i, j}^{d} a_{i j}\left(\frac{x^{\delta}}{L T\left(h_{i}\right) L T\left(g_{i}\right)} L T\left(h_{i}\right) g_{i}-\frac{x^{\delta}}{L T\left(h_{j}\right) L T\left(g_{j}\right)} L T\left(h_{j}\right) g_{j}\right) \\
& =\sum_{i, j}^{d} a_{i j}\left(\frac{x^{\delta}}{x^{\alpha_{i j}}} \frac{x^{\alpha_{i j}} g_{i}}{L T\left(g_{i}\right)}-\frac{x^{\delta}}{\alpha_{i j}} \frac{x^{\alpha_{i j}} g_{j}}{L T\left(g_{j}\right)}\right) \\
& =\sum_{i, j}^{d} \frac{x^{\delta}}{x^{\alpha_{i j}}} a_{i j} S\left(g_{i}, g_{j}\right)
\end{aligned}
$$

where we use that $L T\left(L T\left(h_{k} g_{k}\right)\right)=L T\left(h_{k} g_{k}\right)$ and $L T\left(h_{k} g_{k}\right)=L T\left(h_{k}\right) L T\left(g_{k}\right)$ by wellordering, and $x^{\alpha_{i j}}=\operatorname{lcm}\left(L M\left(g_{i}\right) L M\left(g_{j}\right)\right)$.
Now by assumption, $\overline{S\left(g_{i}, g_{j}\right)}{ }^{G}=0$, so $S\left(g_{i}, g_{j}\right)=\sum_{k=1}^{d} p_{k} g_{k}$ with $\operatorname{deg}\left(p_{k} g_{k}\right)$ less than or equal to $\operatorname{deg}\left(S\left(g_{i}, g_{j}\right)\right)$, which is strictly less than $\alpha_{i j}$, where the first inequality follows from the division algorithm and the second since $S\left(g_{i}, g_{j}\right)$ cancels the leading term.

Consequently we have

$$
\begin{aligned}
\sum_{\operatorname{deg}\left(h_{i} g_{i}\right)=\delta} L T\left(h_{i}\right) g_{i} & =\sum_{i, j}^{d} \frac{x^{\delta}}{x^{\alpha_{i j}}} a_{i j} \sum_{k=1}^{d} p_{k} g_{k} \\
& =\sum_{i, j}^{d} \sum_{k=1}^{d} \frac{x^{\delta}}{x^{\alpha_{i j}}} p_{k} g_{k}
\end{aligned}
$$

with degree strictly less than $\delta$.
This concludes case 2 and the reverse direction.
$\left(\Rightarrow\right.$ :) We need to show that $\overline{S\left(g_{i}, g_{j}\right)}{ }^{G}=0$ for all $\left(g_{i}, g_{j}\right)$ in $G$ under the assumption that $G=\left\{g_{1}, \ldots, g_{d}\right\}$ is Gröbner.
$S\left(g_{i}, g_{j}\right)=\frac{x^{\gamma}}{L T\left(g_{i}\right)} \cdot g_{i}-\frac{x^{\gamma}}{L T\left(g_{j}\right)} \cdot g_{j}$ in $I$ implies $L T\left(S\left(g_{i}, g_{j}\right)\right)$ in $\langle L T(I)\rangle=\left\langle L T\left(g_{i}\right), \ldots, L T\left(g_{d}\right)\right\rangle$, so we can move to the next step of the division algorithm without remainder by dividing by some $L T\left(g_{k}\right)$ by Lemma 1 .

The next step is dividing $S\left(g_{i}, g_{j}\right)-\left(\frac{L T\left(S\left(g_{i}, g_{j}\right)\right)}{L T\left(g_{k}\right)} \cdot g_{k}\right)$ (in $\left.I\right)$ by $G$. By same reasoning as above, we can move to the next step.

The division algorithm terminates by well-ordering, thus we conclude that

$$
{\overline{S\left(g_{i}, g_{j}\right)}}^{G}=0 .
$$

This concludes the proof.

This constructs an algorithm (Buchberger's Algorithm) to construct a Gröbner basis from an original generating set $G$, by adding ${\overline{S\left(g_{i}, g_{j}\right)}}^{G} \neq 0$ to $G$, and repeating the process with this newly defined $G$ until $\overline{S\left(g_{i}, g_{j}\right)}=0$ for all $i, j \in\{1, \ldots, n\}$.

Proposition 2. The Buchberger Algorithm terminates.
Proof. Let $G=\left\{g_{1}, \ldots g_{d}\right\}$ be a potential Gröbner basis. If ${\overline{S\left(g_{i}, g_{j}\right)}}^{G} \neq 0$, let $G^{\prime}=G \cup\left\{{\overline{S\left(g_{i}, g_{j}\right)}}^{G}\right\}$. We note that

$$
S\left(g_{i}, g_{j}\right)=\sum_{i=1}^{d} p_{i} g_{i}+{\overline{S\left(g_{i}, g_{j}\right)}}^{G}
$$

implies ${\overline{S\left(g_{i}, g_{j}\right)}}^{G}$ is in I. So $G \subsetneq G^{\prime}$. We also have that $\langle L T(G)\rangle \subsetneq\left\langle L T\left(G^{\prime}\right)\right\rangle$, otherwise ${\overline{S\left(g_{i}, g_{j}\right)}}^{G}$ would not be the remainder. Suppose that the algorithm never terminates. This would imply that

$$
\langle L T(G)\rangle \subsetneq\left\langle L T\left(G^{\prime}\right)\right\rangle \subsetneq\left\langle L T\left(G^{\prime \prime}\right)\right\rangle \subsetneq \ldots
$$

which contradicts that $K\left[x_{1}, \ldots, x_{n}\right]$ is Noetherian.
Thus, the Buchberger Algorithm terminates.

Corollary 1. A Gröbner basis always exists.
Proof. Directly follows from the Buchburger Algorithm and Proposition 2.

We provide a working example of the construction of a Gröbner basis using grlex ordering.
Example. Let $G_{1}=\left\{x^{3}-2 x y, x^{2} y-2 y^{2}+x\right\}, f_{1}:=x^{3}-2 x y, f_{2}:=x^{2} y-2 y^{2}+x$, our potential Gröbner basis. We test this hypothesis in grlex:
$S\left(f_{1}, f_{2}\right)=y \cdot\left(x^{3}-2 x y\right)-x \cdot\left(x^{2} y-2 y^{2}+x\right)=-x^{2}:=f_{3}$
$\overline{S\left(f_{1}, f_{2}\right)}{ }^{G_{1}}=f_{3} \neq 0$ (we omit this division calculation since the remainder occurs immediately)
Since ${\overline{S\left(f_{i}, f_{j}\right)}}^{G_{1}}$ is non-zero for some $i, j, G_{1}$ is not a Gröbner basis.

Now let $G_{2}=\left\{f_{1}, f_{2}, f_{3}\right\}$.
$\overline{S\left(f_{1}, f_{2}\right)}{ }^{G_{2}}=0$ (due to above)
$S\left(f_{1}, f_{3}\right)=\left(x^{3}-2 x y\right)-(-x) \cdot\left(-x^{2}\right)=-2 x y:=f_{4}$
$\overline{S\left(f_{1}, f_{3}\right)}{ }^{G}=f_{4} \neq 0$ (we omit this division calculation since the remainder occurs immediately)
Similarly, since ${\overline{S\left(f_{i}, f_{j}\right)}}^{G_{2}}$ is non-zero for some $i, j, G_{2}$ is not a Gröbner basis.

Continuing, let $G_{3}=\left\{f_{1}, f_{2}, f_{3}, f_{4}\right\}$
${\overline{S\left(f_{1}, f_{2}\right)}}^{G_{3}}=0$ (as before)
${\overline{S\left(f_{1}, f_{3}\right)}}^{G_{3}}=0$ (due to above)
$S\left(f_{1}, f_{4}\right)=y \cdot\left(x^{3}-2 x y\right)-\left(-\frac{1}{2} x^{2}\right) \cdot(-2 x y)=-2 x y^{2}$
$\overline{S\left(f_{1}, f_{4}\right)^{G_{3}}}=0$ (we omit this division calculation since it is only one step)
$S\left(f_{2}, f_{3}\right)=\left(x^{2} y-2 y^{2}+x\right)-(-y)\left(-x^{2}\right)=-2 y^{2}+x:=f_{5}$
$\overline{S\left(f_{2}, f_{3}\right)}{ }^{G_{3}}=f_{5} \neq 0$ (we omit this division calculation since the remainder occurs immediately)

Let $G_{4}=\left\{f_{1}, f_{2}, f_{3}, f_{4}, f_{5}\right\}$
$\overline{S\left(f_{1}, f_{2}\right)}{ }^{G_{4}}={\overline{S\left(f_{1}, f_{3}\right)}}^{G_{4}}={\overline{S\left(f_{1}, f_{4}\right)}}^{G_{4}}=0$, as before.
By adding $f_{5}$ we have ${\overline{S\left(f_{2}, f_{3}\right)}}^{G_{4}}=0$.

It remains to check that
${\overline{S\left(f_{1}, f_{5}\right)}}^{G_{4}},{\overline{S\left(f_{2}, f_{4}\right)}}^{G_{4}},{\overline{S\left(f_{2}, f_{5}\right)}}^{G_{4}},{\overline{S\left(f_{3}, f_{4}\right)}}^{G_{4}},{\overline{S\left(f_{3}, f_{5}\right)}}^{G_{4}},{\overline{S\left(f_{4}, f_{5}\right)}}^{G_{4}}=0$.
Note that we do not need to check $S\left(f_{2}, f_{1}\right)$ etc. since $S\left(f_{i}, f_{j}\right)=-S\left(f_{j}, f_{i}\right), 1 \leq i<j \leq 5$, and -1 clearly does not affect divisibility as a unit.
$S\left(f_{1}, f_{5}\right)=y^{2} \cdot\left(x^{3}-2 x y\right)-\left(-\frac{1}{2} x^{3}\right) \cdot\left(-2 y^{2}+x\right)=-2 x y+\frac{1}{2} x^{4}$
We now compute $\overline{S\left(f_{1}, f_{5}\right)}{ }^{G_{4}}$ :

$$
\begin{aligned}
& \begin{array}{l}
q_{1}=\frac{1}{2} x \\
q_{2}=-1 \\
q_{3}=0 \\
q_{4}=1 \\
q_{5}=1
\end{array} \\
& \left.\begin{array}{l}
f_{1} \\
f_{2} \\
f_{3} \\
f_{4} \\
f_{5}
\end{array}\right)-2 x y+\frac{1}{2} x^{4} \\
& \frac{-(-2 x y)}{\frac{1}{2} x^{4}} \\
& \frac{-\left(\frac{1}{2} x^{4}-x^{2} y\right)}{-x^{2} y} \\
& \frac{-\quad\left(-x^{2} y+2 y^{2}-x\right)}{-2 y^{2}+x} \\
& \frac{-\quad\left(-2 y^{2}+x\right)}{0}
\end{aligned}
$$

Note that this is not the only way to solve $\overline{S\left(f_{1}, f_{5}\right)}{ }^{G}$. Since we are free to choose the order of divisors to compute the remainder, we see that we also get 0 remainder for

$$
\begin{array}{r}
\left.\begin{array}{r}
q_{4}=1 \\
q_{3}=-\frac{1}{2} x^{2} \\
\ldots \\
f_{4} \\
f_{3}
\end{array}\right)-2 x y+\frac{1}{2} x^{4} \\
-(-2 x y) \\
\frac{-\left(\frac{1}{2} x^{4}\right)^{\frac{1}{2} x^{4}}}{0}
\end{array}
$$

Similarly, one can show the remaining $\overline{S\left(f_{i}, f_{j}\right)}{ }^{G_{4}}=0$, so we conclude that $G_{4}$ is a Gröbner basis for $\left\langle x^{3}-2 x y, x^{2} y-2 y^{2}+x\right\rangle$ with respect to grlex.

## 5 Minimal and Reduced Gröbner Basis

We briefly state and discuss the concept of a minimal and reduced Gröbner basis.

Definition 11. Let $G=\left\{g_{i}, \ldots, g_{n}\right\}$ be a Gröbner basis. A minimal Gröbner basis is one such that all leading coefficients equal 1 and is free of any $g_{i}$ such that $L T\left(g_{i}\right)$ is in $\left\langle L T\left(G \backslash\left\{g_{i}\right\}\right)\right\rangle$.

Remark 3. $G \backslash\left\{g_{i}\right\}$ is still a Gröbner basis since $\langle L T(G)\rangle=\left\langle L T\left(G \backslash\left\{g_{i}\right\}\right)\right\rangle$.
An ideal $I$ does not have a unique minimal Gröbner basis, although given a minimal Gröbner basis $\mathrm{G},\langle L T(G)\rangle$ forms the unique minimal basis of $\langle L T(I)\rangle$ ([1]).

However, for $I \neq\{0\}$, $I$ does have a unique reduced Gröbner basis for a given monomial ordering ([1]):

Definition 12. A reduced Gröbner basis G is one such that all leading coefficients equal 1, and for any $g_{i}$ in $G$, no monomial of $g_{i}$ is an element of $\left\langle L T\left(G \backslash\left\{g_{i}\right\}\right)\right\rangle$.

## 6 The Ideal Membership Problem

How does the existence of a Gröbner basis solve the classical Ideal Membership problem introduced in 2?

Theorem 3. Let $G=\left\{g_{1}, \ldots, g_{d}\right\}$ be a Gröbner basis for an ideal $I$ of $K\left[x_{1}, \ldots, x_{n}\right]$. Then the remainder when dividing an element $f$ of $K\left[x_{1}, \ldots, x_{n}\right]$ by $G$ is independent of the ordering of $g_{i}$.

Proof. The division algorithm gives existence of such a remainder. To show uniqueness, let $f=p_{1} g_{1}+\ldots+p_{d} g_{d}+r_{1}=q_{1} g_{1}+\ldots+q_{d} g_{d}+r_{2}$ in $K\left[x_{1}, \ldots, x_{n}\right]$ have two distinct remainders. This gives that $0=p_{1} g_{1}+\ldots+p_{d} g_{d}-q_{1} g_{1}-\ldots-q_{d} g_{d}=r_{1}-r_{2}$ is in $I$. Thus, $L T\left(r_{1}-r_{2}\right)$ is in $\left\langle L T\left(g_{1}\right), \ldots, L T\left(g_{d}\right)\right\rangle$ by the definition of Gröbner basis, and by Lemma 1 some $L T\left(g_{i}\right)$ divides $L T\left(r_{1}-r_{2}\right)$. This is a contradiction since no terms of $r_{1}$ nor $r_{2}$ are divisible by any $L T\left(g_{i}\right)$. We conclude that $r_{1}-r_{2}$ must be zero.

Corollary 2. The remainder of $f$ after dividing by Gröbner basis of $I$ yields 0 if and only if $f$ is in $I$ itself.

Proof. This is a direct result from Theorem 3, but we write a full proof.
$\left(\Rightarrow\right.$ :) Remainder 0 gives $f=p_{1} g_{1}+\ldots+p_{d} g_{d}$ for $p_{i}$ in $K\left[x_{1}, \ldots, x_{n}\right]$ and $g_{i}$ in $G$, so $f$ is an element of $I$ as written as a linear combination of elements of $G$.
$(\Leftarrow:)$ (This is identical to Remark 1) We can write $f=p_{1} g_{1}+\ldots+p_{d} g_{d}+r$ for some remainder $r$ using the division algorithm. Since $f$ is in $I, r=f-p_{1} g_{1}+\ldots+p_{d} g_{d}$ is an element of $I$. Assume $r$ is non-zero. By the definition of Gröbner basis then, $L T(r)$ is in $\left\langle L T\left(g_{1}\right), \ldots, L T\left(g_{d}\right)\right\rangle$, and by Lemma 1 some $L T\left(g_{i}\right)$ divides $L T(r)$. But this contradicts that $r$ is the smallest remainder, thus $r$ must be zero.

## 7 Some Applications in Number Theory

Theoretically, a Gröbner basis exists for any ideal in a polynomial ring over any field $K$. In particular, this is true for algebraic number fields $\mathbb{Q}(\alpha)$, for algebraic number $\alpha$. The study of algebraic number fields (finite extensions of the field $\mathbb{Q}$ ) is central to algebraic number theory, and is introduced in an undergraduate course in number theory. As we have seen, Gröbner bases allow us to study ideals in polynomial rings over such fields.

Another classical application of Gröbner bases is in finding solutions to systems of polynomial equations. For a system of polynomial equations, we can translate this system into an equivalent system equating to zero. We can then study the ideal generated by polynomials and the affine variety $\mathbf{V}(I)$ :

Definition 13. For an ideal $I$ of $K\left[x_{1}, \ldots, x_{n}\right]$, we define

$$
\mathbf{V}(I)=\left\{\left(a_{1}, \ldots, a_{n}\right) \in K^{n} \mid f\left(a_{1}, \ldots, a_{n}\right)=0 \forall f \in I\right\}
$$

Proposition 3. ([1]) If $I=\left\langle f_{1}, \ldots f_{m}\right\rangle \triangleleft K\left[x_{1}, \ldots, x_{n}\right]$, then $\mathbf{V}(I)=\mathbf{V}\left(f_{1}, \ldots, f_{m}\right)$.
By Proposition 3, solutions are independent of generating basis, so we can equivalently study the Gröbner system of equations. Interestingly, Gröbner bases in lex ordering triangulate (eliminate variables) such systems and can simplify calculations greatly (cf. [1]).

In a 2015 paper ([4]) on the very topic of Gröbner bases over algebraic number fields, a more time efficient algorithm to compute Gröbner bases is proposed. It is well observed that the efficiency of computing Gröbner bases via the Buchberger algorithm is highly dependent on the field arithmetic, where algebraic number fields are relatively difficult.

In another paper ([5]) (2007), Gröbner bases are used to automate a variety of classic proofs of number theoretic nature involving divisibility, congruence, and the notion of coprime, by reduction to our original ideal membership problem.

A final example demonstrates the power of Gröbner bases by reducing computational time.
Definition 14. A Diophantine equation is a polynomial equation in two or more variables with integer coefficients, where integer solutions are of interest.
Definition 15. Pell's equation is a Diophantine equation of the form $x^{2}-A y^{2}=1$ where $A$ is a positive non-square integer. Again, integer solutions are of interest. This equation can be generalized to $x^{2}-A y^{2}=N$ where N is a nonzero integer.

Theorem 2. ([6]) For any integers $a>1$ and $b>1$ such that $a \neq b$, the system

$$
\begin{aligned}
& a x^{2}-c z^{2}=1 \\
& b y^{2}-d z^{2}=1
\end{aligned}
$$

of generalized Pell's equations has at most two integer solutions with $x, y, z>0$.

When showing a restriction on the set of solutions to a system of generalized Pell's equations, the 2008 paper ([6]) proposes a Gröbner basis method which requires half the classical method time.

## References

[1] Cox, David, Little, John, and O'Shea, Donal. Ideals, Varieties, and Algorithms an Introduction to Algebraic Geometry and Commutative Algebra. Fourth edition, Berlin, Heidelberg, Springer, 1998.
[2] Dummit, David, and Foote, Richard. Abstract Algebra. Third edition, John Wiley and Sons, Inc, 2004.
[3] "Gröbner basis". Wikipedia, The Free Encyclopedia, Wikimedia Foundation, 29 January 2023, https://en.wikipedia.org/wiki/Gröbner_basis.
[4] Boku, Dereje Kifle, Decker, Wolfram, Fieker, Claus, and Steenpass, Andreas (2015). "Gröbner Bases Over Algebraic Number Fields." https://arxiv.org/pdf/1504.04564. pdf.
[5] Harrison, John (2007). "Automating Elementary Number-Theoretic Proofs Using Gröbner Bases". In: Pfenning, F. (eds) Automated Deduction - CADE-21. CADE 2007. Lecture Notes in Computer Science(), vol 4603. Springer, Berlin, Heidelberg. https: //doi.org/10.1007/978-3-540-73595-3_5.
[6] Cipu, Mihai. "Gröbner Bases and Diophantine Analysis". Journal of Symbolic Computation, vol. 43, no. 10, October 2008, pp. 681-687. URL.

