# Regularisation of Gröbner basis computations for weighted and determinantal systems, and an application to medical imagery 

Thibaut Verron

## To cite this version:

Thibaut Verron. Regularisation of Gröbner basis computations for weighted and determinantal systems, and an application to medical imagery. Symbolic Computation [cs.SC]. Université Pierre et Marie Curie, 2016. English. NNT: . tel-01404406v1

## HAL Id: tel-01404406 https://theses.hal.science/tel-01404406v1

Submitted on 28 Nov 2016 (v1), last revised 1 Mar 2017 (v2)

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers.

L'archive ouverte pluridisciplinaire HAL, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

THÊSE DE DOCTORAT DE L'UNIVERSITÉ PIERRE ET MARIE CURIE

Spécialité

Informatique
École Doctorale Informatique, Télécommunications et Électronique (Paris)

Présentée par
Thibaut Verron

Pour obtenir le grade de
DOCTEUR de l'UNIVERSITÉ PIERRE ET MARIE CURIE

# Régularisation du calcul de bases de Gröbner pour des systèmes avec poids et déterminantiels, et application en imagerie médicale 

Thèse dirigée par Jean-Charles Faugère et Mohab Safey El Din
soutenue le lundi 26 septembre 2016
après avis des rapporteurs:
M. Laurent Busé Chargé de recherche, Inria, HDR
M. Bruno SAlvy

Directeur de recherche, Inria
devant le jury composé de :
M. Jean-Charles Faugère Directeur de recherche, Inria
M. Mohab Safey El Din Professeur, Université Pierre et Marie Curie
M. Laurent Busé Chargé de recherche, Inria, HDR
M. Bruno Salvy Directeur de recherche, Inria
M. Bernard Bonnard Professeur, Université de Bourgogne
M. Stef Graillat

Professeur, Université Pierre et Marie Curie

## Remerciements

Je ne saurais commencer sans remercier Jean-Charles et Mohab d'avoir accepté de m'encadrer lors de cette thèse. Cette thèse n'aurait pas été la même sans leur suivi précis, leur disponibilité constante, leurs conseils judicieux et leurs relectures toujours très attentives. Mention spéciale aux "blagues" de Mohab, toujours rafraîchissantes, surtout en période de deadline!

Je remercie chaleureusement Laurent Busé et Bruno Salvy pour avoir accepté de rapporter ma thèse, leurs conseils et leurs références ont été très utiles. Merci aussi à Bernard Bonnard et Stef Graillat pour avoir accepté de faire partie du jury.

Je tiens aussi à remercier Bernard pour son aide dans la recherche d'opportunités postdoctorales, et Alain Jacquemard pour ses explications patientes de théorie du contrôle, ainsi que ses conseils toujours très avisés.

La période de la thèse est également l'occasion de goûter à l'expérience de l'enseignement, et je remercie ceux qui m'ont accordé leur confiance pour cette tâche : Valérie Ménissier-Morain, Christoph Lauter, Camélia Constantin, Karine Heydemann, et Fabien Tarissan ; ainsi que mes partenaires lors de ces enseignements : Olga Kupriianova, Lionel Tabourier, Alexandre Wallet, Claude Dutheillet et Dominique Béréziat. Cette expérience a été très enrichissante, et c'est grâce à vous.

Je voudrais aussi remercier tout le personnel administratif que j'ai dû solliciter pendant ma thèse, leur dévouement à nous faciliter le quotidien a été à chaque fois d'une grande aide. J'ai une pensée particulière pour Isabelle Delais à l'ENS et Irphane Khan à l'UPMC, pour l'incroyable aide qu'ils ont pu m'apporter quand j'en avais besoin.

Merci également à toute l'équipe PolSys pour leur accueil toutes ces années : Daniel, et ces longues conversations toujours très enrichissantes, Élias, jamais à court d'encouragements, Guénaël pour toutes ces blagues qui ont égayé le bureau, Jérémy et ses talents de batteur, et Ludovic qui fut mon premier contact avec l'équipe. Merci également aux post-doctorants de l'équipe que j'ai connus, Christian et Brice.

Durant ces années, j'ai partagé le bureau 338 avec de nombreuses personnes, et je les remercie tous pour ces moments partagés : tout d'abord les doctorants, (par ordre alphabétique) Alexandre G., Alexandre W., Aurélien, Cécile, Christopher, Frédéric, Ivan, Jules, Louise, Matias, Mourad, Nastia, Pacôme, Pierre-Jean, Rina, Ulrick et Vincent; ainsi que nos cobureaux de passage, Aaron, Anna, Marta, Nabil et Simone. Merci également à Luk, toujours présent dans nos cours et nos .emacs.d. Je tiens également à remercier nos voisins, stagiaires ou doctorants dans d’autres équipes: Charles, Clothilde, Daisuke, Guillaume, Jérôme C., Jérôme G., Ramon, Romain, Sénan et Vincent G. Enfin, merci aux amis rencontrés au détour de conférences, notamment Tristan et son expertise en culture nippone, et Yacine et sa chemise rose.

La thèse est une période difficile, et toute distraction permettant d'alléger les tensions est la bienvenue. C'est aussi une période où l'on perd de vue beaucoup d'amis, et je remercie particulièrement ceux qui m'ont fait la faveur de s'accrocher : notamment Clémentine, pour les
après-midi jeux et concerts, et Fabien, soutien indéfectible depuis cette rentrée en $6^{e}$, il y a de nombreuses années. Je n'oublie pas non plus tous les autres! Je remercie également toute la troupe de Disharmonies pour toutes ces réunions riches en éclats de rire, et je ne leur tiens pas rigueur pour les soirées passées à maquetter en échange.
La famille est également un soutien important, et sur lequel il est bon de savoir que l'on peut compter. Je remercie donc mes parents et Quentin pour leur confiance, leurs encouragements et leur patience. Je ne serais certainement pas là sans eux.
Enfin, je voudrais remercier du fond du cœur Caroline, pour son soutien et ses encouragements sur lesquels je pouvais toujours compter, jusque dans les moments les plus difficiles. Merci.

## Contents

Introduction ..... 9

1. Motivation ..... 9
2. State of the art ..... 14
3. Contributions ..... 16
4. Perspectives ..... 22
I. Preliminaries ..... 25
5. Algebra and geometry ..... 27
1.1. Commutative algebra ..... 27
1.1.1. Ideals ..... 27
1.1.2. Graded rings ..... 28
1.1.3. Polynomial algebras ..... 28
1.1.4. Combinatorics of monomials ..... 30
1.1.5. Hilbert series ..... 32
1.2. Regularity properties ..... 34
1.2.1. Regular sequences ..... 34
1.2.2. Description of the Hilbert series of a homogeneous regular sequence ..... 35
1.2.3. Noether position ..... 36
1.2.4. Semi-regular sequences ..... 39
1.3. Algebraic geometry ..... 40
1.3.1. Nullstellensatz ..... 40
1.3.2. Zariski topology ..... 40
1.3.3. Morphisms and rational maps ..... 41
1.3.4. Dimension and degree ..... 42
1.3.5. Geometric interpretation of the regularity properties ..... 43
1.3.6. Singularities ..... 45
1.3.7. Real semi-algebraic geometry ..... 46
1.4. Genericity ..... 48
1.4.1. Definition ..... 48
1.4.2. Generic properties of homogeneous systems ..... 49
1.4.3. Generic changes of coordinates ..... 50
1.5. Determinantal varieties ..... 51
1.5.1. Definition ..... 52
1.5.2. Cramer's formula and Schur's complement ..... 54
1.5.3. Incidence varieties ..... 56
6. Gröbner bases ..... 59
2.1. Monomial orderings ..... 59
2.1.1. Definition ..... 59
2.1.2. Lexicographical ordering ..... 60
2.1.3. Graded reverse-lexicographical ordering ..... 61
2.1.4. Elimination orders ..... 63
2.2. Gröbner bases: definition ..... 63
2.3. Applications ..... 65
2.3.1. Zero-dimensional systems ..... 66
2.3.2. Positive-dimensional systems and eliminations ..... 66
2.4. Algorithms ..... 67
2.4.1. A pairs algorithm: Buchberger's algorithm ..... 67
2.4.2. A matrix algorithm: Matrix-F5 ..... 69
2.4.3. Matrix algorithms in the affine case ..... 73
2.4.4. Change of order: FGLM ..... 75
2.5. Complexity results ..... 77
2.5.1. Complexity model and notations ..... 77
2.5.2. Complexity of Matrix- $F_{5}$ and degree of regularity ..... 78
2.5.3. Thin-grained complexity of Matrix- $F_{5}$ ..... 80
2.5.4. Complexity of FGLM ..... 82
II. Contributions ..... 85
7. Weighted homogeneous systems ..... 87
3.1. Introduction ..... 87
3.2. Properties ..... 92
3.2.1. Definitions and properties ..... 92
3.2.2. Degree and Bézout's bound ..... 93
3.2.3. Changes of variables and reverse chain-divisible systems of weights ..... 95
3.2.4. Genericity of regularity properties and $W$-compatibility ..... 97
3.2.5. Characterization of the Hilbert series ..... 100
3.3. $\quad W$-Degree of regularity of regular sequences ..... 104
3.3.1. Macaulay's bound in dimension zero ..... 104
3.3.2. Sharper bound on the degree of regularity ..... 106
3.3.3. Conjectured exact formula ..... 107
3.3.4. Positive-dimensional regular sequences ..... 108
3.4. Overdetermined systems ..... 108
3.4.1. Definition of semi-regularity ..... 108
3.4.2. Characterization with the Hilbert series ..... 110
3.4.3. $W$-degree of regularity of semi-regular sequences ..... 114
3.4.4. Fröberg's conjecture ..... 116
3.5. Algorithms and complexity ..... 117
3.5.1. Critical pairs algorithms ..... 117
3.5.2. Adapting the Matrix- $F_{5}$ algorithm ..... 118
3.5.3. Thin-grained complexity of Matrix- $F_{5}$ ..... 120
3.5.4. FGLM and computational strategy for zero-dimensional systems ..... 122
3.5.5. Other algorithms ..... 123
3.6. Experiments ..... 123
3.6.1. Generic systems ..... 123
3.6.2. Discrete Logarithm Problem ..... 126
3.6.3. Polynomial inversion ..... 127
8. Real roots classification for determinants - Applic. to contrast optimization ..... 133
4.1. Introduction ..... 133
4.2. Modeling the dynamics of the contrast optimization ..... 137
4.3. Algorithm ..... 138
4.3.1. Classification strategy ..... 138
4.3.2. The determinantal problem ..... 139
4.3.3. Incidence varieties ..... 141
4.3.4. Locus of rank exactly $r_{0}$ ..... 143
4.3.5. Singularities ..... 144
4.3.6. Boundary ..... 144
4.4. The contrast problem ..... 145
4.4.1. The case of water ..... 145
4.4.2. The general case ..... 150
III. Appendices ..... 151
A. Source code ..... 153
A.1. Magma source code for Matrix- $\mathrm{F}_{5}$ ..... 153
A.1.1. Algorithm Matrix- $\mathrm{F}_{5}$ ..... 153
A.1.2. Weighted homogeneous systems ..... 166
A.2. Maple source code for determinantal varieties (Chapter 4) ..... 170
A.2.1. Algorithms ..... 170
A.2.2. Case of water ..... 177
A.2.3. General case ..... 184
B. Bibliography ..... 187
Index ..... 197

## Introduction

## 1. Motivation

### 1.1. Polynomial system solving and applications

The problem of solving polynomial systems is very important nowadays, because such systems arise in a wide range of applications. Indeed, polynomials are a basis for models in many areas of research.

We give three examples of general applications, both practical and theoretical, of polynomial system solving that we will encounter in this thesis.

Cryptography Algebraic cryptanalysis models cryptographical systems using polynomial systems, and reduces the problem of breaking the code to solving these systems. We shall encounter an example related to the Discrete Logarithm Problem (DLP) on elliptic curves.

Polynomial inversion The problem is, given some polynomials $f_{1}, \ldots, f_{m}$, to identify the relations between the $f_{i}$ 's. For example, the only nontrivial relation between $f_{1}=X Y, f_{2}=X^{2}$ and $f_{3}=Y^{2}$ is $f_{1}^{2}=f_{2} f_{3}$. It can be used for example to obtain implicit equations describing a parametric object.

Real roots classification This is also related to parametric equations: the input is a polynomial system $F=\left(f_{1}(\mathbf{X}, \mathbf{T}), \ldots, f_{m}(\mathbf{X}, \mathbf{T})\right)$ with coefficients in $\mathbb{R}$, indeterminates $\mathbf{X}=\left(X_{1}, \ldots, X_{n}\right)$ and parameters $\mathbf{T}=\left(T_{1}, \ldots, T_{s}\right)$ such that for any $\mathbf{t} \in \mathbb{R}^{s}$, the equations $\left(f_{1}(\mathbf{X}, \mathbf{t}), \ldots, f_{m}(\mathbf{X}, \mathbf{t})\right)$ with unknowns $\mathbf{X}=\left(X_{1}, \ldots, X_{n}\right)$ have a finite number of real solutions $C_{t}$. The goal is to find a dense covering of $\mathbb{R}^{s}$ with open subsets over which $C_{t}$ is constant.

These examples illustrate the fact that the meaning of solving a polynomial system depends on the application. The simplest question one might ask given a polynomial system is "does it have a solution?". This problem is NP-complete on finite fields. Solving a system usually means computing some more information about the solution set. The exact form of the question depends on the application, and systems are usually split in two categories: they can be zerodimensional or positive-dimensional.
Zero-dimensional systems are systems with a finite number of solutions, in an algebraic closure of the coefficient field. For example, systems arising in applications to cryptography are usually zero-dimensional: there are enough constraints to ensure that there is only one solution to the problem. In this case, "solving" usually means either listing all the solutions of the system, giving a parametrisation of all solutions of the system, or simply giving one solution of the system.

On the other hand, positive-dimensional systems are systems with infinitely many solutions, again in some algebraic closure of the coefficient field. There, "solving" potentially takes on many
different meanings, depending on the system and the application: finding a parametrisation of the solutions, computing the equations of the projection of the solutions on a subspace (or in other words, eliminate variables from the system), computing the dimension of the solution set (that is identifying a subset of unknowns which are independent), or, when applicable given the coefficient field, finding topological information on the set of solutions (for example one point per connected component of the set or its complementary).
For example, the polynomial inversion problem is an elimination problem: if the input polynomials are $\left(f_{1}, \ldots, f_{m}\right) \in \mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$, one can construct the ideal $I=\left\langle F_{1}-f_{1}, \ldots, F_{m}\right.$ $\left.f_{m}\right\rangle$ with extra variables $F_{1}, \ldots, F_{m}$, and the wanted relations shall correspond to polynomials in $I \cap \mathbb{K}\left[F_{1}, \ldots, F_{m}\right]$, thus eliminating $X_{1}, \ldots, X_{n}$.
The real roots classification problem also involves an elimination problem: we need to find the set of values of the parameters where the number of real solutions in X changes. Let $\pi: \mathbb{R}^{n} \times \mathbb{R}^{t} \longrightarrow \mathbb{R}^{s}$ be the projection onto the parameters $T_{1}, \ldots, T_{s}$, the wanted subdivision of the parameter space is contained in the projection of the singular locus of the zeroes of $F$ and of the critical points of $\pi$ restricted to this variety. It is enough to compute one polynomial whose zeroes cover this set, and this too can be done using elimination. In order to complete the resolution of the problem, we then need to compute one point per connected component of the complementary of this set, and for each point, count the solutions in X .
Many tools have been developed for solving polynomial systems, including numerical methods (Newton's method...) or symbolic-numeric methods (homotopy continuation [SWo5]...), computing solutions with arbitrary precision, as well as symbolic methods from computer algebra (triangular sets [ALM99], multivariate resultants [CDo5], geometric resolution [GLSo1]...) computing exact solutions.

In this thesis, we mainly focus on Gröbner bases [Buc76]. They can be used to solve systems with a finite number of solutions as well as eliminate variables, thus covering most needs for polynomial system solving in applications. Furthermore, they can solve a polynomial system regardless of its field of definition.
We shall also make use of regular chains [LMXo5], as an alternative tool for computing polynomial elimination, and Cylindric Algebraic Decompositions [Col75] in order to obtain topological information about sets of real solutions of a system of polynomial equations and inequations.

### 1.2. Regularity of Gröbner basis computations

Gröbner bases are computed by successive reductions of the polynomials defining the system. This can be done in an algebraic way, similar to euclidean reductions in the univariate case, or using linear algebra on the coefficients of the polynomials. The historical algorithm by Buchberger [Buc76] uses the former method, while modern algorithms such as $F_{4}$ [Fau99] or $\mathrm{F}_{5}$ [Fauo2] use linear algebra to group reductions together. This makes the choice of a reduction strategy important, to condition how, at any given step, to pick polynomials for the reduction.

The default strategy is the normal strategy for homogeneous systems (or its inhomogeneous variant, the sugar strategy), which picks at each step the polynomials leading to the smallest degree reductions. For this strategy, the degree of the reductions can be seen as an indicator of progress, and we can characterize regular behaviors for the algorithms: a behavior is regular if
it does show any degree fall, that is a step where the degree of the reductions is not (strictly) increasing. See for example the curves on page 12: the red curve is an example of a regular behavior, whereas the blue curve is an example of an irregular behavior.

We shall see that regular behaviors make it possible to analyze the complexity of the algorithms: indeed, if at each step the degree is strictly increasing, we may bound the number of steps by bounding the maximal degree, and then obtain complexity bounds by bounding the cost of each step.

This can be done under some generic properties, that is algebraic properties which are satisfied by almost all systems. An example of a generic property relevant for Gröbner basis computations is zero-dimensional regular sequences, that is sequences of $n$ polynomials in $n$ variables having only a finite number of solutions (in an algebraic closure of the coefficient field). For homogeneous zero-dimensional regular sequences, or sequences such that the highest degree homogeneous components form a zero-dimensional regular sequence, algorithm $F_{5}$ has a regular behavior, and algorithm $F_{4}$ has as few degree falls as possible.

### 1.3. What to do when systems are not generic?

Problems appear when systems fail to satisfy genericity properties: complexity bounds relying on the regular behavior of the algorithms are no longer valid, and the algorithms waste time computing useless reductions.
For example, the authors of [Fau+13], when working with the DLP on Edwards elliptic curves, were confronted with the following system of 5 equations in 5 unknowns, over a finite field:

$$
\begin{aligned}
& 0=\left[\begin{array}{c}
7871 \\
18574 \\
14294 \\
32775 \\
20289
\end{array}\right] e_{5}^{16}+\left[\begin{array}{l}
53362 \\
50900 \\
36407 \\
58813 \\
20802
\end{array}\right] \tilde{e}_{1}^{8}+\left[\begin{array}{c}
26257 \\
128 \\
3037 \\
38424 \\
41456
\end{array}\right] \tilde{e}_{1}^{7} \tilde{e}_{2}+\left[\begin{array}{l}
25203 \\
23117 \\
28918 \\
29298 \\
56353
\end{array}\right] \tilde{e}_{1}^{6} \tilde{e}_{2}^{2}+\left[\begin{array}{l}
19817 \\
29737 \\
52187 \\
36574 \\
46683
\end{array}\right] \tilde{e}_{1}^{5} \tilde{e}_{2}^{3}+\left[\begin{array}{c}
9843 \\
3752 \\
27006 \\
64195 \\
63059
\end{array}\right] \\
&+\left[\begin{array}{l}
11204 \\
25459 \\
58263 \\
17964 \\
57146
\end{array}\right] \tilde{e}_{1}^{4} \tilde{e}_{2}^{4}
\end{aligned}
$$

This system is not a zero-dimensional regular sequence: the highest degree component of all 5 polynomials is a monomial in $e_{5}^{16}$, so any point with $e_{5}=0$ is a root of these components. And indeed, the behavior of Gröbner basis algorithms using the normal strategy for this system is irregular: see on Figure 1 the blue plot, representing the progress of algorithm $F_{5}$, step by step. Each step at which the plot stalls or reverses course is an irregularity, making it harder to bound the complexity.

On the other hand, assume that we now assign degree 2 to variables $\tilde{e}_{i}, i \in\{1, \ldots, 4\}$ (thus excluding $e_{5}$ ). In other words, we are using the weighted degree for the system of weights $(2,2,2,2,1)$. Then all monomials printed above have weighted-degree 16, together with 486 more monomials. The highest weighted homogeneous components of polynomials of the system


Figure 1: Step-by-step progress of Algorithm $\mathrm{F}_{5}$ on a DLP system
are now much larger than just one monomial, and they actually form a zero-dimensional regular sequence. By evaluating the weighted degree of polynomials for the reduction strategy, instead of the total degree, we recover a regular behavior. The behavior of algorithm $F_{5}$ using the weighted degree was plotted in red on Figure 1, and indeed the algorithm progresses steadily towards its final weighted degree.

Another example is given by determinantal ideals: a determinantal system is a system defined as the set of all $r$-minors of a matrix with polynomial entries. It encodes the fact that the matrix has rank at most $r-1$. These systems appear in a lot of applications, for example in cryptography (the MinRank problem [FLPo8]) or in convex optimization (semi-definite programming [Ott+15]). Computing critical points of projections restricted to a determinantal variety is a frequent problem, which has been thoroughly studied ([BV88]...).

For example, consider a $k \times k$ matrix $M$ with coefficients in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$, and let $V$ be the variety defined by $\operatorname{det}(M)$ (that is the $k-1$ 'th determinantal variety associated with $M$ ). Assume that we want to compute the critical points of the projection onto the first variable $X_{1}$, restricted to $V$.

A natural way of approaching this problem is to use the Jacobian criterion on the determinantal system: compute the partial derivatives $\frac{\partial \operatorname{det}(M)}{\partial X_{i}}, i \in\{2, \ldots, n\}$, saturate by $\frac{\partial \operatorname{det}(M)}{\partial X_{1}}$, and compute a Gröbner basis of this system. In the same way as before, we plot the behavior of this strategy in blue in Figure 2, for $k=5, n=7$ and $\mathbb{K}=\mathbb{F}_{65521}$. Again, it appears that there are some irregularities.

However, it is possible to change the modelization, or in other words find another set of equations describing the same geometrical object. The idea is that the matrix $M$ has rank at most $r-1$ if and only if its kernel has dimension at least $k-r+1$. This can be modelled by adding $k(k-r+1)$ variables to the polynomial algebra, form a $k \times(k-r+1)$-matrix $N$ with these variables as entries, and use the entries of $M \cdot N$ as the polynomial systems. We want $N$


Figure 2: Step-by-step progress of Algorithm $\mathrm{F}_{4}$ on the singularities of a determinantal system
to have full rank $(k-r+1)$, this can be described locally with $(k-r+1)^{2}$ additional equations.
The resulting system $F_{I}$ defines what is called an incidence variety. It models the same geometrical object, and we can compute the critical points of the projection restricted to this variety using the Jacobian criterion: critical points of the projection onto $X_{1}$ are points where the truncated Jacobian matrix $\mathrm{Jac}_{X_{2}, \ldots, X_{n}, Y_{1}, \ldots, Y_{k(k-r+1)}}\left(F_{I}\right)$ has rank less than the codimension of the variety.
This is a determinantal condition too, and we can model it with minors of the Jacobian, or again form an incidence variety using the Jacobian. This last approach is known in optimization as the introduction of Lagrange multipliers. In this case, both these approaches make the computations more regular, and faster by a factor of almost 10 . On Figure 2, we plotted in red (resp. green) the behavior of the computations on the minors of the Jacobian of the incidence system (resp. the system obtained by multiplying the Jacobian with Lagrange multipliers), with the same input as before.

Both weighted-homogeneous and determinantal systems are examples of the more general problem of solving structured systems: systems arising in applications frequently exhibit some structure that they do not share with arbitrary systems. This is both a problem and an advantage: on the one hand, this makes the systems non-generic, and usually the algorithms need to be adapted to recover a regular behavior for these systems; and on the other hand, once dedicated algorithms exist, the additional structure frequently makes it easier to solve generic structured systems than arbitrary generic systems.
In general, working with a new structure raises two deeply interleaved problems:

- identify properties of generic polynomials with the structure, which imply a regular behavior of dedicated algorithms;
- and design dedicated algorithms taking the structure into account and having a regular behavior for generic systems.


Figure 3: Algorithmic strategy for zero-dimensional systems

## 2. State of the art

### 2.1. Computational strategy for Gröbner bases

We now give more details about how Gröbner bases can be used for solving polynomial systems, and why a reduction strategy taking into account the degree of the polynomials is useful in this workflow.

The definition of a Gröbner basis depends on the choice of a monomial ordering, in order to give a formal meaning to the reduction of polynomials. Gröbner bases for different monomial orderings will yield different kinds of information, and be more or less easy to compute. In terms of actually solving the system, two monomial orderings are particularly useful:

- elimination orderings, that is orderings such that any monomial divisible by $X_{1}, \ldots, X_{i}$ is greater than any monomial not divisible by $X_{1}, \ldots, X_{i}$, for some $i$ : using such an ordering, one may eliminate the variables $X_{1}, \ldots, X_{i}$ from the ideal;
- the lexicographical ordering: up to a random change of coordinates, the lexicographical Gröbner basis of a zero-dimensional system is a parameterized representation of all solutions: $X_{n}$ is given as zero of a univariate polynomial, and all other unknowns are expressed as polynomials in $X_{n}$.

On the other hand, computing bases for elimination orderings or the lexicographical ordering is usually more difficult than for other orderings such as the GRevLex ordering. A common strategy for solving a polynomial system is thus to first compute a GREVLEX basis using algorithm $F_{4}$ or $F_{5}$, and then perform a change of ordering to compute an elimination or lexicographical basis from the GRevLex basis. This change of ordering can be done using another run of algorithm $\mathrm{F}_{4}$ or $\mathrm{F}_{5}$, or using a dedicated algorithm such as FGLM [Fau+93] in the zero-dimensional case (Figure 3), or the Gröbner walk [CKM97] in the general case.

This makes computing a Gröbner basis for the GRevLex ordering an important step in most cases. This order orders the monomials according to their degree first, which is why algorithms usually use the normal strategy by default when computing a GRevLex basis: it is a way to consider the smallest monomials first.

### 2.2. Complexity studies for homogeneous systems

Homogeneous systems are maybe the most widely studied structure for Gröbner basis computations: indeed, from the very start, strategies such as the normal strategy or the sugar strategy were designed in order to take into account the degree of homogeneous polynomials, or properties of their highest degree components when the systems were inhomogeneous.

It has been observed that the normal strategy frequently gives good results in practice. This does not necessarily mean that the algorithms are lightning-fast, but monitoring their progress does not show any "irregular" behavior of the kind of our earlier examples: the degree of the polynomials considered grows steadily from start to end.

The study of this behavior led to algebraic characterizations of systems exhibiting this regular behavior. It turned out that these properties are generic, which is a formal way of expressing that "most" systems will satisfy them, and explains why these strategies appeared efficient in practice.

An example of generic property ensuring regularity in Gröbner basis algorithms is regular sequences. It has several equivalent definitions, but for now we will restrict to systems with $n$ equations and $n$ unknowns (so-called square systems). A square system is a regular sequence if and only if its set of solutions is finite.

This property ensures that algorithm $\mathrm{F}_{5}$ has a regular behavior on homogeneous systems. For inhomogeneous systems, the "good" property is regularity in the affine sense, and it requires that the highest degree components of the system form a regular sequence.

Under these genericity hypotheses, the regular behavior of the algorithms makes it possible to obtain complexity bounds. For example, here are complexity bounds for algorithms $F_{5}$ and FGLM for generic square systems.

Theorem 1. Let $F=\left(f_{1}, \ldots, f_{n}\right) \subset \mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$ be a homogeneous system with respective degree $d_{1}, \ldots, d_{n}$. Assume that the system is regular in the affine sense (and thus zero-dimensional).

Then algorithm $F_{5}$ computes a GREvLEX Gröbner basis of $\langle F\rangle$ in time

$$
O\left(d_{\mathrm{reg}}\binom{n+d_{\mathrm{reg}}-1}{d_{\mathrm{reg}}}^{3}\right)
$$

where $d_{\text {reg }}$ is the degree of regularity of the system, that is the largest degree we need to reach in a run of $F_{5}$. It is bounded by Macaulay's bound [BFS14, Cor. 13]:

$$
d_{\mathrm{reg}} \leq \sum_{i=1}^{n}\left(d_{i}-1\right)+1
$$

Algorithm FGLM computes a lexicographical basis of $\langle F\rangle$ in time

$$
O\left(n \operatorname{deg}(I)^{3}\right)
$$

where $\operatorname{deg}(I)$ is the degree of the ideal, that is the number of solutions. It is given by Bézout's bound [Laz83, Prop. 1]:

$$
\operatorname{deg}(I)=\prod_{i=1}^{n} d_{i}
$$

### 2.3. Other structures

The problem of computing Gröbner bases for structured polynomial systems is not new, and a lot of new structures have been successfully exploited from this point of view in the past few
years: examples include multi-homogeneous systems [FSS11], systems invariant under a group action [FS12], determinantal systems [FSS13]...

In particular, this last example will prove useful in this thesis. Recall that given a $k \times k$-matrix $M$ with polynomial entries, and an integer $r \in\{1, \ldots, k\}$, the determinantal system $D_{r}$ is the system given by all $r$-minors of $M$. Its zeroes, the so-called determinantal variety, are points at which the matrix $M$ has rank at most $r-1$.

The structure of these ideals has been thoroughly studied, both from an algebraic standpoint [BV88] and from a computational standpoint [FSS13]. In particular, it is known that generically, the singular locus of $D_{r}$ is the determinantal variety $D_{r-1}$.

As for computations, the idea of modelling the determinantal variety as the projection of an incidence variety is far from new: for example, the classic method of using Lagrange multipliers in optimization relies on modelling rank defects of the Jacobian matrix of a system as an incidence variety, and this technique was central in the work of [ $\mathrm{Nal}_{15}$ ] on real optimization for determinantal varieties. This strategy has been analyzed from a computational point of view in the zero-dimensional case [FSS13].

## 3. Contributions

In this thesis, we worked on two classes of structured polynomial systems: weighted homogeneous systems, and determinantal systems for a root classification problem.

Weighted homogeneity is the structure used in the DLP example on page 11. It generalizes homogeneity, by giving more choice over how the degree is computed. This structure appears naturally in a wide range of applications, for example whenever one applies a polynomial change of variables to a homogeneous system. Algorithmic strategies for weighted homogeneous systems were known, but there was no complexity analysis. Furthermore, the FGLM step of the classic strategy became a bottleneck in the resolution. We identified genericity properties which allowed us to obtain complexity bounds for the algorithms, and described a strategy allowing algorithm FGLM to take advantage of the structure. Under genericity hypotheses, we proved that the algorithms have a completely regular behavior in this strategy.

On the other hand, for determinantal systems, we started with a specific problem from an application to contrast optimization, in medical imagery. While solving this problem, we identified that the underlying problem was the more general problem of real root classification, specialized to singularities of determinantal varieties. However, while implemented algorithms exist for this general problem, they were unable to give a solution to the system from contrast optimization. We proposed a strategy taking advantage of the determinantal structure of the system, in order to spread the computations over several smaller problems within the reach of existing algorithms. Using incidence varieties to model the rank defects, the regularity of the Gröbner basis algorithms used for computing eliminations is improved.

### 3.1. Weighted-homogeneous systems

Computational strategy Weighted homogeneous systems are a generalization of homogeneous systems. Given a system of weights $W=\left(w_{1}, \ldots, w_{n}\right) \in \mathbb{N}_{>0}^{n}$, the $W$-degree of the
monomial $X_{1}^{\alpha_{1}} \cdots X_{n}^{\alpha_{n}}$ is defined as

$$
\operatorname{deg}_{W}\left(X_{1}^{\alpha_{1}} \cdots X_{n}^{\alpha_{n}}\right)=\sum_{i=1}^{n} w_{i} \alpha_{i}
$$

A polynomial is said to be $W$-homogeneous if all its monomials have the same $W$-degree.
Equivalently, a polynomial $f$ is $W$-homogeneous with $W$-degree $d$ if and only if

$$
\operatorname{hom}_{W}(f)=f\left(X_{1}^{w_{1}}, \ldots, X_{n}^{w_{n}}\right)
$$

is homogeneous with degree $d$.
This morphism gives a natural strategy for computing Gröbner bases for a weighted homogeneous system $F$ : we can run the usual algorithms on hom $_{W}(F)$. Some algorithms, such as $\mathrm{F}_{4}$ or $F_{5}$, appears to become faster with this strategy, but algorithm FGLM becomes a bottleneck.

This strategy is also interesting from a theoretical point of view. Homogeneity is certainly the most studied structure with respect to Gröbner basis computations: this includes the complexity studies given above, but also detailed characterizations of the Hilbert series of an homogeneous ideal [Mor96], thin-grained complexity bounds under stronger genericity hypotheses [BFS14], and complexity analyses for overdetermined systems under a property conjectured to be generic [BFSo4].
The question raised by all these results for homogeneous systems is how far these bounds and results can be transposed to the weighted homogeneous case.

It turns out that for a generic $W$-homogeneous system $F, \operatorname{hom}_{W}(F)$ is "sufficiently generic" as a homogeneous system, and a run of Algorithm Matrix- $\mathrm{F}_{5}$ on $\operatorname{hom}_{W}(F)$ computes a Gröbner basis without any reduction to zero. In this case, as in the homogeneous case, its complexity can be determined in terms of the size of the computed matrices and the degree of regularity of the system.

In applications, systems are rarely weighted homogeneous. In this situation, as in the homogeneous case, the behavior of the algorithms is tied to reductions to zero of the highest $W$-degree components of the systems, and so, to whether these components satisfy regularity properties. If this condition is not satisfied, the behavior of algorithm $F_{5}$ is expected to be irregular, with reductions to zero and degree falls, and generic complexity analyses no longer apply: the algorithms perform worse in practice, and theory fails to quantify exactly how much worse.
This also explains why transforming a weighted homogeneous system into a homogeneous one yields significant improvements in practice. Indeed, given a generic weighted homogeneous systems, its highest degree components are unlikely to be large, and thus to satisfy regularity properties: informally, the smaller a class of polynomial systems is, the less likely there will be enough room for interesting generic properties to be satisfied.

We now turn to algorithm FGLM, for zero-dimensional systems: its complexity is bounded in terms of the number of solutions of the system, regardless of any generic property. The morphism $\operatorname{hom}_{W}$ multiplies this number of solutions by $\Pi w_{i}$ : for example consider $F=\left(X^{2}-Y, X-1\right)$, with the weights $(1,2)$; its only solution is $(X=1, Y=1)$, $\operatorname{but~}_{\operatorname{hom}_{W}}(F)=\left(X^{2}-Y^{2}, X-1\right)$ has 2 solutions: $(X=1, Y=1)$ and $(X=1, Y=-1)$.


Figure 4: Algorithmic strategy for a zero-dimensional system, using the $W$-degree

This explains why algorithm FGLM becomes a bottleneck when computing a lexicographical Gröbner basis through $\operatorname{hom}_{W}$ : the transformation artificially demultiplies the number of solutions, leading the algorithms to build and reduce unnecessarily large matrices.

On the other hand, it turns out that applying $\mathrm{hom}_{W}^{-1}$ to a reduced Gröbner basis of hom ${ }_{W}(F)$ yields a Gröbner basis for a weighted equivalent of the GRevLex order. This basis can also be computed directly from the original system by using the weighted degree instead of the total degree in algorithm $F_{4}$ or $F_{5}$.

This new basis can be used as input for algorithm FGLM, and it generates the same ideal as the original system $F$, hence it has the correct degree. Algorithm FGLM can compute a lexicographical basis of the ideal from this basis, with the same complexity bounds as above in terms of the degree of the ideal. The complete computational strategy is summed up on Figure 4.

Regularity properties and complexity bounds The definition of regular sequences does not involve the gradation, so it extends to a weighted setting. Regarding inhomogeneous systems, as in the total degree case, we consider whether the highest $W$-degree components form a regular sequence.

These properties are still generic amongst sequences of polynomials of given $W$-degree, as long as there exists at least one such sequence. It is proved straightforwardly by transposing the genericity proofs in the homogeneous case. However, the existence hypothesis is a necessary one, conditioning the system of weights: take for example $W=(2,3)$ and $D=(5,5)$, the only monomial of $W$-degree 5 being $X Y$, so for any weighted-homogeneous system with respective $W$-degree $D$, the set of solutions is the union of $\{X=0\}$ and $\{Y=0\}$, and it is not a finite set.

We shall give hypotheses on the system of degrees ensuring that this hypothesis is satisfied. For example, if all degrees are divisible by the least common multiple of the weights, or if each degree $d_{i}$ is divisible by the corresponding weight $w_{i}$, then regular sequences exist, and thus are generic. But these hypotheses are not necessary: take for example $W=(2,3)$ and $D=(5,6)$, the sequence $\left(X Y, X^{3}+Y^{2}\right)$ is regular.

Some results about weighted homogeneous systems were known before this thesis. Most notably, Bézout's bound giving the number of solutions of a zero-dimensional ideal admits a weighted version [Spa12, Th. 1.67 and 1.68]:

$$
\operatorname{deg}(I)=\frac{\prod_{i=1}^{n} d_{i}}{\prod_{i=1}^{n} w_{i}}
$$

which gives us immediately the complexity of algorithm FGLM in the strategy described above, for a square system defined by a regular sequence:

$$
C_{\text {FGLM }}=O\left(n\left(\frac{\prod_{i=1}^{n} d_{i}}{\prod_{i=1}^{n} w_{i}}\right)^{3}\right) .
$$

For a run of Algorithm Matrix- $\mathrm{F}_{5}$ on $\operatorname{hom}_{W}(F)$, the usual complexity bounds from the homogeneous case can be applied, but they can actually be improved. The main asymptotical improvement comes from the work of combinatoricians on the Sylvester denumerant of $d$, that is the number of monomials with a given $W$-degree $d$. In particular, it is known [FSoo] that asymptotically, this number of monomials is the number of monomials with total degree $d$, divided by the product of the weights. As such, the complexity bound for Algorithm Matrix- $\mathrm{F}_{5}$ run on $\operatorname{hom}_{W}(F)$ is divided by $\left(\Pi w_{i}\right)^{3}$ when compared to the original bound for homogeneous systems, given on page 15 .
The bound on the degree of regularity can also be improved, with sufficient hypotheses on the system. Assuming that the system form a regular sequence, a first weighted version of Macaulay's bound yields that

$$
d_{\mathrm{reg}} \leq \sum_{i=1}^{n}\left(d_{i}-w_{i}\right)+\max \left\{w_{i}\right\} .
$$

In other words, Algorithm Matrix- $\mathrm{F}_{5}$ may skip the last

$$
\sum_{i=1}^{n}\left(w_{i}-1\right)-\left(\max \left\{w_{i}\right\}-1\right)
$$

degree steps in a run on hom $_{W}(F)$.
This bound can be refined, but we need some stronger genericity hypothesis, containing some information about which variables are involved at each step of the computation. More precisely, we say that $F$ is in simultaneous Noether position if the sequences $\left(f_{1}, X_{2}, \ldots, X_{n}\right)$, $\left(f_{1}, f_{2}, X_{3}, \ldots, X_{n}\right), \ldots,\left(f_{1}, \ldots, f_{n}\right)$ are all regular. Under this hypothesis, the following bound applies:

$$
d_{\mathrm{reg}} \leq \sum_{i=1}^{n}\left(d_{i}-w_{i}\right)+w_{n} .
$$

In particular, when applicable, it implies that for the best complexity, one should order the variables such that the smallest variable (for the monomial order) is the one with the smallest weight.
This bound is still not sharp in full generality. We conjecture that with no hypothesis on the weights and $W$-degrees as long as there exist regular sequences, the following bound is true and optimal:

$$
\begin{equation*}
d_{\mathrm{reg}}=w_{n}\left\lceil\frac{\sum_{i=1}^{n}\left(d_{i}-w_{i}\right)-g}{w_{n}}\right\rceil \tag{1}
\end{equation*}
$$

with $g$ is the Frobenius number of $\left(w_{1}, \ldots, w_{n}\right)$, that is the highest $W$-degree at which there exist no monomial (with the convention that $g=-1$ if there always exists a monomial, that is if $w_{i}=1$ for some $i$ ).


Figure 5: Contrast optimization: both pictures were taken using the same setup of two tubes, the inner one containing deoxygenated blood and the outer one containing oxygenated blood. The left-hand picture is the reference picture, with no contrast optimization; the right-hand picture is contrast-optimized by saturating the inner sample.

Examples This structure appears in many applications, we already gave the example of the work of the authors of $[\mathrm{Fau}+13]$ on the DLP for Edwards elliptic curves.

Another example where weighted homogeneity appears naturally is polynomial inversion systems. Recall that given a set of polynomials $f_{1}, \ldots, f_{n}$ in $X_{1}, \ldots, X_{n}$, we may obtain their relations by adding variables $F_{1}, \ldots, F_{n}$ to the algebra, and eliminating $X_{1}, \ldots, X_{n}$ from $\left\langle F_{1}-\right.$ $\left.f_{1}, \ldots, F_{n}-f_{n}\right\rangle$. There is a natural system of weights to try when working with such systems, as first described in [Tra96]: if we give the $F_{i}$ 's a high-enough weight, we can ensure that the highest $W$-degree components of the system form a regular sequence and a sequence in Noether position. Furthermore, if the $f_{i}$ have additional structure, for example if they are homogeneous, using the degree of $f_{i}$ as the weight of $F_{i}$ makes the highest $W$-degree components even larger, thus improving the regularity of the behavior of Algorithm $\mathrm{F}_{5}$.

In Section 3.6, we shall present benchmarks computing the relations between fundamental invariants of the cyclic group or the dihedral group, between monomials, and between minors of matrices. We could obtain diverse speed-ups for the first step of the computational strategy (computing a GRevLex basis) for these systems, from a modest 1.5 to 4 for matrix minors to $10^{5}$ for monomials.

### 3.2. Real root classification for determinantal systems Application to contrast optimization

Context The second class of structured systems that we studied arised in an application in medical imagery. Nuclear Magnetic Resonance (MNR) is a technique relying on applying a magnetic field to a body, and measuring the response of different matters. Contrast optimization is the problem of ensuring that two biological matters that we wish to distinguish on a picture, for example blood and water, stand apart: ideally, one should be completely black, and the other as bright as possible.

The saturation method is a contrast optimization technique, consisting of making the magnetic field variable in order to progressively obtain the optimal contrast. In order to find the optimal "path" towards the saturated state, control theory has been involved [Lap+12].

This strategy involves alternating so-called bang-arcs and singular-arcs, and estimating the
complexity of the path requires to estimate the number of switching points, in terms of the matters that we want to study. These switching points are characterized as the singularities of the determinant of a matrix whose entries are polynomials in the parameters and some variables, together with some polynomial inequalities [Bon+13].

This problem has two forms: in its most general form, both matters are free, each matter being described by two parameters. The structure of the system allows to eliminate one parameter, and we need to identify parts of $\mathbb{R}^{3}$ where the number of singularities does not change. A useful easier case is the case where one of the matters is water, leaving only 2 parameters.
This is a particular case of a more general problem: given a matrix $M(\mathbf{X}, \mathbf{G})$ whose entries are polynomial in some variables $\mathbf{X}$ and some parameters G , and a target rank $r$, let $V_{r}$ be the locus at which $M$ has rank at most $r$. We wish to identify a dense subset $U$ of the parameter space, together with a covering of $U$ with open subsets, such that on each open subset, the number of singularities of $V_{r}$ in the fiber over the parameters is constant.

In the case of the application, the physics of the system brings an additional constraint on the solutions: they have to be in the Bloch ball, which is described by inequations. Thus the problem is a real roots classification problem on a semi-algebraic set.

Algorithmic strategy The problem is a real roots classification problem for a semi-algebraic set described by a parameterized determinantal system and some inequalities.

The Cylindrical Algebraic Decomposition (CAD) method can be used to solve this problem, but none of its implementations can attack the application directly. On the other hand, it is known that adding equational constraints can dramatically improve the complexity of CAD computations. In our context, these equations should be the equations of the borders of the subdivision areas in the parameter space. This is the crux of real roots classification algorithms such as those presented in [LRo7; YHXo1]. However, none of the existing implementations of these algorithms for this problem were able to tackle the general case of the imagery problem.

We propose a strategy taking advantage of the determinantal structure of the varieties at hand, to guide existing algorithms towards a solution.

Under some hypotheses satisfied by the application, these strategies classify the cardinality of the fibers by identifying critical points of the projections, that is points where several single roots merge into one multiple root, and points where the variety meets the boundary of the relevant domain. In other words, we need to compute the singular locus of a determinantal variety, and then the singularities of this singular locus, the critical points of a projection restricted to it, and its intersection with the boundary of the domain.

The classical tool for computing critical and singular points is the Jacobian criterion, which characterizes singularities in terms of minors of the Jacobian matrix of the system defining the variety. In other words, a singular locus is again a determinantal variety.

Furthermore, the singular locus of a determinantal ideal is closely related to smaller minors of the same matrix. In particular, points where the rank of the matrix is $r-1$ are always singular points of the locus where the rank is $r$.

This is the first property that we use to steer the computations: we compute a rank stratification of the solutions, by separating those singularities where the matrix has exactly rank $r$ (which we can compute using the Jacobian criterion), and those points where the matrix has
rank at most $r-1$, which are automatically singularities of the ideal.
Furthermore, recall that the classification strategy aims at computing a set of hypersurfaces of the parameter space $\mathbb{R}^{t}$, splitting it into open sets over which the number of singularities is constant. So if we know that some component $C$ of the singular locus of the determinantal variety has dimension at most $t-1$, we only need to compute the equation of one hypersurface containing its image by the projection on the parameter space: this hypersurface will necessarily contain the projection of the singular points of $C$, the critical values of the projection restricted to $C$ and the intersection of $C$ with the boundaries of the domain.

Under some generic hypotheses, the dimension of rank strata is determined by the target rank defect and the number of variables and parameters. It turns out that the high rank stratum necessarily has dimension less than $t$, thus reducing the problem to a single elimination for that stratum.

The last regularization step is the use of incidence variety to model singularities of the determinants where needed. Algorithmically, we add $k(k-r)$ variables to the system, build a matrix $N$ with $k$ rows and $k-r$ columns, each entry being one of the new variables, and we build the system with the entries of $M \cdot N$. In order to ensure that the resulting kernel vectors are linearly independent, we pick at random a matrix $U$ with $k-r$ rows and $k$ columns, and we add to the system the entries of $U \cdot N-\operatorname{Id}_{k-r}$.

Results for the application This strategy was used for answering questions related to the contrast imaging problem.

In particular, an open question for the case of water was whether the number of singularities was a global invariant, not depending on the second matter: for all experimental samples, there was only one such singularity.

It turns out that the answer is no: we identified 3 areas of the parameter space where the number of singularities is 1,2 and 3 respectively. See Figure 6: the red area is out of the physically relevant domain; for parameters in the areas containing blue diamonds, there are two singularities in the fiber; for parameters in the areas containing green circles, there are three singularities in the fiber; and for all other areas, there is only one singularity.

These results can be obtained in about 10 s using the new strategy, while a direct approach using Gröbner bases requires 100 s . Furthermore, we were able to use this strategy to obtain the boundaries of the classification in the general case. This computation takes less than 2 h , while it is intractable for the general strategies.

## 4. Perspectives

In this section, we list some questions left open by this thesis, as well as potential extensions and generalizations that we would find interesting.

Open questions and conjectures for $W$-homogeneous systems First, there is the problem of $W$-compatibility: what are the systems of degrees and weights such that there exists $W$-homogeneous regular sequences, and thus that they are generic? We were able to give sufficient conditions, but counter-examples show that they are not necessary.


Figure 6: Decomposition of the parameter space in the case of water: the parameters are the relaxation times for the second matter, $\gamma_{2}$ and $\Gamma_{2}$, and in each area delimited by the black curves, the cardinality of the fiber of the singularities of $V_{r}$ is constant and written as a circled number; the red area corresponds to physically irrelevant parameters.

Genericity results are useful because testing whether a given system is a regular sequence is not easy in general: if one knows that most sequences similar to the system at hand are regular, it makes sense to run the algorithms assuming that the hypothesis is satisfied.

As of now, for systems of degrees and weights falling out of the sufficient conditions for regular sequences to exist, this assumption is riskier. Obtaining a set of necessary and sufficient conditions could complete the characterization of the genericity of weighted homogeneous regular sequences and make this procedure more likely to succeed.
Second, we conjectured formula (1) as a sharp bound of the degree of regularity of a sequence in simultaneous Noether position, without any hypothesis on the weights. So far, we have not been able to prove it.

And last, it would be interesting to be able to characterize semi-regular $W$-homogeneous sequences in terms of their Hilbert series, once again without any hypothesis on the weights.

Several systems of weights Weighted homogeneous systems have been studied considering only one system of weights. However, some systems may be $W$-homogeneous for several systems of weights $W$. Consider for example the polynomial

$$
X^{3} Y^{2} Z+X^{4} Z^{2}+X^{2} Y^{4},
$$

it is ( $1,1,1$ )-homogeneous with degree 6 , and ( $1,2,3$ )-homogeneous with degree 10 .
This new constraint gives extra structure to the Macaulay matrix of the system. Exploiting this structure should yield further improvements on the complexity of Matrix- $F_{5}$ for these systems.

Furthermore, polynomials with are $W_{i}$-homogeneous for all systems of weights $W_{i}$ in some family are also $W$-homogeneous for any $W$ obtained by linear combination of the $W_{i}$ 's: for example, the above polynomial is also ( $2,3,4$ )-homogeneous with degree 16 . It may be interesting to see if this property can be used to find systems of weights leading to easier computations.

Non-positive weights An extension of the previous point could be to consider systems of weights including weights set to zero or a negative integer. With only one system of weights, this leads to problems because homogeneous components do not necessarily have finite dimension. However, with several systems of weights, it is possible to ensure that this does not happen, and it would be useful to know what algorithmic strategies we can use in this situation.

This class of systems is very wide, it would for example include multi-homogeneous systems (systems where the variables are split into groups, and polynomials are homogeneous in total degree for each group): for example, consider a polynomial $f$ in $\mathbb{K}\left[X_{1}, X_{2}, Y_{1}, Y_{2}\right]$, and assume that it is bilinear for the variables $X_{1}, X_{2}$ and $Y_{1}, Y_{2}$. In terms of weights, consider the systems of weights $W_{1}=(1,1,0,0)$ and $W_{2}=(0,0,1,1) ; f$ being bilinear means that $f$ is $W_{1}$-homogeneous with $W_{1}$-degree 1 and $W_{2}$-homogeneous with $W_{2}$-degree 1 .

Weight search Returning to the case of only one system of weights, until now, we have only worked where systems for which the user was able to guess a good system of weights (either because it was natural in the application, or through trial and error). It would be interesting to describe a strategy, or at least a good set of heuristics, giving good candidate systems of weights for an arbitrary inhomogeneous system.

Testing that a system is a regular sequence is not computationally easy, so using that as a criterion is impractical. On the other hand, we could look for systems of weights making the highest $W$-degree component large-enough for it to be generically regular.

Experimentally, we also observed systems which behaved better with a system of weights such that the highest $W$-degree component is small, but with a very large $W$-homogeneous component at smaller $W$-degree. Characterizing these situations could lead to further improvements on a strategy for picking up good systems of weights.

Application to imagery Our work on real roots classification for the imagery problem also left several open questions and possible extensions.

First, we only described algorithms, and showed experimental data confirming that they are more efficient. It would be interesting to complete the analysis with complexity bounds for the whole process.

Furthermore, the algorithmic strategy described for the roots classification problem for the contrast imagery problem builds a stratification by rank. To the best of our knowledge, the question of whether this stratification is physically meaningful has never been raised.

As for the extensions, the root classification problem is only one of the many computational problems posed by specialists in control theory working on constrast optimization for the MNR. For example, they would be interested in a rational parametrisation of the roots, or in a classification of the parameter space according to the number of connected components in the fibers for higher-dimensional varieties. All these questions offer new computational challenges.

## Part I

## Preliminaries

## Chapter 1

## Algebra and geometry

In this chapter, we recall concepts of commutative algebra and algebraic geometry that will be used in the rest of the thesis. In Section 1.1, we give some basic definitions about rings and ideals, and in particular the properties of gradations on polynomial algebras. In Section 1.2, we focus on some properties of polynomial systems, which will be useful in the context of Gröbner bases. In Section 1.3, we briefly recall the correspondence between algebraic and geometric notions, and in particular we give an interpretation of the previous properties. In Section 1.4, we examine the genericity of regularity properties. Finally, in Section 1.5, we give some definitions and properties of varieties characterizing low-rank matrices.

This chapter and the next one do not contain any original contribution. For most results, we give a reference to a point in the literature where a proof may be found. We only give an explicit proof for folklore results, or when the proof itself is useful for subsequent chapters (e.g. if it is generalized later).

### 1.1. Commutative algebra

### 1.1.1. Ideals

We assume that classic definitions about rings, algebras and ideals are known. We shall consider rings and algebras, which we will suppose to be commutative and with a unity.

Definition 1.1. Let $R$ be a ring. A divisor of zero is an element $x \in R$, non-zero, such that

$$
\exists y \in R, y \neq 0, x y=0
$$

Definition 1.2 (Noetherian ring). A ring $R$ is noetherian if and only if any increasing chain of ideals

$$
I_{0} \subseteq I_{1} \subseteq \cdots \subseteq I_{n} \subseteq \ldots
$$

is stationary.
This property is sometimes defined using the following classic characterization:
Proposition 1.3 ([Eis95, Sec. 1.4]). A ring $R$ is noetherian if and only if all ideals of $R$ are finitely generated.

Proof. Let $R$ be a noetherian ring and $I$ an ideal of $R$. If $I$ is not finitely generated, then we can find a family of polynomials $\left(f_{i}\right)_{i \in \mathbb{N}} \subset I$ such that no finite subset of this family generates $I$. Let $I_{i}=\left\langle f_{0}, \ldots, f_{i}\right\rangle$, the sequence of ideals $\left(I_{i}\right)_{i \in \mathbb{N}}$ is strictly increasing, which is impossible since $R$ is noetherian.

Conversely, let $R$ be a ring such that all ideals of $R$ are finitely generated. Let $I_{0} \subseteq I_{1} \subseteq \cdots \subseteq$ $I_{n} \subseteq \ldots$ be an increasing chain of ideals, the limit $\bigcup_{i \in \mathbb{N}} I_{i}$ is an ideal, and by hypothesis it is finitely generated. All of its generators necessarily belong to one of the ideals in the chain, and the sequence stations at that point.

### 1.1.2. Graded rings

Definition 1.4 (Gradation on a ring, degree, homogeneous components). Let $R$ be a ring (resp. algebra), $R$ is graded if and only if there exists a decomposition (called a gradation or grading of $R$ ) into additive groups (resp. vector spaces)

$$
R=\bigoplus_{i=0}^{\infty} R_{i}
$$

such that $R_{i} \cdot R_{j} \subset R_{i+j}$ : the product of elements of degree $i$ and $j$ has degree $i+j$. We say that the elements of $R_{i}$ are homogeneous of degree $i$.

Let $p \in R$, there exists $d \in \mathbb{N}$ and $p_{i} \in R_{i}(i \in\{0, \ldots, d\})$ such that

$$
p=\sum_{i=0}^{d} p_{i} \text { and } p_{d} \neq 0
$$

We then say that $p$ has degree $d$, and $p_{i}$ is called the homogeneous component of $p$ of degree $i$. If all $p_{i}$ 's are zero except for one, $p$ is homogeneous.

Definition 1.5 (Homogeneous ideal). Let $R$ be a graded ring, and $I$ an ideal of $R$. We say that $I$ is homogeneous if it satisfies one of the following equivalent conditions:

1. $I$ is generated by homogeneous elements
2. For any $p \in I$, all homogeneous components of $p$ belong to $I$

Definition 1.6 (Graded morphism). Let $R$ and $S$ be two graded rings. A morphism $\varphi: R \longrightarrow S$ is said to be graded if it sends homogeneous elements of $R$ onto homogeneous elements of $S$ :

$$
\forall d \in \mathbb{N}, \exists e \in \mathbb{N}, \varphi\left(R_{d}\right) \subset S_{e}
$$

### 1.1.3. Polynomial algebras

Definition 1.7. Let $\mathbb{K}$ be a field and $n$ a positive integer, let $A$ be the polynomial algebra $\mathbb{K}[\mathrm{X}]$ with coefficients in $\mathbb{K}$ and $n$ indeterminates $\mathbf{X}=X_{1}, \ldots, X_{n}$.

A product of indeterminates $X_{1}^{\alpha_{1}} \cdots X_{n}^{\alpha_{n}}$ is a monomial, and the integer vector $\left(\alpha_{1}, \ldots, \alpha_{n}\right)$ is its exponent.

A term is the product of a coefficient and a monomial.

In particular, monomials generate a polynomial algebra as a $\mathbb{K}$-vector space, and terms generate a polynomial algebra as an additive group.

We shall use the same notations throughout this section.
Theorem 1.8. The polynomial algebra $A$ is a noetherian ring.
This is a consequence of a more general result:
Theorem 1.9 (Hilbert's basis theorem, [Eis95, Th. 1.2]). Let $R$ be a noetherian ring, then the polynomial algebra $R[X]$ is noetherian.

Remark 1.10. Theorem 1.8 for polynomial rings admits a more elementary proof than the more general Theorem 1.9, see for example [CLOo7, Ch. 2, sec. 5, th. 4] or [Eis95, Ex. 15.15].

An important feature of polynomial algebras is that they can be graded.
Definition 1.11 (Total degree, weighted degree). Let $W=\left(w_{1}, \ldots, w_{n}\right) \in \mathbb{N}^{n}$ be a family of integers. We call $W$ a system of weights, and define the $W$-degree (or weighted degree) of a monomial as

$$
\operatorname{deg}_{W}\left(\prod_{i=1}^{n} X_{i}^{\alpha_{i}}\right)=\sum_{i=1}^{n} w_{i} \alpha_{i}
$$

and let $A_{W, d}$ be the vector space generated by all monomials $m$ such that $\operatorname{deg}_{W}(m)=d$. Polynomials in $A_{W, d}$ are called $W$-homogeneous (or weighted homogeneous).

In the special case $W=(1, \ldots, 1)$, the $W$-degree is called total degree, or simply degree. In this case, $W$-homogeneous polynomials are simply called homogeneous.

Proposition 1.12. Let $A=\bigoplus A_{d}$ be a gradation such that monomials are homogeneous. Then there exists $W \in \mathbb{N}^{n}$ such that

$$
\forall d \in \mathbb{N}, A_{d}=A_{W, d}
$$

Proof. Consider a degree function deg on $A$. For all $i \in\{1, \ldots, n\}$, let $w_{i}=\operatorname{deg}\left(X_{i}\right)$. The multiplicative property of the degree implies that deg is $\operatorname{deg}_{\left(w_{1}, \ldots, w_{n}\right)}$.

Throughout this thesis, we will only consider non-degenerate systems of weights, that is systems of weights such that $w_{i}>0$ for all $i$.

An important concept is that of homogeneization:
Definition 1.13 (Homogenization of polynomials). Let $A_{h}=\mathbb{K}[H, \mathbf{X}]$. The homogeneization morphism is defined as

$$
\text { -h: } \begin{aligned}
A & \longrightarrow A_{h} \\
f & \longmapsto H^{d} f\left(\frac{X_{1}}{H}, \ldots, \frac{X_{n}}{H}\right) \text { where } d \text { is the degree of } f
\end{aligned}
$$

The dehomogenization morphism is defined as

$$
\begin{aligned}
\bullet^{a}: A_{h} & \longrightarrow A \\
f & \longmapsto f(H=1)
\end{aligned}
$$

It has the following immediate properties:

## Proposition 1.14.

- The homogenization morphism is graded and sends polynomials of degree d onto polynomials of degree d;
- For any polynomial $f$ in $A$ of degree $d$ (not necessarily homogeneous), $f^{h}$ is homogeneous of degree d;
- For any $f \in A,\left(f^{h}\right)^{a}=f$;
- For any $f_{h} \in A_{h}$, if $f_{h}$ is not divisible by $H,\left(f_{h}^{a}\right)^{h}=f_{h}$.

Definition 1.15 (Homogenization of ideals). Let $I$ be an ideal of $A$, the homogenization ideal $I^{h}$ is defined as

$$
I^{h}=\left\langle f^{h} \mid f \in I\right\rangle
$$

let $J$ be an ideal of $A_{h}$, the dehomogenization ideal $I^{a}$ is defined as

$$
I^{a}=\left\langle f^{a} \mid f \in J\right\rangle
$$

The element-wise properties above have the following consequences on ideals:
Proposition 1.16. Let $I$ be an ideal of $A$ and $J$ be an homogeneous ideal of $A_{h}$. Then:

- $I^{h}$ is an homogeneous ideal of $A_{h}$;
- $\left(I^{h}\right)^{a}=I$;
- $\left(J^{a}\right)^{h}=J$.


### 1.1.4. Combinatorics of monomials

Let $d \in \mathbb{N}$, the homogeneous component $A_{d}$ forms a finite dimensional $\mathbb{K}$-vector space, and its dimension is the number of monomials at degree $d$.

Counting these monomials in the total degree case is a classic combinatorial exercise:
Proposition 1.17. Let $n \in \mathbb{N}$ and $d \in \mathbb{N}$. The number of monomials in $n$ variables having total degree d is

$$
N_{n, d}=\binom{n+d-1}{d}
$$

The number of monomials in $n$ variables having total degree at most $d$ is

$$
N_{n, d}^{\prime}=N_{n+1, d}=\binom{n+d}{d}
$$

In the more general case, this number is called a denumerant, as introduced by Sylvester:

Definition 1.18 (Sylvester denumerant, [SF82]). Let $W$ be a set of integers and $d \in \mathbb{N}$, the denumerant $N_{W, d}$ is the number of non-negative integer solutions $\left(\alpha_{1}, \ldots, \alpha_{n}\right)$ to the equation

$$
w_{1} \alpha_{1}+\cdots+w_{n} \alpha_{n}=d .
$$

There are formulas expressing the denumerants for some specific weights $W$ [Com74, Sec. 2.6], but in the most general case, there is no known formula expressing a denumerant as a function of $W$ and $d$. The asymptotic behavior of these denumerants is well-known:

Proposition 1.19 ([FSo9, Prop. IV.2]). If $W$ is fixed and $d$ tends towards infinity, then

$$
N_{W, d} \sim \frac{\operatorname{gcd}(W)}{\Pi(W)} N_{n, d}=\frac{\operatorname{gcd}(W)}{\Pi(W)}\binom{n+d-1}{d} .
$$

Additionally, some bounds are known:
Proposition 1.20 ([Agno2, Th. 3.3 and 3.4]). Define

$$
\begin{aligned}
A & :=\sum_{i=2}^{n} w_{i} \frac{\operatorname{gcd}\left(w_{1}, \ldots, w_{i}\right)}{\operatorname{gcd}\left(w_{1}, \ldots, w_{i-1}\right)} \\
B & :=\sum_{i=2}^{n} w_{i}\left(\frac{\operatorname{gcd}\left(w_{1}, \ldots, w_{i}\right)}{\operatorname{gcd}\left(w_{1}, \ldots, w_{i-1}\right)}-1\right)-n \\
\operatorname{gcd}(W) & :=\operatorname{gcd}\left(w_{1}, \ldots, w_{n}\right) \\
\Pi(W) & :=w_{1} \cdots w_{n},
\end{aligned}
$$

where $\operatorname{gcd}\left(a_{1}, \ldots, a_{k}\right)$ denotes the greatest common divisor of the integers $a_{1}, \ldots, a_{k}$. Then the number $N_{W, d}$ of monomials with $W$-degree $d$ is bounded by

$$
\frac{\operatorname{gcd}(W)}{\Pi(W)} N_{n, d-B+1} \leq N_{W, d} \leq \frac{\operatorname{gcd}(W)}{\Pi(W)} N_{n, d+A-n+1} .
$$

Another interesting combinatorial number related to sums of integers is the Frobenius number [Alfo5, Ch. 1].

Definition 1.21 (Frobenius number). Let $W$ be a set of integers, assume that they are coprime. The Frobenius number of $W$ is the largest integer $d$ such that the equation

$$
w_{1} \alpha_{1}+\cdots+w_{n} \alpha_{n}=d .
$$

has no non-negative integer solutions ( $\alpha_{1}, \ldots, \alpha_{n}$ ).
Remark 1.22. Algebraically, it corresponds to the largest $W$-degree $d$ such that there exists no monomials with $W$-degree $d$, or equivalently, such that the homogeneous component of $W$-degree $d$ in $A$ is 0 .
Remark 1.23. Computing this number is in general a hard computational problem: it has been proved to be NP-hard [Alfo5, Th. 1.3.1]. In particular, there is no general closed-form formula if $n>2$ ([Alfo5, Sec. 1.3]).

Remark 1.24. An easy case is if one of the integers is 1 , because then the equation always have solutions, and the Frobenius number does not exist.

Remark 1.25. The Frobenius number is sometimes called coin number, because of its application to the problem of making change in small coins. Another example of application comes from restauration: the McNuggets problem is deciding, given an integer $d$, whether it is possible to obtain $d$ chicken nuggets using boxes of 6,9 and 20 . If it is the case, $d$ is called a McNuggets number, and the Frobenius number of $\{6,9,20\}$ is the largest non-McNuggets number ([Alfo5, Ch. 1]).

### 1.1.5. Hilbert series

In the previous section, we saw that computing the number of monomials at a given $W$-degree directly is a non-trivial problem. However, results such as the asymptotic equivalent of the denumerants (Prop. 1.19) have been known for a long time. These results were proved by considering generating series, instead of individual denumerants.

Definition 1.26 (Generating series). Let $f: \mathbb{N} \longrightarrow \mathbb{N}$ be an integer function. The generating series of $f$ is the formal series

$$
\sum_{d=0}^{\infty} f(d) T^{d} \in \mathbb{N}[[T]]
$$

Generating series are useful because many usual combinatorial operations can be converted into algebraic operations on the generating series. For the case at point, the generating series for denumerants is well known:

Proposition 1.27 ([Com74, Sec. 2.6, Th. A], [FSo9, Prop. I.1]). Let $W$ be a set of integers, the generating series for the denumerants $N_{W, \bullet}$ is given by

$$
\sum_{d=0}^{\infty} N_{W, d} T^{d}=\frac{1}{\prod_{i=1}^{n}\left(1-T^{w_{i}}\right)}
$$

This generating series is a particular case of a Hilbert series:
Definition 1.28 ([Eis95, Ex. 12.12]). Let $R$ be a graded $\mathbb{K}$-algebra such that all homogeneous components, as $\mathbb{K}$-vector spaces, have finite dimension. The Hilbert series of $R$ is the power series defined by

$$
\mathrm{HS}_{R}(T)=\sum_{d=0}^{\infty} \operatorname{dim}_{\mathbb{K}}\left(R_{d}\right) T^{d} \in \mathbb{N}[[T]]
$$

Proposition 1.29. The Hilbert series of $A$, graded with respect to a system of weights $W$, is

$$
\operatorname{HS}_{A}(T)=\frac{1}{\prod_{i=1}^{n}\left(1-T^{w_{i}}\right)}
$$

Hilbert series of quotients $A / I$ where $I$ is a homogeneous ideal will be very useful, because many interesting parameters of the ideal can be read from the series. In section 1.2, we shall see some examples of polynomial systems whose Hilbert series can be computed very easily.

Proposition 1.30 ([Eis95, Ex. 12.12]). Let $I \subset A$ be a $W$-homogeneous ideal. Then there exist a polynomial $P \in \mathbb{Z}[T]$, and an integer $r \in \mathbb{N}$, such that the Hilbert series of $A / I$ can be written

$$
\operatorname{HS}_{A / I}(T)=\frac{P(T)}{\prod_{i=1}^{n}\left(1-T^{w_{i}}\right)}
$$

Remark 1.31. This result will be proved as a by-product of the proof of Prop. 1.38 in Section 1.2.1.
Definition 1.32. With the same notations:

- the codimension of $I$ is the multiplicity $r$ of 1 as a root of $P$;
- the dimension of $I$ is $d=\operatorname{dim}(I)=n-r$;
- the index of regularity of $I$ is $i_{\text {reg }}(I)=\operatorname{deg}(P)-r$;
- if $I$ is homogeneous for the total degree, the degree of $I$ is $\operatorname{deg}(I)=Q(1)$ where $Q=$ $\mathrm{HS}_{A / I}(T) \cdot(1-T)^{d}$.

Remark 1.33. In Section 1.3.4, we shall see that the dimension and degree have geometrical interpretations. For example, if $r=n$, the dimension of $I$ is 0 , and $\mathrm{HS}_{A / I}$ is a polynomial. The degree of $I$ is then the dimension of $A / I$ as a $\mathbb{K}$-vector space. In this case, there are only a finite number of zeroes common to all polynomials of $I$, and this number, counted with multiplicities, is $\operatorname{deg}(I)$.

Hilbert series (and dimension, codimension, etc.) can be defined for non-homogeneous ideals as well, by taking the Hilbert series of $A^{h} / I^{h}$ as a definition.

To conclude this section, we give a last interpretation of the Hilbert series of a homogeneous quotient ring, which will be of immediate use when describing Gröbner basis algorithms.
Definition 1.34 (Macaulay matrix). Let $W$ be a system of weights, $F=\left(f_{1}, \ldots, f_{m}\right) \in \mathbb{K}[\mathbf{X}]^{m}$ a $W$-homogeneous polynomial system with respective $W$-degree $\left(d_{1}, \ldots, d_{m}\right)$, and $d \in \mathbb{N}$ a $W$-degree. For any $i \in\{1, \ldots, m\}$, let $\mathcal{M}_{d}$ be the set of monomials $\mu$ with $W$-degree $d$ : The family $\mathcal{F}_{d}$ of all products $\mu f_{i}^{\prime}$ 's, $i \in\{1, \ldots, m\}, \mu \in \mathcal{M}_{d-d_{i}}$, is a linear basis of the homogeneous component of $\langle F\rangle$ at $W$-degree $d$. The Macaulay matrix of $F$ at $W$-degree $d$ is the matrix of $\mathcal{F}_{d}$ in the canonical monomial basis (whose elements are the monomials of $W$-degree $d$ ).


Remark 1.35. Macaulay matrices are a generalization of Sylvester matrices ([GKZ94, Ch. 12]) to the multivariate case.

Proposition 1.36. If $F$ is a $W$-homogeneous polynomial system, the Hilbert series of $A /\langle F\rangle$ is the generating series of the corank of $\operatorname{Mac}_{d}(F)$, defined as the difference between its number of columns and its rank.

### 1.2. Regularity properties

### 1.2.1. Regular sequences

Definition 1.37 (Regular sequences). Let $m$ be an integer, $W$ a system of weights, and $F=$ $\left(f_{1}, \ldots, f_{m}\right) \in A^{m}$ a family of $W$-homogeneous polynomials. We say that $F$ is a regular sequence if for any $i \in\{1, \ldots, m\}, f_{i}$ is not a zero-divisor in $A /\left\langle f_{1}, \ldots, f_{i-1}\right\rangle$. In this case, we say that the quotient ring $A /\langle f\rangle$ is a complete intersection.

This notion has several equivalent characterizations:
Proposition 1.38 ([Baro4, Prop. 1.7.4(2)]). Let $m$ be an integer, $W$ a system of weights, $F=$ $\left(f_{1}, \ldots, f_{m}\right) \in A^{m}$ a sequence of $W$-homogeneous polynomials, $\left(d_{1}, \ldots, d_{m}\right)$ their respective $W$ degree, and $I=\langle F\rangle$. The Hilbert series of $A / I$ satisfies the inequality (coefficient-wise)

$$
\mathrm{HS}_{A / I}(T) \geq \frac{\prod_{i=1}^{m}\left(1-T^{d_{i}}\right)}{\prod_{i=1}^{n}\left(1-T^{w_{i}}\right)}
$$

with equality if and only if $F$ is a regular sequence.
Proof. For $i \in\{1, \ldots, m\}$, let $I_{i}=\left\langle f_{1}, \ldots, f_{i}\right\rangle$. We will prove the result by recurrence on $i$.
If $i=0$, then $A / I=A$ and the result is a restatement of the formula for the Hilbert series of a polynomial algebra.
Let $i>0$ and assume that the result is proved for $I_{i-1}$. The multiplication by $f_{i}$ gives an exact sequence:

$$
0 \longrightarrow K_{i} \longrightarrow A / I_{i-1} \xrightarrow{\cdot f_{i}} A / I_{i-1} \longrightarrow A / I_{i} \longrightarrow 0
$$

where $K_{i}$ is the kernel of the application. For any $d \in \mathbb{N}$, this exact sequence restricted to $W$-homogeneous components of $W$-degree $d$ is:

$$
0 \longrightarrow\left(K_{i}\right)_{d} \longrightarrow\left(A / I_{i-1}\right)_{d} \xrightarrow{\cdot f_{i}}\left(A / I_{i-1}\right)_{d+d_{i}} \longrightarrow\left(A / I_{i}\right)_{d+d_{i}} \longrightarrow 0
$$

which implies for the dimensions:

$$
\operatorname{dim}\left(\left(A / I_{i}\right)_{d+d_{i}}\right)=\operatorname{dim}\left(\left(A / I_{i-1}\right)_{d+d_{i}}\right)-\operatorname{dim}\left(\left(A / I_{i-1}\right)_{d}\right)+\operatorname{dim}\left(\left(K_{i}\right)_{d}\right) .
$$

Multiplying by $T^{d}$ and taking the sum for $d$ ranging over $\mathbb{N}$, we get

$$
\mathrm{HS}_{A / I_{i}}(T)=\left(1-T^{d_{i}}\right) \mathrm{HS}_{A / I_{i-1}}(T)+\sum_{d=0}^{\infty} T^{d+d_{i}} \operatorname{dim}\left(\left(K_{i}\right)_{d}\right)
$$

By recurrence hypothesis,

$$
\mathrm{HS}_{A / I_{i-1}}(T)=\frac{\prod_{j=1}^{i-1}\left(1-T^{d_{j}}\right)}{\prod_{j=1}^{n}\left(1-T^{w_{j}}\right)}+R_{i-1}(T),
$$

where $R_{i-1}(T)$ is a series whose coefficients are all nonnegative. Hence

$$
\begin{equation*}
\mathrm{HS}_{A / I_{i}}(T)=\frac{\prod_{j=1}^{i}\left(1-T^{d_{j}}\right)}{\prod_{j=1}^{n}\left(1-T^{w_{j}}\right)}+R_{i-1}(T)+\sum_{d=0}^{\infty} T^{d+d_{i}} \operatorname{dim}\left(\left(K_{i}\right)_{d}\right) \tag{1.1}
\end{equation*}
$$

and all coefficients in the rightmost summand are nonnegative as well, so we get the wanted inequality.

Furthermore, if the sequence $\left(f_{1}, \ldots, f_{i}\right)$ is regular, then $K_{i}=0$ and by recurrence hypothesis, $R_{i-1}(T)=0$, so there is indeed equality. Conversely, assume that

$$
\operatorname{HS}_{A / I_{i}}(T)=\frac{\prod_{j=1}^{i}\left(1-T^{d_{j}}\right)}{\prod_{j=1}^{n}\left(1-T^{w_{j}}\right)}
$$

From Equation (1.1), $R_{i-1}(T)=0$ and $K_{i}=0$. By recurrence hypothesis, if $R_{i-1}(T)=0$, then $\left(f_{1}, \ldots, f_{i-1}\right)$ is a regular sequence; and by definition, if $K_{i}=0$, then $f_{i}$ is not a zero divisor in $A / I_{i-1}$, so $\left(f_{1}, \ldots, f_{i}\right)$ is indeed regular.

Corollary 1.39 ([Baro4, Prop. 1.7-4(3)]). If $F$ is a regular sequence, then any permutation of $F$ is also regular.

Proof. The characterization from Proposition 1.38 does not depend on the order of the polynomials.

### 1.2.2. Description of the Hilbert series of a homogeneous regular sequence

In the total degree case, the Hilbert series of a zero-dimensional regular sequence admits the following description.

Theorem 1.40 ([Mor96, Prop. 2.2]). Let $D=\left(d_{1}, \ldots, d_{n}\right) \in \mathbb{N}^{n}$, and let $a_{d}$ be the coefficient of degree d in the series

$$
S_{D}(T)=\frac{\prod_{i=1}^{n}\left(1-T^{d_{i}}\right)}{(1-T)^{n}}
$$

Further let

- $\delta=\sum_{i=1}^{n}\left(d_{i}-1\right) ;$
- $\delta^{*}=\sum_{i=1}^{n-1}\left(d_{i}-1\right)=\delta-d_{n}+1$;
- $\sigma=\min \left(\delta^{*},\lfloor\delta / 2\rfloor\right)$;
- $\mu=\delta-2 \sigma$.

Then:

- the series $S_{D}(T)$ is a polynomial of degree $\delta$
- it is self-reciprocal (i.e. for anyd, $a_{d}=a_{\delta-d}$ )
- the coefficients $a_{d}$ satisfy:

$$
\begin{aligned}
\forall d \in\{0, \ldots, \sigma-1\}, & a_{d}<a_{d+1} \\
\forall d \in\{\sigma, \ldots, \sigma+\mu-1\}, & a_{d}=a_{d+1} \\
\forall d \in\{\sigma+\mu, \ldots, \delta\}, & a_{d}>a_{d+1}
\end{aligned}
$$

Example 1.41. Take a homogeneous regular sequence of $A=\mathbb{K}[X, Y, Z]$ with respective total degree $(2,4,8)$ : for example, $F=\left(X^{2}, Y^{4}, Z^{8}\right)$. The Hilbert series of $A /\langle F\rangle$ is

$$
\begin{aligned}
H(T) & =\frac{\left(1-T^{2}\right)\left(1-T^{4}\right)\left(1-T^{8}\right)}{(1-T)^{3}} \\
& =(1+T)\left(1+T+T^{2}+T^{3}\right)\left(1+T+\cdots+T^{7}\right) \\
& =1+3 T+5 T^{2}+7 T^{3}+8 T^{4}+8 T^{5}+8 T^{6}+8 T^{7}+7 T^{8}+5 T^{9}+3 T^{10}+T^{11} .
\end{aligned}
$$

The value of the coefficient of degree $d$ in this series, as a function of $d$, is plotted in Figure 1.1.
In this case, the parameters from Theorem 1.40 are:

$$
\begin{aligned}
& \delta=(2-1)+(4-1)+(8-1)=11 \\
& \delta^{*}=(2-1)+(4-1)=4 \\
& \sigma=\min (4,10)=4 \\
& \mu=11-2 \cdot 3=3
\end{aligned}
$$

Indeed, the polynomial $H(T)$ is self-reciprocal (the plot is symmetric around $d=5.5$ ) and its coefficients are increasing until degree $\sigma$, then stationary until degree $\sigma+\mu$, then decreasing.

In Section 3.2.5 we shall generalize this theorem in a weighted setting.

### 1.2.3. Noether position

We recall the following definition of ring theory:
Definition 1.42 (Integral element, integral morphism, integral extension). Let $R \subset S$ be two rings, and $x \in S$. An element $x \in S$ is integral over $R$ if there exists a monic polynomial with coefficients in $R$ annihilating $x$ :

$$
\exists r_{0}, \ldots, r_{d-1} \in R, x^{d}+r_{d-1} x^{d-1}+\cdots+r_{0}=0
$$

An integral morphism $\varphi: R \longrightarrow S$ is a morphism such that any element of $S$ is integral over $\varphi(R)$. If the inclusion $R \hookrightarrow S$ is an integral morphism, then $R \subset S$ is an integral extension.


Figure 1.1: Shape of the Hilbert series of a homogeneous regular sequence with total degree $(2,4,8)$

Definition 1.43 (Noether position). Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a sequence of polynomials in $\mathbb{K}[\mathrm{X}]$. The ideal $\langle F\rangle$ is said to be in Noether position w.r.t. the variables $X_{1}, \ldots, X_{m}$ if it satisfies the two following properties:

- for $i \in\{1, \ldots, m\}$, the canonical image of $X_{i}$ in $\mathbb{K}[\mathrm{X}] / I$ is an algebraic integer over $\mathbb{K}\left[X_{m+1}, \ldots, X_{n}\right] ;$
- $\mathbb{K}\left[X_{m+1}, \ldots, X_{n}\right] \cap I=0$.

In other words, the morphism

$$
\mathbb{K}\left[X_{m+1}, \ldots, X_{n}\right] \longrightarrow \mathbb{K}\left[X_{1}, \ldots, X_{n}\right] /\langle F\rangle
$$

is integral and injective.
The system $F$ is said to be in simultaneous Noether position (or in SNP) w.r.t. the order $X_{1}>X_{2}>\cdots>X_{n}$ if for any $1 \leq i \leq m$, the ideal $\left\langle f_{1}, \ldots, f_{i}\right\rangle$ is in Noether position w.r.t. the variables $X_{1}, \ldots, X_{i}$.

To conclude this section, we give several useful characterizations of Noether position. They appear to be folklore (see for example [BFS14, Prop. 6]), we give a proof for completeness.

For $m \in\{1, \ldots, n\}$, let $\theta_{m}$ be the morphism evaluating $X_{m+1}, \ldots, X_{n}$ to 0 :

$$
\begin{array}{cccc}
\theta_{m}: \mathbb{K}[\mathbf{X}] & \longrightarrow & \mathbb{K}[\mathbf{X}] & \\
& X_{i} & \longmapsto & X_{i} \\
(i \leq m) \\
X_{i} & \longmapsto & 0 & (i>m)
\end{array}
$$

Proposition 1.44. Let $m \leq n$, and let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a sequence of polynomials. The following statements are equivalent:

1. the sequence $F$ is in Noether position w.r.t. the variables $X_{1}, \ldots, X_{m}$;
2. the sequence $F_{\text {ext }}:=\left(f_{1}, \ldots, f_{m}, X_{m+1}, \ldots, X_{n}\right)$ is regular;
3. the sequence $\theta_{m}(F)=F\left(X_{1}, \ldots, X_{m}, 0, \ldots, 0\right)$ is in Noether position w.r.t. the variables $X_{1}, \ldots, X_{m}$;
4. the sequence $\theta_{m}(F)$ is regular.

Proof. $(1 \Longrightarrow 2)$. Let $I$ be the ideal generated by $F$. In Section 1.3.5, we shall state a geometric interpretation of Noether position (Proposition 1.65), which shows that the canonical projection onto the $n-m$ last coordinates

$$
\pi: V(I) \longrightarrow V\left(\left\langle X_{1}, \ldots, X_{m}\right\rangle\right)
$$

is a surjective morphism with finite fibers. This implies that the variety $V\left(\left\langle F_{\text {ext }}\right\rangle\right)=\pi^{-1}(0)$ is zero-dimensional, and so the sequence is regular.
$(2 \Longrightarrow 1)$. This statement will be proved using basic notions of Gröbner basis theory, which are introduced in Sections 2.1.1 and 2.1.2. Let $i \leq m$, we want to show that $X_{i}$ is integral over the ring $\mathbb{K}\left[X_{m+1}, \ldots, X_{n}\right]$. Since $F_{\text {ext }}$ defines a zero-dimensional ideal, there exists $n_{i} \in \mathbb{N}$ such that $X_{i}^{n_{i}}=\operatorname{LT}(f)$ with $f \in\left\langle F_{\text {ext }}\right\rangle$ for the Lex ordering with $X_{n}>\cdots>X_{1}$ (Prop. 2.20). By definition of the Lex ordering, all monomials in $f$ are only divisible by $X_{1}, \ldots, X_{i}$, so we may assume that $f$ lies in $I$. This shows that every $X_{i}$ is integral over $\mathbb{K}\left[X_{i+1}, \ldots, X_{n}\right] / I$. We get the requested result by induction on $i$ : first, this is clear if $i=m$. Now assume that we know that $\mathbb{K}\left[X_{i}, \ldots, X_{n}\right] / I$ is an integral extension of $\mathbb{K}\left[X_{m+1}, \ldots, X_{n}\right]$. From the above, we also know that $X_{i-1}$ is integral over $\mathbb{K}\left[X_{i}, \ldots, X_{n}\right]$, and so, since the composition of integral homomorphisms is integral, we get the requested result.

Finally, we want to check the second part of the definition of Noether position. Assume that there is a non-zero polynomial in $\mathbb{K}\left[X_{m+1}, \ldots, X_{n}\right] \cap I$. If this polynomial is constant, then $I$ is not a proper ideal. And if this polynomial has degree at least 1 , it is a non-trivial syzygy between $X_{m+1}, \ldots, X_{n}$, contradicting the regularity hypothesis.
$(2 \Longrightarrow 4)$. For any $i \in\{1, \ldots, m\}$, write $f_{i}^{\prime}=f_{i}\left(X_{1}, \ldots, X_{m}, 0, \ldots, 0\right)$. Since any permutation of a regular sequence is a regular sequence, $\left(X_{m+1}, \ldots, X_{n}, f_{1}, \ldots, f_{m}\right)$ is a regular sequence, that is, for any $1 \leq i \leq m, f_{i}$ is not a zero divisor in

$$
\mathbb{K}\left[X_{1}, \ldots, X_{n}\right] /\left\langle X_{m+1}, \ldots, X_{n}, f_{1}, \ldots, f_{i-1}\right\rangle
$$

As a consequence, factoring in the quotient by $\left\langle X_{m+1}, \ldots, X_{m}\right\rangle, f_{i}^{\prime}$ is no zero-divisor in

$$
\mathbb{K}\left[X_{1}, \ldots, X_{m}\right] /\left\langle f_{1}^{\prime}, \ldots, f_{i-1}^{\prime}\right\rangle
$$

$(4 \Longrightarrow 2)$. For any $i$, write $f_{i}=f_{i}^{\prime}+r_{i}$ with $f_{i}^{\prime} \in \mathbb{K}\left[X_{1}, \ldots, X_{m}\right]$, and $r_{i} \in\left\langle X_{m+1}, \ldots, X_{n}\right\rangle$. Let $1 \leq i \leq n$. Assume that $g f_{i} \in\left\langle X_{m+1}, \ldots, X_{n}, f_{1}, \ldots, f_{i-1}\right\rangle$ :

$$
\begin{aligned}
g f_{i}=g f_{i}^{\prime}+g r_{i} & =\sum_{j=1}^{i-1} g_{j} f_{j}+\sum_{j=m+1}^{n} g_{j} X_{j} \\
& =\sum_{j=1}^{i-1} g_{j} f_{j}^{\prime}+R \quad \text { with } R \in\left\langle X_{m+1}, \ldots, X_{n}\right\rangle .
\end{aligned}
$$

As a consequence, considering only the monomials in $\mathbb{K}\left[X_{1}, \ldots, X_{m}\right]$

$$
g^{\prime} f_{i}^{\prime}=\sum_{j=1}^{i-1} g_{j} f_{j}^{\prime} \text { where } g^{\prime}=g\left(X_{1}, \ldots, X_{m}, 0, \ldots, 0\right)
$$

Since $\theta_{m}(F)$ is regular, $g^{\prime} \in\left\langle f_{1}^{\prime}, \ldots, f_{i-1}^{\prime}\right\rangle$ :

$$
g=g^{\prime}+r \in\left\langle f_{1}^{\prime}, \ldots, f_{i-1}^{\prime}\right\rangle+\left\langle X_{m+1}, \ldots, X_{m}\right\rangle=\left\langle f_{1}, \ldots, f_{i-1}\right\rangle+\left\langle X_{m+1}, \ldots, X_{m}\right\rangle
$$

And indeed, $f_{i}$ is no zero-divisor in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right] /\left\langle X_{m+1}, \ldots, X_{n}, f_{1}, \ldots, f_{i-1}\right\rangle$. It means that ( $X_{m+1}, \ldots, X_{n}, f_{1}, \ldots, f_{m}$ ) is a regular sequence. By permutation, we conclude that the sequence $\left(f_{1}, \ldots, f_{m}, X_{m+1}, \ldots, X_{n}\right)$ is regular.
$\left(4 \Longleftrightarrow\right.$ 3). The sequence $\theta_{m}(F)=\left(\theta_{m}\left(f_{1}\right), \ldots, \theta_{m}\left(f_{m}\right)\right) \in \mathbb{K}\left[X_{1}, \ldots, X_{m}\right]^{m}$ is regular if and only if the sequence $\left(\theta_{m}\left(f_{1}\right), \ldots, \theta_{m}\left(f_{m}\right), X_{m+1}, \ldots, X_{n}\right)$ is regular. The equivalence between 3 and 4 is then a mirror of the equivalence between 1 and 2 .

### 1.2.4. Semi-regular sequences

The length of a regular sequence cannot be more than the number of indeterminates in the polynomial algebra. In applications, it is often required to solve so-called overdetermined systems, that is systems with more equations than unknowns. In the case of homogeneous systems, we use the notion of semi-regularity to characterize systems which are "as regular as they can".

Definition 1.45 (Semi-regular system). Let $m$ be an integer, and let $F=\left(f_{1}, \ldots, f_{m}\right) \in A^{m}$ be a sequence of homogeneous polynomials. For any $i \in \mathbb{N}$, let $F_{i}=\left(f_{1}, \ldots, f_{i}\right)$. We say that the sequence $F$ is semi-regular if and only if for any $i \in\{1, \ldots, m\}$ and any $d \in \mathbb{N}$, the linear map given by the multiplication by $f_{i}$ :

$$
\begin{array}{clc}
s_{i, d}: \quad\left(A /\left\langle F_{i-1}\right\rangle\right)_{d} & \longrightarrow & \left(A /\left\langle F_{i-1}\right\rangle\right)_{d+d_{i}} \\
g & \longmapsto & f_{i} g
\end{array}
$$

is full-rank (either injective or surjective).
Proposition 1.46 ([Pario, Prop. 1]). Let $m$ be an integer, and let $F=\left(f_{1}, \ldots, f_{m}\right) \in A^{m}$ be a sequence of homogeneous polynomials. For any $i \in \mathbb{N}$, let $F_{i}=\left(f_{1}, \ldots, f_{i}\right)$.

The sequence $F$ is semi-regular if and only if for all $i \in\{1, \ldots, m\}$, the Hilbert series of $A /\left\langle F_{i}\right\rangle$ is

$$
\mathrm{HS}_{A /\left\langle F_{i}\right\rangle}(T)=\left\lfloor\frac{\prod_{i=1}^{m}\left(1-T^{d_{i}}\right)}{(1-T)^{n}}\right\rfloor
$$

where $\lfloor\bullet\rfloor$ is the truncation before the first negative coefficient.
In particular, regular sequences are semi-regular, and if $F$ is semi-regular, then all subsequences $F_{i}$ are semi-regular. On the other hand, with this definition, permutations of a semi-regular sequence need not be semi-regular.

### 1.3. Algebraic geometry

Throughout this section, $\mathbb{K}$ is a field, $n \in \mathbb{N}$, and $A$ is the polynomial algebra $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$.

### 1.3.1. Nullstellensatz

Algebraic geometry establishes a correspondence between algebraic objects (ideals of polynomial algebras) and geometric objects (solution sets of polynomial equations). The backbone of this correspondence is Hilbert's Nullstellensatz (or theorem of zeroes).

Definition 1.47. Let $I$ be an ideal of $A$, let $\mathbb{L}$ be a field containing $\mathbb{K}$. The affine variety (or in short variety) defined by $I$ in $\mathbb{L}$ is the set

$$
V_{\mathbb{L}}(I)=\left\{(\mathbf{x}) \in \mathbb{L}^{n} \mid \forall f \in I, f(\mathbf{x})=0\right\} .
$$

When the field $\mathbb{L}$ is the algebraic closure of the field of definition $\mathbb{K}$, we simply write this set $V(I)$.
Let $V$ be a subset of $\mathbb{L}^{n}$, we define the ideal associated with $V$ as

$$
I(V)=\{f \in \mathbb{C}[\mathbf{X}] \mid \forall(\mathbf{x}) \in V, f(\mathbf{x})=0\}
$$

Theorem 1.48 (Hilbert's Nullstellensatz, [CLOo7, Ch. 4, sec. 1, th. 2 and sec. 2, th. 4]). Let $\mathbb{K}$ be an algebraically closed field. Then

$$
I(V(I))=\sqrt{I}
$$

where $\sqrt{I}$ is the radical of I.

### 1.3.2. Zariski topology

Algebraic varieties satisfy all axioms of the closed sets of a topology:
Proposition 1.49 ([Har77, Prop. 1.2]).

- $V_{\mathbb{L}}(0)=\mathbb{L}^{n}$
- $V_{\mathbb{L}}(A)=\varnothing$
- $V_{\mathbb{L}}(I) \cup V_{\mathbb{L}}(J)=V(I \cap J)$
- $\cap V_{\mathbb{L}}\left(I_{i}\right)=V\left(\sum I_{i}\right)$

Definition 1.50 (Zariski topology). Taking all subsets $V_{\mathbb{L}}(I) \subset \mathbb{L}^{n}$ as closed sets, one obtains a topology on $\mathbb{L}^{n}$, called the Zariski topology.

Definition 1.51 (Irreducible subset). A topological space is called irreducible if it cannot be written as the union of two proper closed subsets.
Equivalently, it means that two non-empty open subsets have non-empty intersection: in other words, all non-empty open subsets are dense.

Proposition 1.52 ([Har77, Cor. 1.4]). If $\mathbb{L}$ is algebraically closed, irreducible varieties correspond to prime ideals of $\mathbb{Z}[\mathrm{X}]$.

Definition 1.53 (Irreducible component). Let $V \subset \mathbb{L}^{n}$ be a variety. There exists, up to permutation, a unique decomposition

$$
V=V_{1} \cup \cdots \cup V_{k}
$$

such that all $V_{i}$ are irreducible varieties and none contain another. The $V_{i}$ are called the irreducible components of $V$.

Other topological notions (open set, closure, interior, boundary...) are defined as usual. Locally closed sets are defined as the intersection of an open set and a closed set. The last definition of this section is an interesting special case:

Definition 1.54 (Distinguished open subset). Let $V \subset \mathbb{L}^{n}$ be a variety and let $S \subset \mathbb{L}^{n}$ be a hypersurface, that is a variety defined by a single polynomial $f$. The open subset $V \backslash S$ is called a distinguished open subset of $V$.

### 1.3.3. Morphisms and rational maps

We shall be interested in functions between varieties, or between dense subsets of varieties. Because of the underlying algebraic structure, we can define polynomials on varieties, and the
Definition 1.55 (Morphism). Let $V_{1} \subset \mathbb{K}^{n}$ and $V_{2} \subset \mathbb{K}^{m}$ be varieties. A morphism $V_{1} \longrightarrow V_{2}$ is an application $\varphi$ which can be expressed, coordinate-wise, by polynomials:

$$
\varphi: \begin{array}{ccc}
V_{1} & \longrightarrow & V_{2} \\
& \left(x_{1}, \ldots, x_{n}\right) & \longmapsto
\end{array}\left(\varphi_{1}(\mathbf{x}), \ldots, \varphi_{m}(\mathbf{x})\right) \quad \text { with } \varphi_{i} \in \mathbb{K}\left[X_{1}, \ldots, X_{n}\right] / I\left(V_{1}\right)
$$

We say that $V_{1}$ and $V_{2}$ are isomorphic if there exist morphisms $V_{1} \longrightarrow V_{2}$ and $V_{2} \longrightarrow V_{1}$ which are mutually inverse.

Definition 1.56 (Rational map). Let $V_{1} \subset \mathbb{K}^{n}$ and $V_{2} \subset \mathbb{K}^{m}$ be varieties. A rational map $V_{1} \rightarrow V_{2}$ is an application $\varphi$ which can be expressed, coordinate-wise, by rational fractions:

$$
\varphi: \begin{array}{ccc}
U_{1} & \cdots & V_{2} \\
& \left(x_{1}, \ldots, x_{n}\right) & \longmapsto
\end{array}\left(\varphi_{1}(\mathbf{x}), \ldots, \varphi_{m}(\mathbf{x})\right) \quad \text { with } \varphi_{i} \in \mathbb{K}\left(X_{1}, \ldots, X_{n}\right)
$$

where $U_{1}$ is the Zariski-open subset of $V_{1}$ defined as the non-vanishing locus of the denominators of the $\varphi_{i}$.
Definition 1.57 (Birational equivalence). If $V_{1}$ and $V_{2}$ are irreducible varieties, we say that $V_{1}$ and $V_{2}$ are birationally equivalent (or simply birational) if there exist rational maps $V_{1} \rightarrow V_{2}$ and $V_{2} \longrightarrow V_{1}$.

Recall that on irreducible varieties, non-empty open subsets are dense. So two irreducible varieties being birational means that, apart from some closed sets with empty interior, they are in a one-to-one rational correspondence.

A useful example of application of these definitions is the following classic construction, associating a distinguished open subset of a variety in $\mathbb{K}^{n}$ to an algebraic variety in $\mathbb{K}^{n+1}$ :

Proposition 1.58. Let $V \subset \overline{\mathbb{K}}^{n}$ be an irreducible variety, let $f \in \mathbb{K}[\mathrm{X}]$, and let $W$ be the distinguished open subset $V \backslash V(f)$. Consider the variety $W^{\prime} \subset \overline{\mathbb{K}}^{n+1}$, with coordinates $X_{1}, \ldots, X_{n}, U$ defined by $U \cdot f-1=0$. Then there is an birational map with no pole:

$$
\begin{array}{ccc}
W & \simeq & W^{\prime} \\
\left(x_{1}, \ldots, x_{n}\right) & \longmapsto & \left(x_{1}, \ldots, x_{n}, \frac{1}{f(X)}\right) \\
\left(x_{1}, \ldots, x_{n}\right) & \longleftrightarrow & \left(x_{1}, \ldots, x_{n}, u\right)
\end{array}
$$

The ring $\mathbb{K}[\mathrm{X}, U] /\langle U \cdot f-1\rangle$ is called the localization of $\mathbb{K}[\mathrm{X}]$ at $f$, and denoted $\mathbb{K}\left[\mathrm{X}, f^{-1}\right]$ or $\mathbb{K}[\mathrm{X}]_{f}$.

This is not true for general subsets, for example any rational map from $\mathbb{K}^{2} \backslash\{(0,0)\}$ will necessarily have a curve of poles.

### 1.3.4. Dimension and degree

The dimension of an algebraic variety is defined as a topological dimension. We shall see that this definition coincides with the algebraic notion introduced earlier, and with the usual notion of dimension when working over the reals or the complex (in Section 1.3.7).

Definition 1.59 (Krull dimension, equidimensional varieties). Let $V$ be an irreducible variety of $\mathbb{L}^{n}$. The Krull dimension of $V$ (or simply dimension) is the largest $d \in\{0, \ldots, n\}$ such that there exists a chain of irreducible closed subsets

$$
V=V_{1} \supsetneq V_{2} \supsetneq \cdots \supsetneq V_{d} .
$$

Let $V$ be any variety, its dimension is defined as the maximal dimension of one of its irreducible components. Let $X$ be any subset of $\mathbb{L}^{n}$, its dimension is defined as the dimension of its closure. As usual, the codimension of $V$ is $n-d$.
Assume that $\mathbb{L}$ is algebraically closed. If all irreducible components of a variety $V$ have the same dimension, we say that $V$ is equidimensional. Any variety $V$ is the union of a finite set of equidimensional varieties called equidimensional components: the $d$-equidimensional component of $V$ is the union of all irreducible components of $V$ having dimension $d$. The local dimension of $V$ at a point $\mathbf{x} \in V$ is the largest dimension of an irreducible component of $V$ containing x .

Proposition 1.60. If $\mathbb{L}_{1} \subset \mathbb{L}_{2}$ are extensions of $\mathbb{K}$, then

$$
\operatorname{dim}\left(V_{\mathbb{L}_{1}}(I)\right) \leq \operatorname{dim}\left(V_{\mathbb{L}_{2}}(I)\right)
$$

Proof. Consequence of $V_{\mathrm{L}_{1}}(I) \subset V_{\mathrm{L}_{2}}(I)$.
Proposition 1.61 ([Har77, Ex. 1.9]). Let $f_{1}, \ldots, f_{m} \in \mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$. All irreducible components of the variety $V_{\overline{\mathbb{K}}}$ have dimension at least $n-m$.
Proposition 1.62 ([Har77, Prop. 1.7]). Let I be an ideal with dimension d, then $V_{\overline{\mathbb{K}}}(I)$ has dimensiond.

In particular, if $I$ is a 0 -dimensional ideal, $V(I)$ is a finite set of points. The cardinality of this set is at most $\operatorname{deg}(I)$, and it is exactly $\operatorname{deg}(I)$ if points are counted with a notion of multiplicity and $\mathbb{K}$ is algebraically closed. For example, we shall see in Section 2.3.1 ideals generated by equations of the form

$$
\left\{\begin{aligned}
g_{1}\left(X_{1}, \ldots, X_{n}\right) & =X_{1}+r_{1}\left(X_{n}\right) \\
g_{2}\left(X_{2}, \ldots, X_{n}\right) & =X_{2}+r_{2}\left(X_{n}\right) \\
& \vdots \\
g_{n-1}\left(X_{n-1}, X_{n}\right) & =X_{n-1}+r_{n-1}\left(X_{n}\right) \\
g_{n}\left(X_{n}\right) & =X_{n}^{\delta}+r_{n}\left(X_{n}\right)
\end{aligned}\right.
$$

For such an ideal, the quotient $\mathbb{K}[\mathrm{X}] / I$ is isomorphic to a univariate quotient:

$$
\mathbb{K}[\mathrm{X}] / I \simeq \mathbb{K}\left[X_{n}\right] /\left\langle g_{n}\right\rangle
$$

and the dimension of the right-hand side of the equality, as a $\mathbb{K}$-vector space, is the degree $\delta$ of the univariate polynomial $g_{n}$. So by definition, we have $\delta=\operatorname{deg}(I)$.

The system may not have exactly $\delta$ solutions in an algebraic closure of $\mathbb{K}$, but if it does not, it means that $g_{n}$ has roots with multiplicity, and it makes sense to give the same multiplicity to the corresponding points in $V(I)$.

For positive dimensional systems, the degree is interpreted through Bézout's theorem. We shall give a rigorous statement of this theorem in Section 1.3.5, but informally-speaking, it means that the intersection of a variety of degree $d_{1}$ and a hypersurface of degree $d_{2}$ is expected to have degree $d_{1} d_{2}$. Given a $d$-equidimensional variety $V(I)$ in $\mathbb{K}^{n}$, we may thus cut it with $d$ random hyperplanes, in order to obtain a finite set of points whose cardinality (with multiplicities) is the degree $\operatorname{deg}(I)$.

### 1.3.5. Geometric interpretation of the regularity properties

In this section, we look back at the regularity properties defined in Section 1.2, in particular regular sequences and Noether position. We shall examine these properties through the algebrageometry correspondence.

Proposition 1.63 ([Baro4, Prop. 1.7.4(1)], [Spa12, Prop. 1.43]). Let $m$ be an integer, $W$ a system of weights, $F=\left(f_{1}, \ldots, f_{m}\right) \in A^{m}$ a sequence of $W$-homogeneous polynomials, and $I=\langle F\rangle$. Then $F$ is a regular sequence if and only if $V(I)$ is equidimensional with dimension $n-m$.

Proof. From Prop. 1.61, we know that all irreducible components of $V(I)$ have dimension at least $n-m$, and from Prop. 1.62, with the Hilbert series described Prop. 1.38, $V(I)$ has dimension at most $n-m$, and thus all irreducible components of $V(I)$ as well.

Now assume that $F$ is not regular. Let $i$ be such that $f_{i}$ is a zero divisor in $\mathbb{K}[\mathbf{X}] /\left\langle f_{1}, \ldots, f_{i-1}\right\rangle$, and assume that $i$ is minimal, that is $\left(f_{1}, \ldots, f_{i-1}\right)$ is a regular sequence. Let $I_{i-1}=\left\langle f_{1}, \ldots, f_{i-1}\right\rangle$ and let $V_{i-1}=V\left(I_{i-1}\right)$, this variety is equidimensional with dimension $n-i+1$. Let $I_{f}$ be the annihilator of $f_{i}$ in $A_{i-1}$, that is

$$
I_{f}=\left\{g \in \mathbb{K}[\mathbf{X}] / I_{i-1} \mid g f_{i}=0 \in \mathbb{K}[\mathbf{X}] / I_{i-1}\right\}
$$

Since $f_{i}$ is a zero-divisor in $\mathbb{K}[\mathrm{X}] / I_{i-1}, I_{f}$ is a non-zero ideal of $\mathbb{K}[\mathrm{X}] / I_{i-1}$ such that $f I_{f}=0$ in $\mathbb{K}[\mathrm{X}] / I_{i-1}$. Geometrically, it means that

$$
V_{i-1}=\left(V_{i-1} \cap V(f)\right) \cup\left(V_{i-1} \cap V\left(I_{f}\right)\right)
$$

and since $I_{f}$ is non-zero, $V\left(I_{f}\right) \cap V_{i-1}$ is a proper subvariety of $V_{i-1}$. Hence $V(f) \cap V_{i-1}$ is non-empty, and $f$ vanishes on an irreducible component $V$ of $V_{i-1}$. So $V\left(\left\langle f_{1}, \ldots, f_{i}\right\rangle\right)$ contains $V$, and has dimension $n-i+1$. So $V(F)$ contains

$$
V \cap V\left(\left\langle f_{i+1}, \ldots, f_{m}\right\rangle\right)
$$

and by Prop. 1.61, it has dimension at least $n-m+1>n-m$.
In particular, if $F$ is a regular sequence of length $n$, the variety $V(\langle f\rangle)$ has dimension 0 , which means that it is a finite set of points. Systems defined by a zero-dimensional regular sequences, sometimes called square systems (by analogy with linear algebra, where systems with as many equations as unknowns are defined by a square matrix), play a particular role in resolution strategies.

Theorem 1.64 (Bézout, consequence of [Spa12, Th. 1.67 and 1.68]). Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a system of $W$-homogeneous polynomials with total degree $d_{1}, \ldots, d_{n}$. If $F$ is a regular sequence, then the degree of the ideal $\langle F\rangle$ is given by the Bézout bound:

$$
\operatorname{deg}(I)=\frac{\prod_{i=1}^{m} d_{i}}{\prod_{i=1}^{n} w_{i}}
$$

In particular, if $m=n$, then the dimension of $A /\langle F\rangle$ as $a \mathbb{K}$-vector space is either infinite or equal to $d_{1} \cdots d_{n} / w_{1} \cdots w_{n}$.

Proof. The first statement is a consequence of the formula for the Hilbert series of a regular sequence. For the second statement, from Props. 1.61 and $1.63, V(\langle F\rangle)$ has dimension 0 if and only if $F$ is a regular sequence. In that case, the result is a consequence of the first statement. If $V(\langle F\rangle)$ has positive dimension, the dimension of $A /\langle F\rangle$ is infinite.

All this can be applied in a non-homogeneous setting, by working with the homogenized system $F^{h}$ instead. However, in that case the Bézout bound is really a bound, not necessarily sharp. The reason for that is that $F^{h}$ may have more solutions than $F$. These solutions are called solutions at infinity and correspond to solutions with $H=0$ (where $H$ is the homogenization indeterminate). For example, two affine lines in a plane (degree 1 hypersurfaces) may not intersect if they are parallel.

Equivalently, solutions at infinity are points where the highest degree components of polynomials in $F$ simultaneously vanish. We shall encounter this behavior again when describing algorithms for polynomial system solving, in Section $2.4 \cdot 3$, and we will define regularity in the affine sense, in order to ensure that this phenomenon does not appear.

We conclude this section with the following geometrical interpretation of Noether position.

Proposition 1.65 ([Eis95, Cor. 9.3]). Assume that $\mathbb{K}$ is algebraically closed, and let $I \subset \mathbb{K}[\mathrm{X}]$ be an ideal in Noether position with respect to $X_{1}, \ldots, X_{m}$. Then the projection

$$
\begin{aligned}
\pi: & \mathbb{K}^{n} \\
\left(x_{1}, \ldots, x_{n}\right) & \longmapsto\left(\mathbb{K}^{n-m}\right. \\
& \longmapsto\left(x_{m+1}, \ldots, x_{n}\right)
\end{aligned}
$$

is surjective, and the fibers $\pi^{-1}(\mathrm{y})\left(\mathrm{y} \in \mathbb{K}^{n-m}\right)$ are finite.

### 1.3.6. Singularities

The study of singularities of algebraic varieties and morphisms is of first importance for classification problems. We shall see an example of such application in Section 1.3.7 (Theorem 1.82). For now, in this section, we give some definitions and useful theorems.

In the following section, we assume that $\mathbb{K}$ is algebraically closed.
Definition 1.66 (Jacobian matrix). Let $G=\left(g_{1}, \ldots, g_{m}\right): \mathbb{K}^{n} \longrightarrow \mathbb{K}^{m}$ be a polynomial map, and let $\mathrm{V}=\left\{V_{1}, \ldots, V_{k}\right\} \subset\left\{X_{1}, \ldots, X_{n}\right\}$. The facobian matrix of $G$ with respect to V is the matrix of the differential of $G$ :

$$
\operatorname{Jac}(G)_{\mathrm{V}}=\left[\begin{array}{ccc}
\frac{\partial g_{1}}{\partial V_{1}} & \cdots & \frac{\partial g_{1}}{\partial V_{n}} \\
\vdots & & \vdots \\
\frac{\partial g_{m}}{\partial V_{1}} & \cdots & \frac{\partial g_{m}}{\partial V_{n}}
\end{array}\right] .
$$

The Jacobian matrix of $G$ with respect to $\mathbf{X}$ is written $\operatorname{Jac}(G)$.
Definition 1.67 (Smoothness). Let $V$ be a variety in $\mathbb{K}^{n}$, and let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a system of generators of $I(V)$. Let $\mathbf{x} \in V$, and let $d$ be the local dimension of $V$ at $\mathbf{x}$. The point $\mathbf{x}$ is a regular point of $V$ if the $\operatorname{Jacobian}$ matrix $\operatorname{Jac}(F)(\mathrm{x})$ has rank $n-d$. Otherwise, x is a singular point of $V$.

The regular locus (resp. singular locus) of $V$ is the set $\operatorname{reg}(V)($ resp. $\operatorname{sing}(V))$ of all regular (resp. singular) points of $V$. If $\operatorname{sing}(V)=\varnothing$, we say that $V$ is smooth.
Remark 1.68. This definition of regularity requires $F$ to be a system of generators of $I(V)$. By the Nullstellensatz, it is equivalent to saying that $V=V(\langle F\rangle)$ and $\langle F\rangle$ is radical.
Definition 1.69 (Critical point). With the same notations as above, let $g=\left(g_{1}, \ldots, g_{k}\right)$ define a polynomial map $\mathbb{K}^{n} \longrightarrow \mathbb{K}^{k}$. We consider the vertical joint matrix

$$
\operatorname{Jac}(f, g)=\left[\begin{array}{l}
\operatorname{Jac}(f) \\
\operatorname{Jac}(g)
\end{array}\right] .
$$

Let $\mathbf{x} \in \operatorname{reg}(V)$. The point $\mathbf{x}$ is a critical point of the map $g$ restricted to $V$ if $\operatorname{Jac}(f, g)(\mathbf{x})$ has rank less than $k+n-d$. In that case, its image $g(\mathbf{x})$ is a critical value of $g$ restricted to $V$. A point $\mathrm{y} \in \mathbb{K}^{k}$ which is not a critical value of $g$ restricted to $V$ is a regular value of $g$ restricted to $V$.

The set of critical points of $g$ restricted to $V$ is denoted by crit $(G, V)$.
Remark 1.70. If $\mathbf{x} \in \operatorname{sing}(V)$, and $\mathbf{x}$ lies in a $d$-dimensional irreducible component of $V$, the rank of $\operatorname{Jac}(f, g)(\mathrm{x})$ is necessarily less than $k+n-d$. In other words, if $V$ is $d$-equidimensional, the locus at which $\operatorname{Jac}(f, g)(\mathrm{x})$ has rank less than $k+n-d$ is $\operatorname{sing}(V) \cup \operatorname{crit}(g, V)$.

Remark 1.71. Let $\mathbf{y} \in \mathbb{K}^{k}$; if $\mathbf{y}$ is a regular value of $g$ restricted to $V$, then $\left\langle f,\left(g_{i}-y_{i}\right)\right\rangle$ is radical and $g^{-1}(\mathbf{y}) \cap V \subset \mathbb{K}^{n}$ is smooth. In particular, if $0 \in \mathbb{K}^{k}$ is a regular value of $g$ restricted to $V$, then the variety $V(\langle f, g\rangle) \subset \mathbb{K}^{n}$ is smooth.

The Jacobian criterion is a very important tool linking the Jacobian matrix and properties of the underlying variety:

Theorem 1.72 (Jacobian criterion, [Nal15, Th. 1.9]). Let $V \in \mathbb{K}^{n}$ be an algebraic variety, and $f=\left(f_{1}, \ldots, f_{m}\right) \subset \mathbb{K}[\mathbf{X}]$ a system of generators of $I(V)$.

1. If $V$ is $d$-equidimensional, then for all $\mathbf{x} \in \mathbb{K}^{n}$, $\mathbf{x}$ is a regular point of $V$ if and only if $\operatorname{Jac}(f)(\mathbf{x})$ has rank $n-d$;
2. If for all $\mathbf{x} \in \mathbb{K}^{n}, \operatorname{Jac}(f)(\mathbf{x})$ has rank $n-d$, then $V$ is $d$-equidimensional and smooth.

We conclude this section with two powerful tools when working with critical values.
Sard's theorem states that given a polynomial map, a generic point of the target space is a regular value.

Theorem 1.73 (Sard's theorem, [Nal15, Lemma 1.16]). Let $V \subset \mathbb{K}^{n}$ be a variety, and let $g$ : $\mathbb{K}^{n} \longrightarrow \mathbb{K}^{m}$ be a polynomial map. Then $g(\operatorname{crit}(g, V))$ is contained in a proper hypersurface of $\mathbb{K}^{m}$.

Thom's transversality theorem states that given a smooth variety, the intersection with generic coordinate level lines is smooth.

Theorem 1.74 (Thom's weak transversality theorem, [Nal15, Th. 1.18],[Demoo, Th. 3.7.4]). Let $f: \mathbb{K}^{n} \times \mathbb{K}^{k} \longrightarrow \mathbb{K}^{r}$ be a polynomial map. For $\mathbf{y} \in \mathbb{K}^{k}$, define the partial application

$$
\begin{aligned}
f_{\mathbf{y}}: \mathbb{K}^{n} & \longrightarrow \mathbb{K}^{r} \\
\mathbf{x} & \longmapsto f(\mathbf{x}, \mathbf{y})
\end{aligned}
$$

Let $\mathcal{U} \subset \mathbb{K}^{n}$ be an open subset such that $0 \in \mathbb{K}^{r}$ is a regular value of $f$ restricted to $\mathcal{U} \times \mathbb{C}^{k}$. Then there exists an open subset $\mathcal{V} \subset \mathbb{K}^{k}$ such that for all $\mathbf{y} \in \mathcal{V}, 0 \in \mathbb{K}^{r}$ is a regular value of $f_{\mathbf{y}}$ restricted to $\mathcal{U}$.

### 1.3.7. Real semi-algebraic geometry

In this section, we focus on the case where the base field is $\mathbb{R}$. The definitions of the previous sections allow us to work with real algebraic sets (defined as subsets of $\mathbb{R}^{n}$ defined as the sets of zeroes of some polynomial ideal).

However, geometry over the reals is richer if we take into account the fact that $\mathbb{R}$ is an ordered field, by allowing inequalities in the definition of the basic objects of the geometry: that is the definition of a real semi-algebraic set.

Definition 1.75. A real semi-algebraic set is a finite union $V$ of subsets $V_{i}$, with each $V_{i}$ defined by a set of polynomial equations $\left\{F_{i}=0 \mid F_{i} \subset \mathbb{R}[\mathrm{X}]\right\}$ and a set of polynomial inequations $\left\{G_{i}>0 \mid G_{i} \subset \mathbb{R}[\mathbf{X}]\right\}$, as

$$
V_{i}=\left\{\begin{array}{l|l}
\mathbf{x} \in \mathbb{R}^{n} & \begin{array}{l}
\forall f \in F_{i}, f(\mathbf{x})=0 \\
\forall g \in G_{i}, g(\mathbf{x})>0
\end{array}
\end{array}\right\}
$$

Like algebraic varieties, real semi-algebraic sets are stable under finite union and intersection. Furthermore, they are also stable under complement. Finally, they are stable under projection onto linear subspaces:

Theorem 1.76 (Tarski-Seidenberg, [Nal15, Th. 1.13]). Let $V$ be a semi-algebraic subset of $\mathbb{R}^{n+1}$ and $\pi: \mathbb{R}^{n+1} \longrightarrow \mathbb{R}^{n}$ the projection onto the first $n$ coordinates. Then $\pi(V)$ is a semi-algebraic subset of $\mathbb{R}^{n}$.

Let $V$ be a real semi-algebraic set of $\mathbb{R}^{n}$. Locally (up to intersection with an open subset of $\mathbb{R}^{n}$ ), $V$ is a submanifold of $\mathbb{R}^{n}$, as defined by differential geometry:

Definition 1.77 (Submanifold). A submanifold $M$ is a subset of $\mathbb{R}^{n}$ such that, for all $\mathbf{x} \in M$, there exists an open subset $U \subset \mathbb{R}^{n}$ such that $M \cap U$ is diffeomorphic to $\mathbb{R}^{d}$, for some integer $d \in\{0, \ldots, n\}$. In this case, $d$ is the local dimension of $M$ at $\mathbf{x}$. The dimension of $M$ is the largest local dimension at a point of $M$.

Remark 1.78. Let $V$ be a real algebraic set. This definition of dimension coincides with the Krull dimension, as defined in Section 1.3.4.
Remark 1.79. One can define submanifolds of $\mathbb{C}^{n}$ in a similar way. Then complex algebraic varieties are locally submanifolds of $\mathbb{C}^{n}$. Moreover, the dimension of an algebraic variety $V$ seen as a submanifold is the Krull dimension of $V$, and if $V$ is $d$-equidimensional (for example irreducible), for all points at which $V$ is a submanifold, $V$ has local dimension $d$.

Remark 1.8o. Over the reals, it is no longer true: Whitney's umbrella (Figure 1.2) is a classic example of an irreducible subset where points may have different local dimensions. It is the algebraic set $W$ of $\mathbb{R}^{3}$ defined as the real zeroes of $X^{2}-Y^{2} Z$. This polynomial is irreducible (over $\mathbb{C}$ ), so the algebraic variety $W$ is irreducible.

However, $W$ is the disjoint union of two semi-algebraic sets:

$$
\begin{aligned}
& W_{+}=\left\{\begin{array}{l|l}
(x, y, z) & \begin{array}{l}
x^{2}=y^{2} z \\
z \geq 0
\end{array}
\end{array}\right\} \\
& W_{-}=\left\{\begin{array}{l|l}
(x, y, z) & \begin{array}{l}
x=y=0 \\
z<0
\end{array}
\end{array}\right\}
\end{aligned}
$$

The semi-algebraic set $W_{+}$has dimension 2 (the projection onto $(x, y)$ is a local diffeomorphism), and the semi-algebraic set $W_{-}$has dimension 1 (it is a half-line).

The final result of this section is an algebraic version of Ehresmann's theorem, or Thom's isotopy lemma, giving a description of fibers of the projection of a real semi-algebraic set under some hypotheses.

We first recall the following topological definition:
Definition 1.81 (Proper map). Let $A$ and $B$ be two topological spaces, and $f: A \longrightarrow B$ a map. The map $f$ is proper if and only if for all compact subsets $K \subset B, f^{-1}(K)$ is compact in $A$.

Theorem 1.82 (Thom's isotopy lemma, [CS95]). Let $V$ be a real semi-algebraic set of $\mathbb{R}^{n} \times \mathbb{R}^{t}$, and let $\pi: \mathbb{R}^{n} \times \mathbb{R}^{t} \longrightarrow \mathbb{R}^{t}$ be the projection onto the last $t$ coordinates. Assume that:


Figure 1.2: Whitney's umbrella

- $V$ is smooth;
- $V$ is locally closed;
- the restriction $\pi_{I V}$ has no critical points;
- $\pi_{I_{V}}$ is proper (for the euclidean topology);
- $V$ is $t$-equidimensional.

Then there exists a finite set $F$ and a semi-algebraic diffeomorphism

$$
h=\left(h_{0}, \pi\right): V \xrightarrow{\sim} F \times \mathbb{R}^{t} .
$$

### 1.4. Genericity

### 1.4.1. Definition

Genericity formalizes the intuitive notion that some property is almost always true.
Definition 1.83. Let $m$ be an integer, and $S=\left(S_{i}\right)_{i \in\{1, \ldots, m\}}$ a family of finite subsets of monomials of $A$. Let $A_{m, S}$ be the set of families of $m$ polynomials $f_{i} \in \mathbb{K}[\mathbf{X}]$, with $\operatorname{Supp}\left(f_{i}\right) \subset S_{i}$ for all $i \in\{1, \ldots, m\}$. Elements of $A_{k, S}$ can be seen as points in an affine space over $\overline{\mathbb{K}}$, the coordinates being the coefficients of the polynomials. Finally, let $P$ be a property of elements of $A_{m, s}$.

We say that $P$ is generic amongst elements of $A_{m, S}$ if the set of families satisfying $P$ in $A_{m, S}$ contains a non-empty Zariski-open subset of $A_{m, s}$.

A typical choice for the support $S_{i}$ is the set of polynomials with degree less than a given integer $d$, or homogeneous polynomials with degree exactly $d$.

### 1.4.2. Generic properties of homogeneous systems

We will now prove genericity results for the regularity properties defined in Section 1.1.5, for homogeneous systems (with the total degree). More generally, the proofs that these properties cover Zariski-open subsets hold for any gradation, but for non-trivial systems of weights, these subsets need not be non-empty (see Section 3.2.4).
We assume that $A$ is graded in total degree. For $d \in \mathbb{N}$, let $A_{d}$ be the set of homogeneous generic polynomials with degree exactly $d$.

Proposition 1.84 ([Baro4, Prop. 1.7.4(5)], [Par10]). Let $m \in\{1, \ldots, n\}$ and let $\left(d_{1}, \ldots, d_{m}\right) \in$ $\mathbb{N}^{m}$ be a family of degrees. Homogeneous regular sequences with degree $\left(d_{1}, \ldots, d_{m}\right)$ are generic amongst families of homogeneous polynomials of degree $\left(d_{1}, \ldots, d_{m}\right)$, that is in $\prod_{i=1}^{m} A_{d_{i}}$.

Proof. Assume that $m=n$ and let $F=\left(f_{1}, \ldots, f_{m}\right) \in \mathbb{K}[\mathbf{A}][\mathbf{X}]$ be a sequence of polynomials in $\mathbf{X}$, with generic coefficients A. Let $I=\langle F\rangle$. Write

$$
H_{\mathrm{reg}}=\frac{\prod_{i=1}^{n}\left(1-T^{d_{i}}\right)}{(1-T)^{n}}=\sum_{i=1}^{\infty} a_{i} T^{i} \text { with } a_{\delta+1}=a_{\delta+2}=\cdots=0
$$

and let $b_{i}$ be the coefficient of degree $i$ in the Hilbert series of $F$.
By Prop. 1.38, for all $i \in \mathbb{N}, b_{i} \geq a_{i}$ and $F$ is regular if and only if they are equal. For each $i$, this equality means that the rank of the Macaulay matrix of $F$ at degree $i$ is $a_{i}$, and this can be encoded by the non-vanishing of the $a_{i} \times a_{i}$-minors of this matrix. Let $U_{i}$ be the Zariski-open subset defined by the non-vanishing of at least one of these minors.
Furthermore, if $b_{\delta+1}=0$, it means that $I$ contains all monomials of degree $\delta+1$, and it implies that $b_{\delta+2}=\cdots=0$. So we only need to check a finite number of equalities $a_{i}=b_{i}$, for $i$ ranging from 0 to $\delta+1$. Geometrically, it means that the set of regular sequences is the finite intersection $U$ of the open sets $U_{i}$.

Consider the sequence ( $X_{1}^{d_{1}}, \ldots, X_{n}^{d_{n}}$ ), it is regular, and so $U$ is non-empty.
If $m<n, F$ is regular if and only if there exist linear forms $l_{m+1}, \ldots, l_{n}$ such that the sequence $\left(f_{1}, \ldots, f_{m}, l_{m+1}, \ldots, l_{n}\right)$ is regular. In particular, the set of regular sequences of homogeneous polynomials with degree $\left(d_{1}, \ldots, d_{m}\right)$ contains the set of sequences $F=\left(f_{1}, \ldots, f_{m}\right)$ such that $\left(f_{1}, \ldots, f_{m}, X_{m+1}, \ldots, X_{n}\right)$ is regular. From the above, regular sequences with degree $\left(d_{1}, \ldots, d_{m}, 1, \ldots, 1\right)$ form a non-empty Zariski-open subset $U$ of $\prod_{i=1}^{m} A_{d_{i}} \times A_{1}^{n-m}$, and the closed subset $C$ defined by $f_{m+1}=X_{m+1}, \ldots, f_{n}=X_{n}$ does meet $U$. Consider the projection onto the $m$ first components

$$
\pi: \prod_{i=1}^{m} A_{d_{i}} \times A_{1}^{n-m} \longrightarrow \prod_{i=1}^{m} A_{d_{i}}
$$

It is surjective, and its restriction to $C$ is an isomorphism. Since $C \cap U$ is non-empty, $\pi(C \cap U)$ is a non-empty open subset of $\prod_{i=1}^{m} A_{d_{i}}$, defining exactly those sequences $F=\left(f_{1}, \ldots, f_{m}\right)$ such that $\left(f_{1}, \ldots, f_{m}, X_{m+1}, \ldots, X_{n}\right)$ is regular.

Corollary 1.85. Homogeneous sequences in (simultaneous) Noether position with respective degree ( $d_{1}, \ldots, d_{m}$ ) are generic amongst families of homogeneous polynomials of degree $\left(d_{1}, \ldots, d_{m}\right)$.

Proof. This is a direct consequence of the proof of the case $m<n$ of Prop. 1.84.
Corollary 1.86. Letm be an integer, and let $\left(d_{1}, \ldots, d_{m}\right)$ be a family of degrees. Assume that there exists at least one semi-regular sequence with degree $\left(d_{1}, \ldots, d_{m}\right)$, then homogeneous semi-regular sequences with degree $\left(d_{1}, \ldots, d_{m}\right)$ are generic amongst families of homogeneous polynomials of degree $\left(d_{1}, \ldots, d_{m}\right)$.

Proof. If $m \leq n$, semi-regular sequences are actually regular, and so Prop. 1.84 applies. Otherwise, semi-regular sequences define zero-dimensional ideals and are characterized by their Hilbert series, so the proof that zero-dimensional regular sequences form a Zariski-open subset in Prop. 1.84 can be repeated for semi-regular sequences.

This corollary states that semi-regular sequences form a Zariski-open subset of all sequences, and that the property is generic if and only if this Zariski-open subset is non-empty. This is an open question:

Conjecture 1.87 (Fröberg, [Frö85]). Let $m$ be an integer, and let $\left(d_{1}, \ldots, d_{m}\right)$ be a family of degrees. Then there exists at least one semi-regular sequence of homogeneous polynomials with respective degrees $d_{1}, \ldots, d_{m}$.

This conjecture is only proved in a handful of cases [Moro3, Th. 1.5]:

- $m \leq n$ : then semi-regular sequences are regular sequences;
- $n=2$;
- $n=3$ and the field $\mathbb{K}$ is infinite;
- $m=n+1$ and the field $\mathbb{K}$ has characteristic 0 .


### 1.4.3. Generic changes of coordinates

Some properties are also true up to a generic change of coordinates: given a property $P$, this means that given any system $F$, the set of linear changes of coordinates $L$ such that $F \circ L$ satisfies $P$ is generic amongst all linear changes of coordinates.

A classic result of this kind is Noether's normalization lemma: it is a fundamental result of dimension theory, stating that given a variety $V$ of dimension $d$, up to a generic change of coordinates, the projection $V \longrightarrow \mathbb{K}^{n-m}$ is surjective and finite.

We give here an algebraic transcription of this theorem for a regular sequence:
Theorem 1.88 (Noether's normalization theorem, [Eis95, Th. 13.3]). Let $\mathbb{K}$ be an infinite field, $F$ a regular sequence in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$, and $k \in\{1, \ldots, n\}$. Then, for a generic choice of $n-k$ linear forms $l_{i}\left(X_{i+1}, \ldots, X_{n}\right), i \in\{1, \ldots, k\}$, the change of variables

$$
X_{i}=X_{i}^{\prime}+l_{i}\left(X_{i+1}, \ldots, X_{n}\right)
$$

is such that $F\left(X_{1}\left(\mathrm{X}^{\prime}\right), \ldots, X_{k}\left(\mathrm{X}^{\prime}\right), X_{k+1}, \ldots, X_{n}\right)$ is in Noether position with respect to the variables $X_{k+1}, \ldots, X_{n}$.

More generally, if $F$ has length $n$, then for a generic choice of $n$ linear forms $l_{i}\left(X_{i+1}, \ldots, X_{n}\right)$, $i \in\{1, \ldots, n\}$, the change of variables

$$
X_{i}=X_{i}^{\prime}+l_{i}\left(X_{i+1}, \ldots, X_{n}\right)
$$

is such that $F\left(X_{1}\left(\mathrm{X}^{\prime}\right), \ldots, X_{n}\left(\mathrm{X}^{\prime}\right)\right)$ is in simultaneous Noether position with respect to the order $X_{1}^{\prime}>\cdots>X_{n}^{\prime}$.

Noether's normalization theorem is really made of two parts: first, up to a substitution $X_{i}=P(\mathbf{X})$, any system can be put in Noether position; and second, that under some hypotheses, we can make additional assumptions on the substitution. Here, the second part is that if $\mathbb{K}$ is infinite, then any generic linear change of variables gives a suitable transformation; this part of the proof of [Eis95, Th. 13.3] relies on the following lemma:

Lemma 1.89 (Noether's normalization lemma, [Eis95, lem. 13.2.c]). Let $\mathbb{K}$ be an infinite field and $f \in A=\mathbb{K}\left[X_{1}, \ldots, X_{r}\right]$ a non-constant polynomial. Then there are homogeneous elements $X_{1}^{\prime}, \ldots, X_{r-1}^{\prime} \in A$ with degree 1 such that $A$ is a finitely generated module over $\mathbb{K}\left[X_{1}^{\prime}, \ldots, X_{r-1}^{\prime}, f\right]$. Furthermore, there exists a dense Zariski-open subset $U \subset \mathbb{K}^{r-1}$ such that for all $\left(a_{i}\right) \in U$, one can choose $X_{i}^{\prime}=X_{i}-a_{i} X_{r}$.

Proof. For any $1 \leq i \leq r-1$, let $a_{i} \in \mathbb{K}$, and let $X_{i}^{\prime}=X_{i}-a_{i} X_{r}$. If $f$ has degree $d$, let $f_{d}$ be the degree $d$ component of $f$. We need to show that for generic $a_{i}$, under this change of variables, $f$ is monic in $X_{r}$, and it is enough to prove it for $f_{d}$ :

$$
\begin{aligned}
f_{d}\left(X_{1}, \ldots, X_{r}\right) & =f_{d}\left(X_{1}^{\prime}+a_{1} X_{r}, X_{2}^{\prime}+a_{2} X_{r}, \ldots, X_{r-1}+a_{r-1} X_{r}\right) \\
& =f_{d}\left(a_{1}, \ldots, a_{r-1}, 1\right) X_{r}^{d}+\ldots
\end{aligned}
$$

So the set of all $a_{i}$ 's such that $f$ is monic in $X_{r}$ is exactly the set of all $a_{i}$ 's such that

$$
f_{d}\left(a_{1}, \ldots, a_{r-1}, 1\right) \neq 0
$$

and since $f$ is homogeneous non-constant, this is a non-empty open subset of $\mathbb{K}^{r-1}$.
If $F$ is homogeneous, it is still homogeneous after any linear change of variables, so Noether's normalization theorem states that a regular homogeneous sequence can be made into a homogeneous sequence in Noether position, up to a generic change of variables. We shall identify suitable changes of variables for the same property to hold in a weighted setting in Section 3.2.3.

### 1.5. Determinantal varieties

In this section, we will define determinantal varieties, that is varieties defined as the locus of rank defects of a matrix with polynomial entries. These systems arise in many applications, for example in optimization, when using the Jacobian criterion to model singularities and critical points in terms of the minors of the Jacobian matrix.

### 1.5.1. Definition

Definition 1.90 (Determinantal ideal, determinantal variety). Let $r, k, n$ be integers such that $r \leq k$. Let $M \in \mathbb{K}\left[X_{1}, \ldots, X_{n}\right]^{k \times k}$ be a matrix with polynomial coefficients. The $r$ 'th determinantal ideal associated with $M$ is the ideal $\mathcal{D}_{r}$ generated by its $(r+1)$-minors. The corresponding determinantal variety is

$$
V_{r}=\left\{\mathbf{x} \in \mathbb{K}^{n} \mid \operatorname{rank}(M(\mathbf{x})) \leq r\right\}
$$

In Chapter 4, we shall be interested in the study of the singularities of these varieties. Because of the structure of determinantal ideals, these singularities are strongly related to other determinantal varieties.
Proposition 1.91. If $\mathcal{D}_{r}$ is radical, then

$$
V_{r-1} \subset \operatorname{sing}\left(V_{r}\right)
$$

Proof. We will prove the following stronger statement: let $\mathbf{x} \in V_{r-1}$, then all partial derivatives of all $(r+1)$-minors of $M$ vanish at $\mathbf{x}$. This will prove that the Jacobian of $V_{r}$ at $\mathbf{x}$ is the zero matrix, and in particular has rank 0 . Then, by hypothesis the ideal of all $(r+1)$-minors of $M$ is radical, so we can use the Jacobian criterion to characterize $\operatorname{sing}(V)$, and we can conclude that $\mathbf{x} \in \operatorname{sing}\left(V_{r}\right)$.

To prove this statement, let $\mathbf{x} \in V_{r-1}$. For any $(r+1) \times(r+1)$-submatrix of $M$, the result we want to prove only depends on the coefficients of the submatrix. So w.l.o.g., we may assume that $r=k-1$, let $D=\operatorname{det}(M)$. Let $\mathbf{x} \in V_{k-2}$, this means that all $(k-1)$-minors of $M$ vanish at $\mathbf{x}$. Consider a matrix of polynomial indeterminates $\mathbf{U}$ :

$$
\tilde{M}=\left[\begin{array}{ccc}
U_{1,1} & \ldots & U_{1, k}  \tag{1.2}\\
\vdots & & \vdots \\
U_{k, 1} & \ldots & U_{k, k}
\end{array}\right]
$$

and let $\tilde{D}$ be its determinant. Then for any $i, j \in\{1, \ldots, k\}$,

$$
\frac{\partial \tilde{D}}{\partial U_{i, j}}=(-1)^{i+j} \cdot \tilde{\mathcal{M}}_{i, j}(\mathbf{U})
$$

where $\tilde{\mathcal{M}}_{i, j}$ is the $(k-1)$-minor of $\tilde{M}$ obtained by removing row $i$ and column $j$.
For any $i, j \in\{1, \ldots, k\}$, let $m_{i, j}$ (resp. $\mathcal{M}_{i, j}$ ) be the coefficient at row $i$ and column $j$ of the matrix $M$ (resp. the minor obtained by removing row $i$ and column $j$ from $M$ ). By the derivation chain rule, for any $X \in\left\{X_{1}, \ldots, X_{n}\right\}$,

$$
\begin{aligned}
\frac{\partial D}{\partial X} & =\sum_{i, j=1}^{k}(-1)^{i+j} \frac{\partial m_{i, j}}{\partial X} \cdot \tilde{\mathcal{M}}_{i, j}(\mathbf{m}) \\
& =\sum_{i, j=1}^{k}(-1)^{i+j} \frac{\partial m_{i, j}}{\partial X} \cdot \mathcal{M}_{i, j}(\mathbf{X})
\end{aligned}
$$

Since by hypothesis all $(k-1)$ minors of $M$ vanish at $\mathbf{x}$, all partial derivatives $\frac{\partial \Delta}{\partial X}$ vanish at x.

Generically, this inclusion is an equality.
Proposition 1.92 ([HNS15b, Prop. 2]). Let $d \in \mathbb{N}$. Generically amongst matrices with coefficients of degree at most $d$ :

1. $V_{r-1}=\operatorname{sing}\left(V_{r}\right)$;
2. the variety $V_{r}$ is equidimensional with codimension $c=(k-r)^{2}$ if $n \geq(k-r)^{2}$, and empty otherwise;
3. the ideal $D_{r}$ is radical.

Proof. Let $\tilde{D}_{r}$ be the $r$ 'th determinantal ideal associated to the matrix $\tilde{M}$ as defined in Equation (1.2), and, as before, let $m_{i, j}$ be the coefficient at row $i$ and column $j$ in $M$. In $\mathbb{K}[\mathbf{X}, \mathbf{U}]$, the two determinantal ideals are related as

$$
\left\langle D_{r}\right\rangle=\left\langle\tilde{D}_{r}\right\rangle+\left\langle U_{i, j}-m_{i, j}(\mathbf{X})\right\rangle_{1 \leq i, j \leq k}
$$

Now, all three statements of the proposition are unconditionally true for $\tilde{M}$ : see [BV88, Th. 2.10] for the radicality of $\tilde{D}_{r}$, and then [BV88, Prop. 1.1] for the other two statements. We shall prove that generically, these statements transfer to $M$.

Consider the application

$$
\begin{aligned}
\mathcal{M}_{\mathrm{m}}: \mathbb{K}^{n} & \longrightarrow \mathbb{K}^{k^{2}} \\
\mathbf{x} & \longmapsto\left(m_{i, j}(\mathbf{x})\right)_{i, j}
\end{aligned}
$$

By Thom's weak transversality lemma, generically over the space of polynomials ( $m_{i, j}$ ),

$$
V_{r} \backslash V_{r-1}=\mathcal{M}_{\mathrm{m}}^{-1}\left(\tilde{V}_{r} \backslash V_{r-1}\right)
$$

is either empty or smooth and equidimensional with codimension $c=(k-r)^{2}$. This proves statements 1 and 2. It also proves that $V_{r}=V\left(D_{r}\right)$ satisfies the hypotheses of [Eis95, Th. 18.18].

Furthermore, add the substitution equations $U_{i, j}-m_{i, j}(\mathbf{X})$ to the generators of $\tilde{D_{r}}$ in $\mathbb{K}[\mathbf{U}, \mathbf{X}]$, and consider the Jacobian of this system. It writes as a block matrix

$$
\operatorname{Jac}_{\mathbf{U}, \mathbf{X}}\left(D_{r}\right)=\left[\begin{array}{cc}
\operatorname{Jac}_{\mathbf{U}}\left(\tilde{D}_{r}\right) & 0 \\
\operatorname{Id}_{k^{2}} & \operatorname{Jac}_{\mathbf{X}}\left(m_{i, j}\right)
\end{array}\right]
$$

where the top-left block has $\binom{k}{r}$ rows and $k^{2}$ columns, and the bottom-right block has $k^{2}$ rows and $n$ columns. By Sard's lemma, the set of $\mathbf{x}$ such that $\operatorname{Jac}_{\mathbf{X}}\left(m_{i, j}\right)$ has non-full rank has codimension at least 1 in $\mathbb{K}^{n}$; and from the above, outside of $\tilde{V_{r-1}}, \mathrm{Jac}_{\mathrm{U}}\left(\tilde{D}_{r}\right)$ has rank $c$. So by Serre's Criterion for radicality ([Eis95, Th. 18.15(a) with Th. 18.18]), under the same genericity condition as above, the ideal $D_{r}$ is radical, which completes the proof.

### 1.5.2. Cramer's formula and Schur's complement

In particular, results from the previous section imply that the system of all $(r+1)$-minors of $M$ is not a regular sequence (it has length $\binom{k}{r+1}^{2}$ and, generically, defines a variety with codimension $\left.(k-r)^{2}\right)$. However, it is possible to cover determinantal varieties with dense open subsets, each being an open subset of a complete intersection.

These constructions stem from the linear algebraic interpretation of the rank: a $k \times k$ matrix has rank at most $r$ if and only if its kernel has dimension at least $k-r$. Algebraically, the correspondence between linear algebra and minors comes from Cramer's formula, and its generalization Shur's complement. Cramer's formula is a classical formula for solving invertible linear system. Here, we give an equivalent statement for matrices with corank 1.

Proposition 1.93 (Cramer's formula). Let $R$ be a ring, $k \in \mathbb{N}, M \in R^{k \times k}$ an invertible matrix, $b \in R^{k}$ a vector, and consider the linear system in $x={ }^{\top}\left(x_{1}, \ldots, x_{k}\right)$ :

$$
M x=b
$$

Then this system has a unique solution, defined by

$$
\forall i \in\{1, \ldots, k\}, x_{i}=\frac{\operatorname{det}\left(M_{i}\right)}{\operatorname{det}(M)}
$$

where $M_{i}$ is the matrix obtained by replacing the column $i$ in $M$ with $b$.
Amongst its consequences, we recall the classic formula for the inverse of a matrix:
Proposition 1.94. With the same notations and hypotheses,

$$
\begin{equation*}
M^{-1}=\frac{1}{\operatorname{det}(M)}^{\top} \operatorname{Com}(M) \tag{1.3}
\end{equation*}
$$

where $\operatorname{Com}(M)$ is the cofactor matrix of $M$, whose coefficient at row $i$ and column $j$ is

$$
c_{i, j}=(-1)^{i+j} \cdot(k-1) \text {-minor of } M \text { obtained by removing row } i \text { and column } j .
$$

Conversely, if $M$ is such that $\operatorname{det}(M) \neq 0$, then $M$ is invertible in the localized ring $R_{\operatorname{det}(M)}=$ $R[1 / \operatorname{det}(M)]$, and its inverse is given by the same formula.

Cramer's rule can be used to parameterize the kernel of non-invertible matrices. For example, assume that $M$ is a $k \times k$ singular matrix such that the submatrix $M_{0}$ obtained by removing row $k$ and column $k$ is invertible, so that $M$ has rank $k-1$. Let $x={ }^{\top}\left(x_{1}, \ldots, x_{k}\right) \in \operatorname{Ker}(M)$ :

$$
M \cdot x=0
$$

Keeping only the $k-1$ first rows of $M$, the system becomes

$$
M_{0} \cdot\left[\begin{array}{c}
x_{1} \\
\vdots \\
x_{k-1}
\end{array}\right]=-x_{k} C_{k}
$$

where $C_{k}$ is the $k^{\prime}$ th column in $M$. Cramer's rule implies that the kernel of $M$ is parameterized by $x_{k}$ as

$$
\forall i \in\{1, \ldots, k-1\}, x_{i}=\frac{-x_{k} \operatorname{det}\left(M_{i}\right)}{\operatorname{det}\left(M_{0}\right)}
$$

where $M_{i}(i \in\{1, \ldots, k-1\})$ is the submatrix obtained by replacing column $i$ in $M_{0}$ with the last column $C_{k}$ of $M$, less row $i$.

This construction is what we use to fill the gap between determinantal ideals and kernels of matrices. For matrices of smaller rank, it can be generalized:
Definition 1.95 (Schur complement). Let $M \in R^{k \times k}$ be a square matrix, $p \in\{1, \ldots, k\}$, and $A$, $B, C$ and $D$ be respectively $p \times p, p \times k-p, k-p \times p$ and $k-p \times k-p$ matrices such that

$$
M=\left[\begin{array}{ll}
A & B \\
C & D
\end{array}\right]
$$

Further assume that $A$ is invertible. Then the Schur complement of the block $A$ in $M$ is the $k-p \times k-p$ matrix

$$
\begin{equation*}
M / A=D-C A^{-1} B . \tag{1.4}
\end{equation*}
$$

Proposition 1.96. The Schur complement of $A$ in $M$ appears in the following block diagonal decomposition of $M$ :

$$
\left[\begin{array}{ll}
A & B \\
C & D
\end{array}\right]=\left[\begin{array}{cc}
\mathrm{Id}_{p} & 0 \\
C A^{-1} & \mathrm{Id}_{k-p}
\end{array}\right] \cdot\left[\begin{array}{cc}
A & 0 \\
0 & M / A
\end{array}\right] \cdot\left[\begin{array}{cc}
\mathrm{Id}_{p} & A^{-1} B \\
0 & \mathrm{Id}_{k-p}
\end{array}\right]
$$

Proposition 1.97. The entry at row $i$ and column $j$ in $M / A$ is $\mathcal{M}_{p, i, j} / \operatorname{det}(A)$, where $\mathcal{M}_{p, i, j}$ is the ( $p+1$ )-minor of $M$ obtained by taking rows $1,2, \ldots, p, p+i$ and columns $1,2, \ldots, p, p+j$, or equivalently, the $(p+1)$-minor of $M$ obtained by appending row $i$ and column $j$ to $A$.
Proof. Let $a_{i, j}, a_{i, j}^{\prime}, b_{i, j}, c_{i, j}, d_{i, j}, d_{i, j}^{\prime}$ be the coefficient at row $i$ and column $j$ in respectively $A$, $A^{-1}, B, C, D, M / A$. Formula (1.4) yields that

$$
d_{i, j}^{\prime}=d_{i, j}-\sum_{u=1}^{p} c_{i, u} \sum_{v=1}^{p} a_{u, v}^{\prime} b_{v, j}
$$

By Formula (1.3),

$$
a_{u, v}^{\prime}=\frac{(-1)^{u+v} \mathcal{A}_{v, u}}{\operatorname{det}(A)}
$$

where $\mathcal{A}_{v, u}$ is the $(p-1)$-minor of $A$ obtained by removing row $v$ and column $u$ ( $u$ and $v$ are swapped because of the transposition in Formula (1.3)). So, by decomposition over column $j$, $\operatorname{det}(A) \sum_{v=1}^{p} a_{u, v}^{\prime} b_{v, j}$ is the $p$-minor of $M$ obtained by removing column $u$ in $A$ and appending column $j$ in $B$, and by decomposition over row $i$,

$$
\operatorname{det}(A) d_{i, j}^{\prime}=\operatorname{det}(A) d_{i, j}-\operatorname{det}(A) \sum_{u=1}^{p} c_{i, u} \sum_{v=1}^{p} a_{u, v}^{\prime} b_{v, j}
$$

is the $(p+1)$-minor of $M$ obtained by adding row $p+i$ and column $p+j$ in $M$ to $A$.

To conclude this section, we give an immediate consequence of these last two results, which is a first local description of determinantal systems.

Proposition 1.98. Assume that $V_{r-1}$ has empty interior in $V_{r}$. Let A be a $r \times r$ submatrix of $M$. Then the open subset $W$ of $V_{r}$ at which $\operatorname{det}(A)$ does not vanish is defined by the vanishing of all $(r+1)$-minors of $M$ containing $A$, and generically, they form a regular sequence.

Proof. Up to permutation of the rows and columns, we may assume that $A$ is the top-left $r \times r$ submatrix of $M$. If $A$ is invertible, $M$ has rank at least $r$, and by Proposition 1.96, $M$ has rank $r$ if and only if $M / A$ is 0 . By Proposition 1.97 , the entries of this matrix are precisely the $(r+1)$ minors containing $A$.

Since generically, the variety $V_{r}$ is equidimensional with dimension $(k-r)^{2}$, the same holds for the non-empty open subset $W$. The sequence of $(r+1)$-minors containing $A$ has length $(k-r)^{2}$, so by the characterization of Propostion 1.63 , they form a regular sequence.

### 1.5.3. Incidence varieties

Incidence varieties offer an alternative modelling of determinantal varieties.
Definition 1.99. Let $r \in\{0, \ldots, k-1\}$. The incidence variety of rank $r$ associated with $M$ is the variety $\mathcal{V}_{r} \subset \mathbb{K}^{n} \times\left(\mathbb{P}^{k-1}(\mathbb{K})\right)^{k-r}$ defined by:

$$
M \cdot\left[\begin{array}{ccc}
Y_{1,1} & \ldots & Y_{1, k-r}  \tag{1.5}\\
\vdots & & \vdots \\
Y_{k, 1} & \ldots & Y_{k, k-r}
\end{array}\right]=\left[\begin{array}{ccc}
0 & \ldots & 0 \\
\vdots & & \vdots \\
0 & \ldots & 0
\end{array}\right]
$$

with the additional condition that the matrix $\left(Y_{i, j}\right)$ has rank $k-r$.
The projection of $\mathcal{V}_{r}$ onto the affine space with coordinates $(\mathbf{X})$ is $V_{r}$. For

$$
\mathbf{u}=\left(u_{1,1}, \ldots, u_{k-r, k}\right) \in \mathbb{K}^{k(k-r)},
$$

we define the variety $\mathcal{V}_{r, \mathbf{u}}^{\prime}$ as the intersection of $\mathcal{V}_{r}$ and the complex solutions of

$$
\left[\begin{array}{ccc}
u_{1,1} & \ldots & u_{1, k}  \tag{1.6}\\
\vdots & & \vdots \\
u_{k-r, 1} & \ldots & u_{k-r, k}
\end{array}\right] \cdot\left[\begin{array}{ccc}
Y_{1,1} & \ldots & Y_{1, k-r} \\
\vdots & & \vdots \\
Y_{k, 1} & \ldots & Y_{k, k-r}
\end{array}\right]=\operatorname{Id}_{k-r}
$$

Proposition 1.100. For any $r \in\{0, \ldots, k-1\}$, the varieties $V_{r}$ and $\mathcal{V}_{r, \mathbf{u}}^{\prime}$ are birational.
Proof. We need to define a morphism $f: W \longrightarrow \mathcal{V}_{r, \mathbf{u}}^{\prime}$ with $W$ a non-empty Zariski-open subset of $V_{r}$, and such that $f$ is inverse to the projection $\mathcal{V}_{r, \mathbf{u}}^{\prime} \longrightarrow V_{r}$ onto the affine space with coordinates (X). Let $W$ be the open subset of $V_{r}$ defined as the non-vanishing locus of the top-left $r$-minor of $M$. Consider the block decompositions, where $A, Y_{(1)}$ and $U_{(1)}$ are $r \times r$ matrices:

$$
M=\left[\begin{array}{ll}
A & B \\
C & D
\end{array}\right] \quad Y=\left[\begin{array}{l}
Y_{(1)} \\
Y_{(2)}
\end{array}\right] \quad U=\left[\begin{array}{ll}
U_{(1)} & U_{(2)}
\end{array}\right]
$$

Over $W, A$ is invertible, let $\Delta=\operatorname{det}(A)$. Let $M / A$ be the Schur complement of $A$ in $M$, Equations (1.5) and (1.6) can be rewritten as

$$
\left[\begin{array}{cc}
\Delta \mathrm{Id}_{r} & A^{-1} B \\
0 & M / A \\
U_{(1)} & U_{(2)}
\end{array}\right] \cdot\left[\begin{array}{c}
Y_{(1)} \\
Y_{(2)}
\end{array}\right]=\left[\begin{array}{c}
0 \\
0 \\
\mathrm{Id}_{r}
\end{array}\right]
$$

We may restrict to the open subset of $\mathcal{V}_{r, \mathbf{u}}^{\prime}$ where $Y_{(2)}$ is invertible, then eliminating $Y_{(1)}$ yields that

$$
\left\{\begin{array}{l}
Y_{(2)}=\left(U_{(2)}-U_{(1)} A^{-1} B\right)^{-1} \\
Y_{(1)}=\frac{-1}{\Delta} A^{-1} B Y_{(2)}
\end{array}\right.
$$

which defines the wanted morphism $W \longrightarrow \mathcal{V}_{r, \mathbf{u}}^{\prime}$.

## Chapter 2

## Gröbner bases

This chapter deals with Gröbner bases, and gives definitions and properties which will be used throughout the thesis. It is a powerful tool for solving polynomial systems, and is defined as a particular set of generators of a polynomial ideal. This set of generators depends on the choice of an order on the monomials, and this choice has consequences on what information does the basis yield, and how complicated it is to compute the basis.

This chapter is organized as follows: in Section 2.1, we define monomial orderings, and list a few examples which shall be useful in the rest of the thesis. In Section 2.2, we recall the definition of a Gröbner basis and a normal form, and in Section 2.3 we examine two common applications of Gröbner bases (namely listing all the solutions of a system with a finite-number of solutions and computing equations of the Zariski-closure of the projection of a variety) and we use these examples to illustrate how each application leads to the choice of a specific monomial order.

Then we move on to more algorithmic considerations. In Section 2.4, we describe 3 algorithms for computing Gröbner bases: the historical algorithm from Buchberger, the more modern $F_{5}$ algorithm through its variant Matrix- $F_{5}$, and the FGLM algorithm, for change of order. Section 2.5 concludes this chapter with complexity results for algorithms Matrix- $\mathrm{F}_{5}$ and FGLM.

This chapter, like the previous one, is purely bibliographical.

### 2.1. Monomial orderings

### 2.1.1. Definition

Gröbner bases generalize euclidean division to multivariate polynomials, and gaussian elimination to nonlinear polynomials. In both cases, we eliminate some monomials until we reach a normal form: in the univariate case, we eliminate high degree monomials; in the linear case, we choose an order on the indeterminates, and we eliminate them in that order.

Definition 2.1 (Monomial ordering). Let $A=\mathbb{K}[\mathrm{X}]$ be a polynomial algebra with $n$ indeterminates; let $\mathcal{M}$ be the set of its monomials. A monomial ordering is a total ordering $<_{\text {mon }}$ on $\mathcal{M}$ such that:

$$
\begin{aligned}
& \forall x \in \mathcal{M}, 1<_{\operatorname{mon}} x \\
& \forall x, y, z \in \mathcal{M}, x<_{\operatorname{mon}} y \Longrightarrow x z<_{\operatorname{mon}} y z
\end{aligned}
$$

Equivalently, a monomial ordering can be defined on the set of exponents $\mathbb{N}^{n}$, in which case the requirements are that:

$$
\begin{aligned}
& \forall v \in \mathbb{N}^{n}, 0<_{\operatorname{mon}} v \\
& \forall u, v, w \in \mathbb{N}^{n}, u<_{\operatorname{mon}} v \Longrightarrow u+w<_{\operatorname{mon}} v+w
\end{aligned}
$$

Given a monomial ordering, we may define:
Definition 2.2 (Leading monomial, leading term, leading coefficient). Let $f \in \mathbb{K}[\mathbf{X}]$, and let $<_{\text {mon }}$ be a monomial ordering. Then

- the leading monomial $\mathrm{LM}_{<_{\operatorname{mon}}}(f)$ of $f$ is the largest monomial of its support ;
- the leading coefficient $\mathrm{LC}_{<_{\text {mon }}}(f)$ of $f$ is the coefficient of $\mathrm{LM}(f)$;
- the leading term $\mathrm{LT}_{<_{\text {mon }}}(f)$ of $f$ is the product $\mathrm{LC}_{<_{\text {mon }}}(f) \mathrm{LM}_{<_{\text {mon }}}(f)$.

Let $I$ be an ideal of $\mathbb{K}[\mathrm{X}]$, the leading ideal (or initial ideal) of $I$ is the ideal $\left\langle\mathrm{LT}_{<_{\text {mon }}}(I)\right\rangle$ generated by the leading terms of all elements of $I$.

For all these notions, we shall omit the subscripted monomial order when clear by the context.
In the multivariate nonlinear case, unlike the linear and the univariate cases, the choice of an order on the monomials is not unique.

In the rest of this section, we give several examples of monomial orderings, with basic algebraic properties. We shall examine them again in Section 2.3, from the angle of applications. Let $\alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right)$ and $\beta=\left(\beta_{1}, \ldots, \beta_{n}\right)$ be two exponents in $\mathbb{N}^{n}$.

### 2.1.2. Lexicographical ordering

The lexicographical ordering is maybe the most "natural" monomial order: it is the dictionary order for the "words" $\alpha$ and $\beta$.

Definition 2.3. The lexicographical ordering Lex is defined as:

$$
\alpha<_{\operatorname{lex}} \beta \Longleftrightarrow \alpha_{1}=\beta_{1}, \ldots, \alpha_{j-1}=\beta_{j-1}, \alpha_{j}<\beta_{j}
$$

Its main property is the following:
Proposition 2.4. Let $f \in \mathbb{K}[\mathrm{X}]$. If $L T_{<\text {lex }}(f) \in \mathbb{K}\left[X_{i}, \ldots, X_{n}\right]$ for $i<n$, then $f \in \mathbb{K}\left[X_{i}, \ldots, X_{n}\right]$.
Proof. Let $\left(\alpha_{1}, \ldots, \alpha_{n}\right)$ be the exponents of $\operatorname{LT}(f)$ for the lexicographical order. If $\operatorname{LT}(f) \in$ $\mathbb{K}\left[X_{i}, \ldots, X_{n}\right]$, then $\alpha_{1}=\cdots=\alpha_{i-1}=0$. Any monomial $X^{\beta}$ such that $\beta_{j}>0$ with $j<i$ is larger than $\mathrm{X}^{\alpha}$ for the lexicographical order.

We shall see in Section 2.3.1 that the lexicographical ordering is particularly useful for solving zero-dimensional systems.

### 2.1.3. Graded reverse-lexicographical ordering

The next monomial orders that we present are degree orders: the monomials are ordered first by degree, then according to some other monomial order. We shall see that these orders are particularly suitable for working with homogeneous polynomials.

Definition 2.5. The graded reverse-lexicographical (or degree reverse lexicographical, or in short DRL) ordering GRevLex is defined as:

$$
\alpha<_{\text {grevlex }} \beta \Longleftrightarrow\left\{\begin{array}{l}
\operatorname{deg}\left(\mathbf{X}^{\alpha}\right)<\operatorname{deg}\left(\mathbf{X}^{\beta}\right) \\
\text { or } \\
\operatorname{deg}\left(\mathbf{X}^{\alpha}\right)=\operatorname{deg}\left(\mathbf{X}^{\beta}\right) \text { and } \alpha_{i}>\beta_{i}, \alpha_{i+1}=\beta_{i+1}, \ldots, \alpha_{n}=\beta_{n}
\end{array}\right.
$$

This definition extends straightforwardly to a weighted setting:
Definition 2.6. Let $W \in \mathbb{N}^{n}$ be a system of weights.The weighted graded reverse-lexicographical ordering $W$-GRevLex is defined as:

$$
\alpha<_{W \text {-grevlex }} \beta \Longleftrightarrow\left\{\begin{array}{l}
\operatorname{deg}_{W}\left(\mathrm{X}^{\alpha}\right)<\operatorname{deg}_{W}\left(\mathrm{X}^{\beta}\right) \\
\text { or } \\
\operatorname{deg}_{W}\left(\mathrm{X}^{\alpha}\right)=\operatorname{deg}_{W}\left(\mathrm{X}^{\beta}\right) \text { and } \alpha_{i}>\beta_{i}, \alpha_{i+1}=\beta_{i+1}, \ldots, \alpha_{n}=\beta_{n}
\end{array}\right.
$$

When working with homogeneous polynomials, the leading term for these orders contains a lot of information on the remaining terms. The next elementary propositions shall be useful in Section 2.5.3, in order to refine the complexity analyses for Algorithm $\mathrm{F}_{5}$.

Proposition 2.7. Let $f \in \mathbb{K}[\mathrm{X}]$ be a homogeneous (resp. W-homogeneous) polynomial. If $X_{i}^{k} \mid \mathrm{LT}_{<_{\text {greverex }}}(f)\left(\right.$ resp. $X_{i}^{k} \mid \mathrm{LT}_{<W \text {-grevex }}(f)$ ), then $f \in\left\langle X_{i}^{k}, X_{i+1}, \ldots, X_{n}\right\rangle$. In particular, if $\mathrm{LT}(f)$ is divisible by $X_{n}$, then $f$ is divisible by $X_{n}$.

Proof. It suffices to prove it in the weighted case, the total degree case is obtained by setting all weights to 1.
First assume that $\operatorname{LT}(f)=X_{1}^{\alpha_{1}} \cdots X_{i}^{\alpha_{i}}$, or in other words that $i$ is the largest index of a variable dividing $\operatorname{LT}(f)$. Let $\mu$ be a monomial with the same $W$-degree, and such that $\mu \notin$ $\left\langle X_{i}^{\alpha_{i}}, X_{i+1}, \ldots, X_{n}\right\rangle$, this means that

$$
\mu=X_{1}^{\beta_{1}} \cdots X_{i}^{\beta_{i}} \text { with } \beta_{i}<\alpha_{i},
$$

so $\mu<_{\text {lex }} \operatorname{LT}(f)$, and thus $\mu>_{W \text {-grevlex }} \operatorname{LT}(f)$, and it may not appear in the support of $f$.
If $i$ is not the largest index of a variable dividing $\operatorname{LT}(f)$, then $\operatorname{LT}(f)$ is divisible by $X_{j}$ for some $j>i$, and from the above

$$
f \in\left\langle X_{j}, \ldots, X_{n}\right\rangle \subset\left\langle X_{i}^{\alpha_{i}}, X_{i+1}, \ldots, X_{j}, \ldots, X_{n}\right\rangle .
$$

For $m \in\{1, \ldots, n\}$, let $\theta_{m}$ be the morphism evaluating $X_{m+1}, \ldots, X_{n}$ to 0 :

$$
\begin{array}{ccccl}
\theta_{m}: \mathbb{K}[\mathbf{X}] & \longrightarrow & \mathbb{K}[\mathbf{X}] & \\
& X_{i} & \longmapsto & X_{i} & (i \leq m) \\
& X_{i} & \longmapsto & 0 & (i>m)
\end{array}
$$

Corollary 2.8. Let $f \in \mathbb{K}[\mathrm{X}]$ be a homogeneous (resp. W-homogeneous) polynomial, and let $<$ be the GREVLEX (resp. $W$-GREVLEX) order. If $m \in\{1, \ldots, n\}$ is such that $\theta_{m}(f) \neq 0$, then

$$
\mathrm{LT}_{<}(f)=\mathrm{LT} T_{<}\left(\theta_{m}(f)\right)
$$

Proof. The evaluation $\theta_{m}(f)$ is nonzero if and only if $f \notin\left\langle X_{m+1}, \ldots, X_{n}\right\rangle$. So if $\theta_{m}(f) \neq 0$, $\mathrm{LT}(f)$ is not divisible by $X_{m+1}, \ldots, X_{n}$, and

$$
\operatorname{LT}\left(\theta_{m}(f)\right)=\theta_{m}(\operatorname{LT}(f))=\operatorname{LT}(f)
$$

Proposition 2.9 ([BFS14, Prop. 7]). Let $\left(f_{1}, \ldots, f_{m}\right)$ be a homogeneous (resp. W-homogeneous) polynomial system in Noether position with respect to the variables $X_{m+1}, \ldots, X_{n}$, let < the GREVLEX (resp. $W$-GREvLEx) order, and $I=\left\langle f_{1}, \ldots, f_{m}\right\rangle$. Then:

$$
\mathrm{LT}_{<}(I)=\mathrm{LT}_{<}\left(\theta_{m}(I)\right) \cdot\left\langle X_{m+1}, \ldots, X_{n}\right\rangle
$$

Proof. It suffices to prove it in the weighted case, the total degree case is obtained by setting all weights to 1.

From Corollary 2.8, for the W-GRevLex ordering, $\mathrm{LT}\left(\theta_{m}(I)\right) \subset \mathrm{LT}(I)$, this implies the reverse inclusion

$$
\operatorname{LT}(I) \supset \operatorname{LT}\left(\theta_{m}(I)\right) \cdot\left\langle X_{m+1}, \ldots, X_{n}\right\rangle
$$

Conversely, let $f \in I$ be a $W$-homogeneous polynomial, write $\operatorname{LT}(f)=X_{1}^{\alpha_{1}} \cdots X_{n}^{\alpha_{n}}$. We want to prove that there exists $g \in I$ such that $\operatorname{LT}(g)=X_{1}^{\alpha_{1}} \cdots X_{m}^{\alpha_{m}}$. Let $i$ be the largest index of a variable dividing $\operatorname{LT}(f)$ (that is the largest integer such that $\alpha_{i} \neq 0$ ). If $i \leq m$, there is nothing to prove.

Otherwise, from Proposition 2.7, this implies that $f \in\left\langle X_{i}^{\alpha_{i}}, X_{i+1}, \ldots, X_{i}\right\rangle$, which we may write

$$
f=X_{i}^{\alpha_{i}} g_{i}+\sum_{j=i+1}^{n} X_{j} g_{j}
$$

with $g_{j}(j \leq i) W$-homogeneous and $0 \neq g_{i} \in \mathbb{K}\left[X_{1}, \ldots, X_{i}\right]$ (all monomials divisible by $X_{j}$, for $j>i$, are pushed in the term $X_{j} g_{j}$ ). This implies that

$$
\begin{equation*}
X_{i}^{\alpha_{i}} g_{i} \in I+\left\langle X_{i+1}, \ldots, X_{n}\right\rangle \tag{2.1}
\end{equation*}
$$

By Proposition 1.44 , since $\left(f_{1}, \ldots, f_{m}\right)$ is in Noether position, $\left(f_{1}, \ldots, f_{m}, X_{m+1}, \ldots, X_{n}\right)$ is a regular sequence, and by Corollary 1.39, $\left(f_{1}, \ldots, f_{m}, X_{n}, \ldots, X_{m+1}\right)$ is also regular. Equation (2.1) means that $X_{i}^{\alpha_{i}}$ is a zero divisor in $\mathbb{K}[\mathbf{X}] /\left\langle f_{1}, \ldots, f_{m}, X_{n}, \ldots, X_{i+1}\right\rangle$, so by regularity,

$$
g_{i} \in\left\langle f_{1}, \ldots, f_{m}, X_{n}, \ldots, X_{i+1}\right\rangle=I+\left\langle X_{i+1}, \ldots, X_{n}\right\rangle
$$

Since $g_{i} \in \mathbb{K}\left[X_{1}, \ldots, X_{i}\right]$, we deduce that $g_{i} \in I$, and its leading term is $X_{1}^{\alpha_{1}} \cdots X_{i-1}^{\alpha_{i-1}}$. Repeating the process until $i \leq m$, we obtain the wanted polynomial.

### 2.1.4. Elimination orders

The last order that we present is an elimination ordering.
Definition 2.10. Let $k$ be a integer such that $1 \leq k \leq n$, the $k$ 'th elimination ordering $\operatorname{Eum}_{k}$ is defined as:

$$
\alpha<_{\text {Eим }_{k}} \beta \Longleftrightarrow\left\{\begin{array}{l}
\left(\alpha_{1}, \ldots, \alpha_{k}\right)<_{\text {grevlex }}\left(\beta_{1}, \ldots, \beta_{k}\right) \\
\text { or } \\
\left(\alpha_{1}, \ldots, \alpha_{k}\right)=\left(\beta_{1}, \ldots, \beta_{k}\right) \text { and }\left(\alpha_{k+1}, \ldots, \alpha_{n}\right)<_{\text {grevlex }}\left(\beta_{k+1}, \ldots, \beta_{n}\right)
\end{array}\right.
$$

The name comes from the fact that this ordering can be used to eliminate variables.
Definition 2.11. Let $<_{\text {mon }}$ be a monomial ordering, we say that $<_{\text {mon }}$ eliminates the first $k$ variables if and only if for any $\left(\alpha_{1}, \ldots, \alpha_{n}\right) \in \mathbb{N}^{n}$ and $\left(\beta_{k+1}, \ldots, \beta_{n}\right) \in \mathbb{N}^{n-k}$,

$$
\left(0, \ldots, 0, \beta_{k+1}, \ldots, \beta_{n}\right)<_{\operatorname{mon}}\left(\alpha_{1}, \ldots, \alpha_{n}\right)
$$

Equivalently, for $f \in \mathbb{K}[\mathrm{X}]$, if $\mathrm{LT}_{<_{\text {mon }}}(f) \in \mathbb{K}\left[X_{k+1}, \ldots, X_{n}\right]$, then $f \in \mathbb{K}\left[X_{k+1}, \ldots, X_{n}\right]$.
From Proposition 2.4, the lexicographical order eliminates the first $k$ variables for all $k$. On the other hand, the graded reverse lexicographical ordering does not eliminate any group of variables in general (there may be monomials with lower degree in the trailing terms). The elimination orderings are in some sense an intermediate ground between the lexicographical and the graded reverse-lexicographical orderings:

Proposition 2.12. Let $k \in\{1, \ldots, n\}$, the $k$ 'th elimination ordering eliminates the first $k$ variables.

Proof. Immediate consequence of the definition.

### 2.2. Gröbner bases: definition

We can now define Gröbner bases.
Definition 2.13 (Gröbner basis). Let $I$ be an ideal of $\mathbb{K}[\mathrm{X}]$. A Gröbner basis of $I$ is a subset $G=\left\{g_{1}, \ldots, g_{k}\right\} \subset I$ such that

$$
\left\langle\operatorname{LT}\left(g_{1}\right), \ldots, \operatorname{LT}\left(g_{k}\right)\right\rangle=\langle\operatorname{LT}(I)\rangle .
$$

This definition can be rephrased as follows: a Gröbner basis of $I$ is a finite subset $G \subset I$ and such that

$$
\forall f \in I, f \neq 0 \Longrightarrow \exists g \in G, \operatorname{LT}(g) \mid \operatorname{LT}(f) .
$$

The fundamental property of Gröbner bases is that they define a normal form:

Proposition/Definition 2.14 (Normal Form application, [BW93, Th. 5.35]). Let $I$ be an ideal of $\mathbb{K}[\mathrm{X}]$ and let $G$ be a Gröbner basis of $I$. Let $f \in \mathbb{K}[\mathrm{X}]$ be a polynomial. There exists a unique polynomial $N F(f) \in \mathbb{K}[X]$ such that:

- $f-\mathrm{NF}(f) \in\langle G\rangle$ (we say that $f$ reduces to $\mathrm{NF}(f)$ modulo $G$, written $f \xrightarrow{G} \mathrm{NF}(f)$ );
- For all $\mu$ in the support of $\mathrm{NF}(f), \mu$ is not divisible by any $\operatorname{LT}(g), g \in G$.

The polynomial $\operatorname{NF}(f)$ is called the normal form of $f$ modulo $G$. If $f=\operatorname{NF}(f)$, we say that $f$ is in normal form. Otherwise, $f$ is reducible modulo $G$.

If $\operatorname{LT}(f)$ is reducible modulo $G$, we say that $f$ is head-reducible modulo $G$. If $f$ is reducible modulo $G$ but $\mathrm{LT}(f)$ is irreducible, we say that $f$ is tail-reducible.

This notion allows to prove that a Gröbner basis of $I$ is a set of generators of $I$.
Proposition 2.15. Let $G$ be a Gröbner basis of $I$, then $\langle G\rangle=I$.
Proof. By definition, $\langle G\rangle \subset I$. Let $f \in I, f-\mathrm{NF}(f) \in\langle G\rangle \subset I$, so $\mathrm{NF}(f) \in I$. Assume that $\operatorname{NF}(f) \neq 0$, then by definition of a Gröbner basis, $\operatorname{LT}(\operatorname{NF}(f))$ is divisible by $\operatorname{LT}(g)$ for some $g \in G$, thus contradicting the definition of $\operatorname{NF}(f)$. Hence $\operatorname{NF}(f)=0$ and $f \in\langle G\rangle$.

The normal form application has a lot more properties which we list now.
Proposition 2.16 ([CLOo7]). Let $f, g \in \mathbb{K}[\mathrm{X}]$ and let NF be a normal form application for a Gröbner basis $G$ of some ideal I. Then:

1. $\mathrm{NF}(f)$ is in normal form with respect to $G$;
2. $\operatorname{NF}(f)=0$ if and only if $f \in I$;
3. $\mathrm{NF}(f)=\mathrm{NF}(g)$ if and only if $f-g \in I$;
4. $\mathrm{NF}(f+g)=\mathrm{NF}(f)+\mathrm{NF}(g)$;
5. $\mathrm{NF}(f g)=\operatorname{NF}(\mathrm{NF}(f) \mathrm{NF}(g))$;

In particular, NF defines a surjective ring morphism $\mathbb{K}[\mathrm{X}] \longrightarrow \mathbb{K}[\mathrm{X}] / I$.
Proof. Item 1 is a consequence of the unicity of the normal form. Item 2 is a consequence of the proof of Proposition 2.15.

For item 3, write

$$
f-g=(f-\mathrm{NF}(f))+(\mathrm{NF}(f)-\mathrm{NF}(g))+(\mathrm{NF}(g)-g),
$$

so $f-g \in I$ if and only if $\mathrm{NF}(f)-\mathrm{NF}(g) \in I$. In particular, if $\mathrm{NF}(f)-\mathrm{NF}(g)=0$, then $f-g=0$.
Conversely, the support of $\operatorname{NF}(f)-N F(g)$ is a subset of the union of the supports of $N F(f)$ and $\mathrm{NF}(g)$, so all monomials of $\mathrm{NF}(f)-\mathrm{NF}(g)$ are irreducible modulo $G$. So $\mathrm{NF}(f)-\mathrm{NF}(g)$ is in normal form, and so if $\mathrm{NF}(f)-\mathrm{NF}(g) \in I$, then it is zero.

Items 4 and 5 are proved similarly: since

$$
\mathrm{NF}(f+g)-\mathrm{NF}(f)-\mathrm{NF}(g)=(\mathrm{NF}(f+g)-f-g)+(f-\mathrm{NF}(f))+(g-\mathrm{NF}(g)) \in I,
$$

the same argument involving the support shows that

$$
\mathrm{NF}(f+g)-\mathrm{NF}(f)-\mathrm{NF}(g)=0
$$

And since

$$
f g-\mathrm{NF}(f) \mathrm{NF}(g)=f(g-\mathrm{NF}(g))-(f-\mathrm{NF}(f)) \mathrm{NF}(g) \in I,
$$

taking the normal form of both terms of the left hand side yields that

$$
\mathrm{NF}(f g)-\mathrm{NF}(\mathrm{NF}(f) \mathrm{NF}(g))=0
$$

Item 2 is the main property of Gröbner bases: the normal form application that they define give a practical way of testing whether any polynomial belong to the ideal.

The last two propositions of this section are existence statements for Gröbner bases. We shall describe algorithms computing Gröbner bases later, which will give a constructive proof of these results.

Proposition 2.17 ([CLOo7, Ch. 2, sec. 5, cor. 6]). Let I be an ideal of $\mathbb{K}[\mathrm{X}]$, then I admits a Gröbner basis for any given monomial order.

Gröbner bases are not unique in general, since we may add polynomials from the ideal to any Gröbner basis and still obtain a Gröbner basis. However, with additional properties, unicity can be ensured.

Definition 2.18. Let $I$ be an ideal of $\mathbb{K}[X]$, and let $G$ be a Gröbner basis of $I$. We say that $G$ is reduced if the following two conditions are satisfied:

1. all leading coefficients of polynomials of $G$ are 1 ;
2. all polynomials in $G$ are in normal form with respect to the others.

Proposition 2.19 ([CLOo7, Ch. 2, sec. 7, prop. 6]). Let I be an ideal of $\mathbb{K}[\mathrm{X}]$, then I admits a unique reduced Gröbner basis for any given monomial order.

### 2.3. Applications

What information one can obtain from a Gröbner basis depends on the monomial order that is used. The most useful orders in that regard are the lexicographical order and the elimination orders.

### 2.3.1. Zero-dimensional systems

Recall that zero-dimensional systems are systems with a finite number of solutions. A first property is that we can check this property using any Gröbner basis.

Proposition 2.20. Let $I$ be an ideal of $\mathbb{K}[\mathrm{X}]$, and let $G$ be a Gröbner basis of I. The ideal I is zero-dimensional if and only if for all $X \in\{\mathbf{X}\}$, there exists a polynomial in $G$ with leading term $X^{\alpha}$ for some $\alpha$.

Proof. Recall that $I$ is zero-dimensional if and only if $\mathbb{K}[\mathrm{X}] / I$ is a $\mathbb{K}$-vector space with finite dimension. Let

$$
M_{I}=\{\mu \text { monomial of } \mathbb{K}[\mathbf{X}] \mid \forall g \in G, \operatorname{LT}(g) \nmid \mu\},
$$

elements of $M_{I}$ are all linearly independent modulo $I$. Hence $M_{I}$ is finite.
Let $X \in\{\mathbf{X}\}$, there exists $\beta \in \mathbb{N}$ such that $X^{\beta} \notin M_{I}$, and so there exists $\alpha \in \mathbb{N}, \alpha \leq \beta$ such that $X^{\alpha}$ is the leading term of an element of $G$.

The lexicographical order allows one to solve generic zero-dimensional systems.
Definition 2.21 (Shape position). Let $I$ be an ideal in $\mathbb{K}[\mathbf{X}]$, and let $G$ be the reduced lexicographical Gröbner basis of $I$. We say that $I$ is in shape position if $G$ has the following shape:

$$
\begin{cases}g_{1}\left(X_{1}, \ldots, X_{n}\right) & =X_{1}+r_{1}\left(X_{n}\right) \\ g_{2}\left(X_{2}, \ldots, X_{n}\right) & =X_{2}+r_{2}\left(X_{n}\right) \\ & \vdots \\ g_{n-1}\left(X_{n-1}, X_{n}\right) & =X_{n-1}+r_{n-1}\left(X_{n}\right) \\ g_{n}\left(X_{n}\right) & =X_{n}^{\operatorname{deg}(I)}+r_{n}\left(X_{n}\right)\end{cases}
$$

In particular, if $I$ is in shape position, it is zero-dimensional, and its $\operatorname{deg}(I)$ roots can be recovered by solving the univariate polynomial $g_{n}$.

This property is generic in the following sense:
Lemma 2.22 (Shape lemma, [GM89]). Up to a generic change of coordinates, any zero-dimensional ideal is in shape position.

This makes lexicographical Gröbner bases a very powerful tool for solving zero-dimensional systems.

### 2.3.2. Positive-dimensional systems and eliminations

Gröbner bases can also be used to work with positive-dimensional systems. In particular, elimination bases allow to "compute" projections: the projection of an algebraic variety onto an affine subspace is not in general an algebraic variety, but its Zariski closure is, and one can compute equations defining this algebraic set by polynomial elimination.

Proposition 2.23 ([CLOo7, Chap. 3, Sec. 1, Th. 2]). Let I be an ideal of $\mathbb{K}[\mathrm{X}]$, let $k \in\{1, \ldots, n\}$, and let $G$ be a Gröbner basis of I for an order eliminating the first $k$ variables (for example ELIM ${ }_{k}$ or LEX).

Then the elimination basis $G_{k}:=G \cap \mathbb{K}\left[X_{k+1}, \ldots, X_{n}\right]$ is a Gröbner basis of $I \cap \mathbb{K}\left[X_{k+1}, \ldots, X_{n}\right]$. Geometrically, $G_{k}$ is a system defining the Zariski closure of the projection of $V(I)$ onto the coordinates $\left(X_{k+1}, \ldots, X_{n}\right)$.

Remark 2.24. In [CLOo7, Chap. 3, Sec. 1, Th. 2], the theorem is written for the lexicographical order, but the proof only uses the fact that this order eliminating the first $k$ variable.
Remark 2.25. The shape position is an example of successive eliminations: with the notations of Definition 2.21, for any $i \in 1, \ldots, n$, the polynomials $\left\{g_{i}, \ldots, g_{n}\right\}$ generate the ideal $I \cap$ $\mathbb{K}\left[X_{i}, \ldots, X_{n}\right]$. In particular, the last polynomial $g_{n}$ generates $I \cap \mathbb{K}\left[X_{n}\right]$.

Geometrically, a zero-dimensional ideal $I$ has finitely-many solutions. Their projection on the $X_{n}$ line is a finite set, and $I$ being in shape position means that each of these roots $x_{n}$ is the projection of only 1 point $\left(x_{1}, \ldots, x_{n}\right) \in V(I)$ onto the $X_{n}$ line, and this point is given by

$$
\left(-r_{1}\left(x_{n}\right), \ldots,-r_{n}\left(x_{n}\right), x_{n}\right)
$$

Another example of how elimination techniques can be useful is saturation:
Proposition 2.26 ([CLOo7, Sec. 4.4, Ex. 9]). Let I be an ideal of $\mathbb{K}[\mathrm{X}]$, let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a system of generators of $I$, and let $g \in \mathbb{K}[\mathrm{X}]$. Consider the saturated system

$$
F_{\text {sat }}:=\left(f_{1}, \ldots, f_{m}, U \cdot g-1\right) \subset \mathbb{K}\left[U, X_{1}, \ldots, X_{n}\right] .
$$

Let $G_{\text {sat }}$ be a Gröbner basis of $\left\langle F_{\text {sat }}\right\rangle$ eliminating the first variable $U$. Then $G_{\text {sat }}$ is a Gröbner basis of the saturated ideal $\left(I: g^{\infty}\right)$ defined as

$$
\left(I: g^{\infty}\right):=\left\{f \in \mathbb{K}[\mathbf{X}] \mid \exists n \in \mathbb{N}, g^{n} f \in I\right\} .
$$

Geometrically, $G_{\text {sat }}$ is a system defining the Zariski closure of the complementary of $V(\langle g\rangle)$ in $V(I)$.

### 2.4. Algorithms

Many algorithms have been developed to compute Gröbner bases, we will present three examples here.

### 2.4.1. A pairs algorithm: Buchberger's algorithm

The first algorithm for computing Gröbner bases was Buchberger's algorithm [Buc76]. Its building block is pairwise reductions.

Definition 2.27 (S-polynomial). Let $f, g \in \mathbb{K}[\mathbf{X}]$, the $S$-polynomial of $f$ and $g$ is defined as

$$
S-\operatorname{pol}(f, g):=\operatorname{lcm}(\mathrm{LM}(f), \operatorname{LM}(g))\left(\frac{f}{\operatorname{LM}(g)}-\frac{g}{\operatorname{LM}(f)}\right)
$$

In particular, if $f$ is not in normal form with respect to $g, S-\operatorname{pol}(f, g)$ is a reduction of $f$ modulo $g$ :

$$
S-\operatorname{pol}(f, g)=f-\frac{\operatorname{lcm}(\mathrm{LM}(f), \mathrm{LM}(g))}{\mathrm{LM}(f)} \cdot g
$$

This allows for computing polynomial reductions algorithmically.

```
Algorithm 2.1 Reduction
Input: \(f \in \mathbb{K}[\mathbf{X}], G \subset \mathbb{K}[\mathbf{X}]\)
Output: \(h \in \mathbb{K}[\mathbf{X}]\) such that \(f \xrightarrow{G} h\) and \(h\) is in normal form with respect to \(G\)
    \(h \leftarrow f\)
    res \(\leftarrow 0\)
    while \(h \neq 0\) do
        while \(\exists g \in G\) such that \(\operatorname{LT}(g) \mid \operatorname{LT}(h)\) do
            \(h \leftarrow S-\operatorname{pol}(h, g)\)
        end while
        res \(\leftarrow\) res \(+\mathrm{LT}(h)\)
        \(h \leftarrow h-\operatorname{LT}(h)\)
    end while
    return res
```

Remark 2.28. If $G$ is a Gröbner basis and $f$ is a polynomial, Reduction $(f, G)$ is the normal form of $f$ with respect to $G$.
Remark 2.29. If $G$ is a Gröbner basis of some ideal $I$, let

$$
G^{\prime}:=\{\operatorname{Reduction}(g, G \backslash g) \mid g \in G\} \backslash\{0\} .
$$

Then $G^{\prime}$ is a reduced Gröbner basis of $I$ for the same monomial order.
Buchberger's algorithm uses S-polynomials and Algorithm Reduction to compute a Gröbner basis of an ideal: given a set of generators of the ideal, it creates all $S$-polynomials of pairs of these generators, then reduces them modulo the initial set of generators. The result is a new, larger set of generators, and it iterates these steps until all $S$-polynomials reduce to zero.

```
Algorithm 2.2 Buchberger ([Buc76])
Input: \(F \subset \mathbb{K}[\mathrm{X}]\)
Output: \(G \subset \mathbb{K}[\mathrm{X}]\) such that \(G\) is a Gröbner basis of \(\langle F\rangle\)
    \(G \leftarrow F\)
    \(Q \leftarrow F \times F\)
    while \(Q \neq \varnothing\) do
        Pick \((f, g)\) in \(Q\), remove \((f, g)\) from \(Q\)
        \(h \leftarrow S-\operatorname{pol}(f, g)\)
        \(h \leftarrow\) Reduction \((h, G)\)
        if \(h \neq 0\) then
            \(Q \leftarrow Q \cup G \times\{h\}\)
            \(G \leftarrow G \cup\{h\}\)
        end if
    end while
    return \(G\)
```

Theorem 2.30 (Correctness of Buchberger's algorithm, [Buc76],[CLOo7, Ch. 2, sec. 7, th. 2]). Buchberger's algorithm terminates and returns a Gröbner basis of I.

Remark 2.31. Buchberger's algorithm being correct gives another proof that Gröbner bases exist for all ideals.

This version of Buchberger's algorithm is only of theoretical interest. In practice, many pairs will eventually reduce to the zero polynomial, and thus are effectively useless for the rest of the computations. To mitigate this effect, two algorithmic criteria were given in [Buc79], filtering out some pairs leading to reductions to zero. The algorithm can also be improved by defining strategies for the choice of the next pair (for example the sugar strategy, see [Gio+91]) Buchberger's algorithm is currently implemented, for example in Macaulay2 [Macaulay2], Singular [Singular] or for computations over arbitrary rings in Magma[Magma] or Maple[Maple].

Buchberger's algorithm belongs to a wide family of Gröbner basis algorithms, called pairs algorithms: they construct a set of pairs, reduce them, and build new pairs with the result.

Other examples include the HDriven algorithm [Tra96] (Buchberger with an additional criterion based on information on the Hilbert series of the ideal, currently implemented in most systems implementing Buchberger), the $\mathrm{F}_{4}$ algorithm [Fau99] (selecting several pairs at each step and using linear algebra instead of $S$-polynomials to compute the reductions, currently implemented in the library FGb [Fau1o], in Magma and in Maple), and the $\mathrm{F}_{5}$ algorithm [Fauo2] ( $F_{4}$ with the additional $F_{5}$ criterion, see Section 2.4.2).

### 2.4.2. A matrix algorithm: Matrix- $\mathrm{F}_{5}$

The complexity of the previous algorithms is difficult to estimate, because it mainly depends on the length of the queue at any given step. However, the complexity of these algorithms can be bounded by studying matrix variants, replacing $S$-polynomials and critical pairs with the construction and reduction of Macaulay matrices.

For example, the matrix variant of Buchberger's algorithm is Lazard's algorithm [Laz83]: at each degree $d$, it builds the Macaulay matrix of degree $d$ of the system, reduces it to echelon form, and extract a new set of generators from the new matrix.

In this section, we describe a matrix variant of $\mathrm{F}_{5}$ algorithm, called Matrix- $\mathrm{F}_{5}$ ([BFS14]).
This algorithm works by keeping track, for each computed polynomial, of "where it came from". It uses this information to exclude some rows from the constructed matrix: the matrix built is a submatrix of the Macaulay matrix, whose rows span the same vector space as the rows of the Macaulay matrix.

Definition 2.32 (Signature). Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a family of homogeneous polynomials in $\mathbb{K}[\mathrm{X}]$, and let $I=\langle F\rangle$. Consider a polynomial $g \in I$ such that

$$
g=g_{i} f_{i}+\sum_{j=1}^{i-1} g_{j} f_{j}
$$

with $g_{i} \neq 0$ and $i$ minimal for this property. The signature of $g$ is the pair $\left(i, \operatorname{LT}\left(g_{i}\right)\right)$.
Signatures are ordered using the lexicographical ordering on the two components:

$$
(i, \mu)<(j, v) \Longleftrightarrow\left\{\begin{array}{l}
i<j \\
\text { or } \\
i=j \text { and } \mu<v
\end{array}\right.
$$

Remark 2.33. Decomposing the $g_{j}$ 's into sums of monomials, the definition means that a polynomial with signature $(i, \mu)$ can be written as a linear combination of $\mu f_{i}$ and polynomials with smaller signatures.

Lemma 2.34 ( $\mathrm{F}_{5}$ criterion, [Fauo2]). Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a family of homogeneous polynomials in $\mathbb{K}[\mathrm{X}]$, and let $I=\langle F\rangle$. Let $i \in\{1, \ldots, m\}$ and $\mu$ be a monomial of $\mathbb{K}[\mathrm{X}]$. If $\mu$ is the leading term of a polynomial in $\left\langle f_{1}, \ldots, f_{i-1}\right\rangle$, then the polynomial $\mu f_{i}$ is a linear combination of polynomials with signature less than $(i, \mu)$. In other words, no polynomial in I has signature $(i, \mu)$.

Proof. Let $\mu=\mathrm{LT}(g)$, and let $(j, \bullet)$ be the signature of $g$. By hypothesis, $j<i$. Write

$$
g=g_{j} f_{j}+\sum_{k=1}^{j-1} g_{k} f_{k}
$$

and let $r=g-\mu$. All monomials in $r$ are smaller than $\mu$. Write

$$
\begin{aligned}
\mu f_{i} & =g f_{i}-r_{2} f_{i} \\
& =-r f_{i}+\sum_{k=1}^{j} g_{k} f_{i} f_{k}
\end{aligned}
$$

Hence $\mu f_{i}$ has signature at most $(i, \mathrm{LT}(r))$ with $\mathrm{LT}(r)<\mu$.

Unlike pairs-based algorithms which stop when the pair queue is empty, there is no obvious stopping criterion for matrix-based algorithms. This means that the algorithm takes an additional parameter along with the system, a degree $d_{\max }$ at which to stop. Formally, what the algorithm computes is a Gröbner basis truncated to $d_{\max }$ (or a $d_{\max }$-Gröbner basis), that is a set of polynomials containing the polynomials of the reduced Gröbner basis of $I$ of degree at most $d_{\text {max }}$.
In Algorithm 2.3, EchelonForm denotes a routine performing Gaussian elimination on a matrix, with restrictions:

- rows are not exchanged;
- rows are only reduced by rows above them.

These conditions will be used to ensure that any polynomial is only reduced by polynomials with smaller signatures.

The notation $\operatorname{row}(M, i)$ is a shorthand for the polynomial associated with the $i$ 'th row of the matrix $M$. When talking about matrices, the notation $A \cup B$ means that, both matrices having the same number of columns, we append matrix $B$ to the bottom of matrix $A$.
The algorithm proceeds by looping over the degree $d$, and then over the polynomials $f_{i}$, and then over the monomials $\mu$ with degree $d-d_{i}$, as if to build the Macaulay matrix with rows $\mu f_{i}$. For each row $\mu f_{i}$, it tests if it satisfies the $\mathrm{F}_{5}$-criterion (line 11), and if not, it computes its signature and adds the row to the matrix.

Theorem 2.35 (Correctness of algorithm Matrix- $\mathrm{F}_{5}$, [Fauo2; BFS14]). Algorithm Matrix- $F_{5}$ terminates and computes a $d_{\max ^{-G}}$-Gröbner basis of $F$. Furthermore, if $F$ is a regular sequence, then no rows are reduced to zero in the echelon forms.

Proof. For all $i \in\{1, \ldots, m\}$, let $F_{i}=\left\{f_{1}, \ldots, f_{i}\right\}$ and $I_{i}=\left\langle F_{i}\right\rangle$.
We will prove by recurrence over $(d, i)$ that the rows of $M_{d, i} \operatorname{span}\left(I_{i}\right)_{d}$. Note that for any ( $i, d$ ) such that it is true, after Gaussian elimination, the rows of $\tilde{M}_{d, i}$ still span $\left(I_{i}\right)_{d}$, and the leading terms of these rows are the leading terms of $\left(I_{i}\right)_{d}$.
The base cases $d=\min \left(d_{j}\right)$ and $i=0$ for any $d$ are clear.
Assume that $d>\min \left(d_{j}\right)$ and $i>0$. The matrix $M_{d, i}$ is a submatrix of the Macaulay matrix $\mathrm{Mac}_{d}$ of $F_{i}$. More precisely, it is obtained by removing from $\mathrm{Mac}_{d}$ those rows failing to match the $\mathrm{F}_{5}$ criterion.
By induction hypothesis, the rows with signature $(j, \bullet), j<i \operatorname{span}\left(I_{i-1}\right)_{d}$, so we need only check those rows with signature $(i, \bullet)$. If $M_{d, i}$ does not span $\left(I_{i}\right)_{d}$, then there exists a polynomial

$$
f=\sum_{j=1}^{i} g_{j} f_{j}
$$

which does not lie in $\operatorname{Vect}\left(M_{d, i}\right)$. Let $\mu=\operatorname{LT}\left(g_{j}\right)$, and assume that $f$ is such that $\mu$ is minimal. Write $g_{i}=\mu+r_{i}$, all monomials in the support of $r_{i}$ are smaller than $\mu$, and by minimality of $\mu$

$$
r:=r_{i} f_{i}+\sum_{j=1}^{i-1} g_{j} f_{j} \in \operatorname{Vect}\left(M_{d, i}\right)
$$

```
Algorithm 2.3 Matrix- \(\mathrm{F}_{5}\) ([BFS14])
Output: \(F \subset \mathbb{K}[\mathrm{X}]\) homogeneous system, \(d_{\max } \in \mathbb{N}\) an implicit monomial ordering
Input: \(G \subset \mathbb{K}[\mathrm{X}]\), truncated Gröbner basis of \(\langle F\rangle\) at degree \(d_{\text {max }}\)
    \(G \leftarrow \varnothing\)
    for \(d=0\) to \(d_{\text {max }}\) do
        \(N_{d} \leftarrow\) number of monomials of degree \(d\)
        \(M_{d, 0}, \tilde{M}_{d, 0} \leftarrow\) matrix with o rows and \(N_{d}\) columns
        for \(i=1\) to \(m\) do
            \(M_{d, i} \leftarrow M_{d, i-1}\)
            if \(d=d_{i}\) then
                \(M_{d, i} \leftarrow M_{d, i} \cup\) row \(f_{i}\) with signature ( \((, 1\) )
                else if \(d>d_{i}\) then
                    for all \(\mu\) monomial of degree \(d-d_{i}\) do
                    if \(\mu\) is not leading term of a row in \(\tilde{M}_{d-d_{i}, i-1}\) then
                                    \(j \leftarrow\) largest \(j\) such that \(X_{j}\) divides \(\mu\)
                                    \(\mu_{\tilde{f}}^{\prime} \leftarrow \mu / X_{j}\)
                                    \(\tilde{f} \leftarrow\) row of \(\tilde{M}_{d-1, i}\) with signature ( \(i, \mu^{\prime}\) )
                                    \(M_{d, i} \leftarrow M_{d, i} \cup\) row \(X_{j} \tilde{f}\) with signature \((i, \mu)\)
                    end if
                    end for
            end if
        end for
        \(\tilde{M}_{d, i} \leftarrow\) EchelonForm \(\left(M_{d, i}\right)\)
        \(R \leftarrow\) number of rows in \(M_{d, i}\)
        \(G \leftarrow G \cup\left\{\operatorname{row}\left(\tilde{M}_{d, i}, j\right) \mid j \in\{1, \ldots, R\}, \operatorname{row}\left(\tilde{M}_{d, i}, j\right) \neq \operatorname{row}\left(M_{d, i}, j\right)\right\}\)
    end for
    return \(G\)
```

Furthermore, since signatures are inserted in increasing order, $r$ lies in the vector space spanned by all rows inserted before considering signature $(i, \mu)$.

If the row with signature $(i, \mu)$ is inserted in the matrix, then $f$ is the sum of this line and $r$, and so it lies in $\operatorname{Vect}\left(M_{d, i}\right)$. So the signature $(i, \mu)$ was rejected by the $F_{5}$ criterion, which implies that the row $\mu f_{i}$ is a linear combination of rows with signature less than $(i, \mu)$, and so again it lies in $\operatorname{Vect}\left(M_{d, i}\right)$. So we have reached a contradiction, and we conclude that

$$
\operatorname{Vect}\left(M_{d, i}\right)=\left(I_{i}\right)_{d}
$$

For the second part of the statement, assume that there is a row reduced to zero in a run of Matrix- $F_{5}$, and let $(i, \mu)$ be the signature of such a row. It means that there exist polynomials $g_{1}, \ldots, g_{i}$ such that

$$
\sum_{j=1}^{i} g_{j} f_{j}=0
$$

and so

$$
g_{i} f_{i} \in\left\langle f_{1}, \ldots, f_{i-1}\right\rangle
$$

If $F$ is regular, this implies that $g_{i} \in I_{i-1}$, and so $\mu=\operatorname{LT}\left(g_{i}\right) \in \operatorname{LT}\left(I_{i-1}\right)$. Thus this row is rejected by the $F_{5}$ criterion, and we reached a contradiction.

### 2.4.3. Matrix algorithms in the affine case

In the previous section, we described a matrix algorithm computing Gröbner bases for homogeneous systems. The idea was to build a full-rank submatrix of the Macaulay matrix of the system, and reduce it to echelon form.

Macaulay matrices can be defined in an affine setting, in the following way:
Definition 2.36. If $F=\left(f_{1}, \ldots, f_{m}\right) \in \mathbb{K}[\mathrm{X}]^{m}$ is any polynomial system (not necessarily homogeneous), the Macaulay matrix of $F$ is defined as the Macaulay matrix of the homogenized system $F^{h}$. Equivalently, it is the matrix of a basis of the set of products $\mu f_{i}$ of degree at most $d$.

In general, this set of products is not a basis of the set of polynomials in $\langle F\rangle$ with degree at most $d$.
Example 2.37. In $\mathbb{K}[X, Y]$, let $f_{1}=X Y+Y$ and $f_{2}=X$. At total degree 1 , the only polynomial in the Macaulay matrix is $f_{2}=X$. However, $f_{1}-Y f_{2}=Y$ is also a degree 1 polynomial in the ideal.

Matrix-based Gröbner basis algorithms can be adapted to the affine case in a way similar to the construction of the affine Macaulay matrix: compute $F^{h}$, run the algorithm on this homogeneous system, and dehomogenize the result. But for the same reason, Matrix- $F_{5}\left(F^{h}, d_{\max }\right)$ does not necessarily return a $d_{\max }$-Gröbner basis of $\langle F\rangle^{h}$.
Example 2.38. With the system of Example 2.37, the output of Matrix- $\mathrm{F}_{5}(F, 1)$ is $X$, which does not generate $\langle F\rangle=\langle F\rangle^{h}=\langle X, Y\rangle$.

This behavior is called a degree fall:

Definition 2.39. Let $F$ be a polynomial system and let $F^{h}$ be the homogenized system (with an homogenization variable $H$ ). We say that Matrix- $\mathrm{F}_{5}$ has a degree fall if it computes a row corresponding to a polynomial divisible by $H^{k}$ for some $k$.

Ignoring degree falls is not efficient, because it requires the algorithm to reach an unnecessarily high degree. Instead, degree falls can be handled in the following way:

1. Detect them in all lines computed at a given step;
2. Reinject them into the matrix at the relevant degree

Note that in order to detect whether some polynomial $f$ represents a degree fall, for the GRevLex order, it suffices to check whether $\operatorname{LT}(f)$ is divisible by $H$.
This technique mitigates the additional complexity induces by degree falls, and it is quite efficient in some cases. However, estimating the complexity of the algorithms now requires to be able to count the number of degree falls and their magnitude.
For this reason, the complexity of matrix Gröbner basis algorithms for affine systems is often evaluated under some generic assumption ensuring that there is no degree fall.

Definition 2.40. Let $F$ be a polynomial system. We say that $F$ is regular in the affine sense if the sequence $\left(h_{1}, \ldots, h_{m}\right)$, where for all $i \in\{1, \ldots, m\}, h_{i}$ is the highest degree component of $f_{i}$, is regular.

Proposition 2.41. Let $d_{1}, \ldots, d_{m} \in \mathbb{N}_{>0}$. Sequences regular in the affine sense are generic amongst sequence of polynomials of respective degree less than $d_{i}$.

Proof. There exists a non-empty Zariski-open subset $U$ of the space $A_{=}$of families of homogeneous polynomials of respective degree $d_{i}$ such that $U$ is contained in the set of regular sequences. The space of all families of arbitrary polynomials of respective degree $d_{i}$ is the product of $A_{=}$and the space $A_{<}$of all polynomials of respective degree less than $d_{i}$, and the non-empty Zariski-open set $U \times A_{<}$is contained in the set of sequences regular in the affine sense.

Proposition 2.42 ([Baro4]). Let $F$ be a polynomial system, assume that $F$ is regular in the affine sense. Let $d_{\max } \in \mathbb{N}$, let $G_{h}=$ Matrix- $\mathrm{F}_{5}\left(F^{h}, d_{\max }\right)$ and let $G=G_{h}^{a}$. Then Matrix- $F_{5}$ computes $G_{h}$ without any degree fall, and $G$ is a $d_{\text {max }}$-Gröbner basis of $\langle F\rangle$;
Furthermore, assume that we are using the GREVLEX monomial order, and let $d_{\mathrm{reg}}$ be the degree of regularity of the sequence $F_{h}$ of the highest degree components of $F$. Then if $d_{\max } \geq d_{\mathrm{reg}}$, then $G$ is a Gröbner basis of $\langle F\rangle$.

Proof. A degree fall is a reduction to zero of the highest degree components. Since by hypothesis these highest degree components form a regular sequence, the $F_{5}$ Criterion ensures that no such reduction happen.
If there is no degree fall, the rows of the $\mathrm{F}_{5}$ matrix at any degree $d$ represent a linear basis of the set of polynomials of degree $d$ in the homogenized ideal $\langle F\rangle^{h}$, so $G_{h}$ is a $d_{\text {max }}$-Gröbner basis of $\langle F\rangle^{h}$, and $G$ is a $d_{\text {max }}$-Gröbner basis of $\langle F\rangle$.

For the last part of the proposition, note that the GRevLex order is a degree order, and so the leading term of $f$ is the leading term of its highest degree component. Since there is no degree fall in the algorithm, all pivots used for computing the echelon form, and thus all leading terms of a polynomial in $G_{h}$, are monomials in $X_{1}, \ldots, X_{n}$ (they are not divisible by $H$ ). This shows that $\mathrm{LT}(\langle F\rangle)=\operatorname{LT}\left(\left\langle F_{h}\right\rangle\right)$, and since generators of $\operatorname{LT}\left(\left\langle F_{h}\right\rangle\right)$ can be computed at degree $d_{\text {reg }}$, the same follows for $\langle F\rangle$.

Remark 2.43. If $F$ is a polynomial system which forms a regular sequence in the affine sense, then the homogenized ideal of $\langle F\rangle$ is generated by the homogenization of $F$ :

$$
\langle F\rangle^{h}=\left\langle F^{h}\right\rangle
$$

This is false in general. For example, the system from Example 2.37 is equal to $\langle X, Y\rangle$, and so it is homogeneous. But $\left\langle f^{h}, g^{h}\right\rangle=\langle X, H Y\rangle$.

In general, the relationship between the homogenized ideal and the homogenized generators involves a saturation by $H$ :

$$
\langle F\rangle^{h}=\left\langle F^{h}\right\rangle: H^{\infty}
$$

### 2.4.4. Change of order: FGLM

In this section, we present algorithm FGLM ( [Fau+93]), which allows for change of ordering in the zero-dimensional case. Informally, it converts a Gröbner basis for some monomial order into a Gröbner basis for another order.

The reason such algorithms were developed is because the direct algorithms such as $F_{5}$ are generally much more efficient at computing bases for some orders than others. Typically, the GRevLex is usually the easiest order, and the Lex order is the hardest. Elimination orders stand in between, depending on how many variables we wish to eliminate.

By using a change of order algorithm, we may first compute a Gröbner basis for an "easy" order and then use the change of order to compute the wanted basis.

Algorithm FGLM proceeds as follows: let $I$ be a zero-dimensional of $\mathbb{K}[\mathrm{X}]$, let $G_{1}$ be a Gröbner basis of $I$ for some monomial order $<_{1}$ and let $<_{2}$ be another order for which we want a Gröbner basis of $I$. Since $I$ is zero-dimensional, $\mathbb{K}[X] / I$ is a vector space of finite dimension $D=\operatorname{deg}(I)$. The algorithm will consider two monomial bases for this vector space, namely the two staircases for the monomial orders at hand.

Definition 2.44. Let $I$ be a zero-dimensional ideal of $\mathbb{K}[X]$, and let $<_{\text {mon }}$ be a monomial order. The staircase of $I$ with respect to $<_{\text {mon }}$ is the set of monomials which are not leading monomial of a polynomial in $I$ :

$$
\mathcal{S}_{I}=\{m \in \mathcal{M} \mid m \notin \mathrm{LT}(I)\}
$$

If $G$ is a Gröbner basis of $I$ with respect to $<_{\text {mon }}$, then

$$
\mathcal{S}_{I}=\{m \in \mathcal{M} \mid \forall g \in G, \mathrm{LT}(g) \nmid m\} .
$$

The main idea of Algorithm FGLM (Algorithm 2.5) is to use linear algebra to model the structure of the quotient ring:

1. Use $G_{1}$ to compute the staircase $\mathcal{S}_{1}$ of $I$ with respect to $<_{1}$;
2. Compute the matrices of the multiplication by $x_{1}, \ldots, x_{n}$ in the basis $\mathcal{S}_{1}$;
3. Compute the staircase $\mathcal{S}_{2}$ of $I$ with respect to $<_{2}$

The two first steps are performed by the subroutine MultiplicationMatrix (Algorithm 2.4): it is a routine taking as input a Gröbner basis $G$ of an ideal $I$ and returning the matrices of all linear maps $\mathbb{K}[\mathrm{X}] / I \longrightarrow \mathbb{K}[\mathrm{X}] / I$ given by the multiplication by $X_{i}$ followed by the normal form in the monomial basis $\mathcal{S}$, for $i$ ranging in $\{1, \ldots, n\}$. The number of normal forms to compute can be limited by noticing that any monomial $X_{i} \mu, \mu \in \mathcal{S}$ falls in one of the following cases:

- $X_{i} \mu \in \mathcal{S}$, in which case $\operatorname{NF}\left(X_{i} \mu\right)=X_{i} \mu$;
- $X_{i} \mu=\operatorname{LT}(g), g \in G$, in which case $\operatorname{NF}\left(X_{i} \mu\right)=g-X \mu$;
- otherwise, if $\operatorname{NF}(\mu)=\sum_{k} \alpha_{k} \mu_{k}$, then $\operatorname{NF}\left(X_{i} \mu\right)=\sum_{k} \alpha_{k} \operatorname{NF}\left(X_{i} \mu_{k}\right)$.

In the two former cases, the normal form can be computed for free. In the latter case, it is given as a linear combination of normal forms of smaller monomials (according to the monomial order).

```
Algorithm 2.4 MultiplicationMatrix
Input: \(G\) a Gröbner basis of a zero-dimensional ideal \(I\) for a monomial order \(<_{\text {mon }}\)
Output: \(\mathbf{M}=\left(M_{1}, \ldots, M_{n}\right)\) multiplication matrices of \(X_{1}, \ldots, X_{n}\)
    \(\mathcal{S} \leftarrow\{\mu\) monomial of \(\mathbb{K}[\mathrm{X}] \mid \mu \notin \mathrm{LT}(\langle G\rangle)\}\)
    \(F \leftarrow\left\{X_{i} \mu, i \in\{1, \ldots, n\}, \mu \in \mathcal{S}\right\} \backslash \mathcal{S}\)
    \(N \leftarrow[]\)
    for \(\mu \in \mathcal{S}\) do \(N[\mu] \leftarrow \mu\) end for
    Sort \(F\) by increasing \(<_{\text {mon }}\) order
    for \(\mu \in F\) do
        if \(\mu \in \mathrm{LT}(G)\) then
            Find \(g\) such that \(g \in G\) and \(\mu=\operatorname{LT}(g)\)
            \(N[\mu] \leftarrow g-\mu\)
            else
                Find ( \(i, \mu^{\prime}\) ) such that \(\mu=X_{i} \mu^{\prime}, \mu^{\prime} \in F\)
                Write \(N\left[\mu^{\prime}\right]=\sum_{\mu^{\prime \prime}<\mu^{\prime}} \alpha_{\mu^{\prime \prime}} \mu^{\prime \prime}\)
                \(N[\mu] \leftarrow \sum_{\mu^{\prime \prime}<\mu^{\prime}} \alpha_{\mu^{\prime \prime}} N\left[X_{i} \mu^{\prime \prime}\right]\)
        end if
    end for
    Copy the normal forms \(N\) into the multiplication matrices \(\mathbf{M}\)
    return M
```

Algorithm FGLM (Algorithm 2.5) then computes the staircase $\mathcal{S}_{2}$ of the ideal with respect to $<_{2}$ maintaining a association between the monomials $\mu$ and their normal form $v_{\mu}$ w.r.t. $G_{1}$.

For this purpose, it considers, in increasing $<_{2}$ order, all monomials already known to be in $\mathcal{S}_{2}$ (starting with 1) and considers their smallest multiples. Let $\mu$ be a monomial in the staircase, and let $X_{i}$ be a variable. The normal form of $X_{i} \mu$ with respect to $G_{1}$ can be computed by matrix multiplication:

$$
v_{X_{i} \mu}:=\mathrm{NF}\left(X_{i} \mu\right)=M_{i} \cdot v_{\mu} .
$$

If this normal form is linearly independent from all previously computed normal forms, this means that $X_{i} \mu$ lies in the staircase $\mathcal{S}_{2}$. Otherwise, $X_{i} \mu$ lies in the initial ideal of $I$ w.r.t. $<_{2}$, and the linear relation between $\left(v_{s}\right)_{s \in \mathcal{S}_{2}}$ and $v_{X_{i} \mu}$ gives a polynomial of $I$ with leading term $X_{i} \mu$ : this polynomial may be added to the Gröbner basis of I w.r.t. $<_{2}$.

```
Algorithm 2.5 FGLM ([Fau+93])
Input: \(G_{1}\) a Gröbner basis of \(I\) for a monomial order \(<_{1}\), another monomial order \(<_{2}\)
Output: \(G_{2}\) Gröbner basis of \(I\) for \(<_{2}\)
    if \(G_{1}=\{1\}\) then return \(\{1\}\) end if
    \(\mathbf{M}=\left(M_{1}, \ldots, M_{n}\right) \leftarrow\) MultiplicationMatrix \(\left(G_{1}\right)\)
    \(G_{2} \leftarrow \varnothing\)
    \(\mathcal{S}_{2} \leftarrow\{1\}\)
    \(v_{1} \leftarrow(1,0, \ldots, 0)\)
    while \(\exists i \in\{1, \ldots, n\}, \mu \in \mathcal{S}_{2}, X_{i} \mu \notin \mathcal{S}_{2}\) do
        \(\left(X_{i}, \mu\right) \leftarrow \min \left\{\left(X_{i}, \mu\right) \mid i \in\{1, \ldots, n\}, \mu \in \mathcal{S}_{2}, X_{i} \mu \notin \mathcal{S}_{2} \cup \mathrm{LT}_{<_{2}}\left(G_{2}\right)\right\}\)
        \(v_{X_{i} \mu} \leftarrow M_{i} \cdot v_{\mu}\)
        if \(v_{X_{i} \mu} \in \operatorname{Vect}\left\{v_{s} \mid s \in \mathcal{S}_{2}\right\}\) then
            Find \(\left(\lambda_{s}\right)_{s \in \mathcal{S}_{2}}\) such that \(v_{X_{i} \mu}=\sum_{s \in \mathcal{S}_{2}} \lambda_{s} v_{s}\)
            Add \(X_{i} \mu-\sum_{s \in \mathcal{S}_{2}} \lambda_{s} s\) to \(G_{2}\)
        else
            Add \(\mu\) to \(\mathcal{S}_{2}\)
        end if
    end while
    return \(G_{2}\)
```


### 2.5. Complexity results

### 2.5.1. Complexity model and notations

In this section, we give complexity bounds for Algorithms Matrix- $F_{5}$ and FGLM. We measure the arithmetic complexity of algorithms, that is the number of operations (additions, multiplications, inversions) in the coefficient field $\mathbb{K}$.

This corresponds to time complexity if we assume that field operations are done in constant time. This is a reasonable assumption for finite fields. On the other hand, for infinite fields, it means that we ignore the growth of the coefficients of the polynomials.

Because of the central role of linear algebra in Gröbner basis algorithms, complexity estimates shall frequently depend on the constant $\omega$ defined as the smallest positive number such that one can compute the product of two $k \times k$ matrices in time $O\left(k^{\omega}\right)$.

Most operations on $k \times k$ matrices can be reduced to matrix multiplications, and thus performed in time $O\left(k^{\omega}\right)$ : this includes the row echelon form, the characteristic polynomial...

The exact value of $\omega$ is unknown. The best known theoretical bound at the time of writing is $\omega<2.3728639$ ([LeGalli4]). In practice, matrix multiplication is typically done either

- using the naïve sequential implementation, with a complexity exponent of 3 , or
- using Strassen's algorithm ([Str69]), with a complexity exponent of $\log _{2}(7) \approx 2.81$.

Finally, it is important to note that complexity estimates are always given under genericity assumptions. The worst-case complexity of computing a Gröbner basis is known to be doublyexponential in the degree of the polynomials ([MM82; BS88]), but this bound is not reached for generic systems.

### 2.5.2. Complexity of Matrix- $F_{5}$ and degree of regularity

In the case of Matrix- $\mathrm{F}_{5}$, at each degree $d$, we build incrementally a matrix $M_{d}$ whose columns are indexed by the monomials of degree $d$ in $\mathbb{K}[\mathrm{X}]$. This matrix is then reduced to echelon form, and this step is the most expensive of the algorithm. Hence the complexity of the algorithm depends mainly on the size of the matrices being built, and on the number of such matrices. Both depend on the highest degree reached when running the algorithm.

Recall that if the system $F$ forms a regular sequence, then the $\mathrm{F}_{5}$ criterion eliminates all reductions to zero. As a consequence, the number of rows of the matrices is necessarily less than the number of columns.

As for the number of matrices that we need to compute and reduce, it is bounded by the stopping degree $d_{\text {max }}$, given as an input to the algorithm. All in all, it gives the following simple estimate for the complexity of algorithm Matrix- $\mathrm{F}_{5}$ :

Theorem 2.45 ([BFS14, Prop. 1]). Let $F \subset \mathbb{K}[\mathrm{X}]$ be a homogeneous polynomial system and $d_{\text {max }} \in \mathbb{N}$, then Matrix $-F_{5}\left(F, d_{\max }\right)$ computes a $d_{\max }-$ Gröbner basis of $\langle F\rangle$ in time

$$
O\left(d_{\max }\binom{n+d_{\max }-1}{d_{\max }}^{3}\right)
$$

We now focus on a typical usage of algorithm $F_{5}$ : computing a GRevLex basis for a system defined by a sufficiently generic system. In this case, we want to estimate how large $d_{\text {max }}$ needs to be in order for Matrix- $\mathrm{F}_{5}$ to return a Gröbner basis of the ideal.

Definition 2.46. Let $F$ be a system of homogeneous polynomials. The degree of regularity of $F$ is the smallest degree $d_{\text {reg }}$ such that Matrix $-F_{5}\left(F, d_{\text {reg }},<_{\text {grevlex }}\right)$ gives a Gröbner basis of $\langle F\rangle$.

Recall that we defined the index of regularity of an ideal as the degree of its Hilbert series (Definition 1.32). This notion gives a bound for the degree of regularity of regular sequences, under genericity hypotheses:

Theorem 2.47 (Macaulay's bound). Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a system of homogeneous polynomials with respective degree $\left(d_{1}, \ldots, d_{m}\right)$, and assume that $F$ is in Noether position with respect to the variables $X_{m+1}, \ldots, X_{n}$. Then

$$
d_{\mathrm{reg}}(F) \leq i_{\text {reg }}(\langle F\rangle)+1=\sum_{i=1}^{m}\left(d_{i}-1\right)+1 .
$$

Remark 2.48. Recall that if $m=n$, that is if the ideal $\langle F\rangle$ is zero-dimensional, the condition of being in Noether position is equivalent to $F$ being a regular sequence.

Proof. First assume that $m=n$. Then the ideal $\langle F\rangle$ is zero-dimensional, and the Hilbert series of $\mathbb{K}[\mathrm{X}] /\langle F\rangle$ is a polynomial with degree $i_{\text {reg }}$. If $d \geq i_{\text {reg }}+1$, all monomials in $\mathbb{K}[\mathrm{X}]$ are in $\langle F\rangle$, and so any polynomial with degree greater than $i_{\text {reg }}+1$ is reducible modulo some polynomial in $\langle F\rangle$ with lower degree.

Now assume that $m<n$. Since $F$ is in Noether position, from Proposition 2.9, if $G$ is a Gröbner basis of $F$, then $G \cup\left\{X_{m+1}, \ldots, X_{n}\right\}$ is a Gröbner basis of $\left\langle F, X_{m+1}, \ldots, X_{n}\right\rangle$. Since this is a zero-dimensional regular sequence, the bound above applies.

Theorem 2.49 ([ $\mathrm{BFS}_{14}$, Prop. 1]). Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a system of homogeneous polynomials. Then a run of Matrix- $\mathrm{F}_{5}\left(F, d_{\max }\right)$ computes a $d_{\max }-$ Gröbner basis of $\langle F\rangle$ in time

$$
\begin{equation*}
O\left(d_{\max }\binom{n+d_{\max }-1}{d_{\max }}^{3}\right) \tag{2.2}
\end{equation*}
$$

If $d_{\max } \geq d_{\mathrm{reg}}$ and the order is GREVLEX, then Matrix- $\mathrm{F}_{5}\left(F, d_{\max }\right)$ computes a Gröbner basis of $F$. If $F$ is in Noether position with respect to the variables $X_{m+1}, \ldots, X_{n}$, then $d_{\mathrm{reg}}$ is bounded by

$$
d_{\mathrm{reg}} \leq \sum_{i=1}^{m}\left(d_{i}-1\right)+1
$$

In the affine case, for systems regular in the affine sense, Proposition 2.42 ensures that we can run AlgorithmMatrix- $F_{5}$ on the homogenized system up to the degree of regularity of the highest degree components, and obtain a Gröbner basis. So the same complexity bounds hold, using the number of monomials of degree at most $d$ for the size of the matrix at degree $d$.

Theorem 2.50. Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a system of polynomials which are regular in the affine sense, and assume that we are using the GREvLEX order. Then Matrix- $\mathrm{F}_{5}\left(F, d_{\max }\right)$ computes a $d_{\text {max }}$-Gröbner basis of $\langle F\rangle$ in time

$$
O\left(d_{\max }\binom{n+d_{\max }}{d_{\max }}^{3}\right)
$$

If $d_{\max } \geq d_{\mathrm{reg}}$, then Matrix- $\mathrm{F}_{5} F, d_{\max }$ computes a Gröbner basis of $F$. If the highest degree components of $F$ are in Noether position with respect to the variables $X_{m+1}, \ldots, X_{n}$, then $d_{\mathrm{reg}}$ is bounded by

$$
d_{\mathrm{reg}} \leq \sum_{i=1}^{m}\left(d_{i}-1\right)+1
$$

### 2.5.3. Thin-grained complexity of Matrix- $F_{5}$

The complexity bound in the previous section is far from optimal: it is obtained by taking the complexity of reducing one dense matrix as large as the largest one the algorithm sees, and multiplying this complexity by the number of matrices that the algorithm needs to reduce.

There are two ways this bound can be improved: by taking into account the fact that the matrices are not dense, and by looking at the size of all matrices instead of the largest one.

Both refinements come as a consequence of the following structure lemma:
Lemma 2.51 (Structure lemma, [BFS14, Prop. 11]). Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a homogeneous system of polynomials of respective degree $\left(d_{1}, \ldots, d_{m}\right)$, in simultaneous Noether position with respect to the order $X_{1}>\cdots>X_{m}$. For all $i \in\{1, \ldots, m\}$, let $G_{i}$ be the reduced GREvLEX Gröbner basis of $\left\langle f_{1}, \ldots, f_{i}\right\rangle$. Let $i \in\{1, \ldots, m\}$, and $g \in G_{i}$ with signature $(j, \mu)$. Then:

- $j \leq i$
- $\operatorname{LT}(g) \in A_{i}:=\left[X_{1}, \ldots, X_{i}\right]$
- $\mu \in A_{i-1}=\mathbb{K}\left[X_{1}, \ldots, X_{i-1}\right]$

Proof. Since $G_{i}$ is a Gröbner basis of $\left\langle f_{1}, \ldots, f_{i}\right\rangle$, any polynomial in $G_{i}$ can be written as an algebraic combination of these polynomials, and thus $j \leq i$.

For the second statement, without loss of generality, we may assume that $i=j$. Indeed, if $g$ has signature $(j, \mu)$ with $j<i$, then $g \in G_{j}$ and applying the proposition to $G_{j}$ yields that $\operatorname{LT}(g) \in \mathbb{K}\left[X_{1}, \ldots, X_{j}\right] \subset \mathbb{K}\left[X_{1}, \ldots, X_{i}\right]$. From Corollary 2.9, $\operatorname{LT}\left(\left\langle f_{1}, \ldots, f_{i}\right\rangle\right)$ is generated by monomials in $A_{i}$, and since $g \in G_{i}, \mathrm{LT}(g) \in A_{i}$.

In order to prove the third statement, we proceed by induction on the signature $(i, \mu)$. For $i=1$, there is no reduction so the result holds. Let us assume $i>1$, then the lemma is true for the minimal valid signature $(i, 1)$. Let now $f$ be a polynomial appearing in the computation of $G_{i}$, in $W$-degree $d>d_{i}$. Before any reduction, $f$ can be written as $X_{l} f_{0}$, where $f_{0}$, with signature $(i, \mu)$ is reduced against the basis $G_{i-1}$, and where $X_{l}$ is a variable with $l \leq i$.

If there is a polynomial in $G_{i}$ with signature $(i, \mu)$, it means that the row corresponding to $f$ was reducible (in the linear sense) by the rows above it. Algebraically, it means that there exists $g$ in $G_{k}, k \leq i$, having signature $(k, \beta)<(i, \mu)$, such that $f$ is head-reducible modulo $g$. By induction hypothesis, we find that $\beta \in A_{k-1}$, and from the second statement, we deduce that $\mathrm{LT}(g) \in A_{k}$. Since $g$ head-reduces $f, \mathrm{LT}(g)$ divides $\operatorname{LT}(f)=X_{l} \mathrm{LT}\left(f_{0}\right)$, and since $f_{0}$ is reduced, $X_{l}$ divides $\mathrm{LT}(\mathrm{g})$, and so $X_{l} \in A_{k}$. Now we are in one of the three following cases:

1. $k<i$ : then $X_{l} \in A_{i-1}$, and the label of $f$ is $\left(i, X_{l} \mu\right)$ with $X_{l} \mu \in A_{i-1}$;
2. $k=i, l<i$ : same;
3. $k=i, l=i$ : then we can write $X_{i} \operatorname{LT}\left(f_{0}\right)=c \mathrm{LT}(g)$, with $c \in A_{i}$. If $X_{i} \mid c, \operatorname{LT}(g)$ divides $\frac{X_{i}}{X_{i}} \operatorname{LT}\left(f_{0}\right)$, which is absurd because we assumed $f_{0}$ to be reduced. So we must have $c \in A_{i-1}$. Then the label of $c g$ is $(i, c \beta)$, which is greater than the label of $f: c \beta$ is in $A_{i-1}$, unlike $X_{i} \mu$. But that is absurd, because we assumed $c g$ to reduce $f$.

The structure lemma gives a correspondence between signatures and polynomials in the Gröbner basis. Counting the possible signatures thus gives us a bound on the number of polynomials in the Gröbner basis:

Proposition 2.52 ([BFS14, Th. 12]). With the same notations, let $i \in\{1, \ldots, m\}$ and $d \in \mathbb{N}$. The number of polynomials of degree $d$ in $G_{i}$ whose leading term does not belong to $\operatorname{LT}\left(G_{i-1}\right)$ is bounded by the generic term $b_{d}^{(i)}$ of the following generating series:

$$
B_{d}^{(i)}(T):=\sum_{d=0}^{\infty} b_{d}^{(i)} T^{d}:=T^{d_{i}} \prod_{k=1}^{i-1} \frac{1-T^{d_{k}}}{1-T}
$$

Remark 2.53. This proposition bounds the number of new elements added to the Gröbner basis by each polynomial $f_{i}$. It is tempting to write it as $G_{i} \backslash G_{i-1}$, but elements of $G_{i-1}$ may be tail-reducible modulo $f_{i}$, changing their value in $G_{i}$. Another way of stating this proposition is that we count the number of polynomials in $G_{i}$ with signature $(i, \bullet)$.

Proof. By Lemma 2.51, any polynomial in $G_{i}$ with signature $(i, \mu)$ is such that $\mu \in A_{i-1}=$ $\mathbb{K}\left[X_{1}, \ldots, X_{i-1}\right]$. Furthermore, the $\mathrm{F}_{5}$ criterion ensures that $\mu \notin F_{i-1}:=\left\langle f_{1}, \ldots, f_{i-1}\right\rangle$. The number of monomials of degree $d$ in $A_{i-1}$ such that $\mu \notin\left\langle f_{1}, \ldots, f_{i-1}\right\rangle$ is precisely the dimension of $\left(A_{i-1} /\left\langle\theta_{i-1}\left(F_{i-1}\right)\right\rangle\right)_{d}$ as a $\mathbb{K}$-vector space. In other words, the number of available monomials for a signature of degree $d$ is bounded by the generic term of the Hilbert series of $A_{i-1} /\left\langle\theta_{i-1}\left(F_{i-1}\right)\right\rangle$. Since $F$ is in simultaneous Noether position, by Proposition 1.44, $\theta_{i-1}\left(F_{i-1}\right)$ is a regular sequence, and this Hilbert series is:

$$
\mathrm{HS}_{A_{i-1} /\left\langle\theta_{i-1}\left(F_{i-1}\right)\right\rangle}=\frac{\prod_{k=1}^{i-1} 1-T^{d_{k}}}{(1-T)^{i-1}}
$$

and the result follows, shifting this series by $T^{d_{i}}$ to count the number of polynomials that we can form with these signatures.

Theorem 2.54 ([BFS14, Sec. 3.2, Eq. 10]). Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a system of homogeneous polynomials with respective degrees $d_{1}, \ldots, d_{m}$, and such that $F$ is in simultaneous Noether position with respect to the order $X_{1}>\cdots>X_{n}$.

Then a run of Matrix- $\mathrm{F}_{5}\left(F, d_{\max }\right)$ computes a $d_{\max }$-Gröbner basis of $\langle F\rangle$ in time

$$
\begin{equation*}
O\left(\sum_{i=1}^{m} \sum_{d=0}^{d_{\mathrm{reg}}} b_{d}^{(i)}\binom{i+d-1}{d}\binom{n+d-1}{d}\right) \tag{2.3}
\end{equation*}
$$

Proof. For each polynomial $f_{i}$, at each degree $d$, we reduce at most $b_{d}^{(i)}$ rows in the matrix. Each reduction uses at most $\binom{i+d-1}{d}$ pivots (because the resulting leading term is in $\mathbb{K}\left[X_{1}, \ldots, X_{i}\right]$ ), and each elementary reduction costs at most $\binom{n+d-1}{d}$ arithmetic operations (the number of columns in the matrix).

Both bounds are plotted for comparison in Figure 2.1.


Figure 2.1: Number of operations in a run of algorithm Matrix- $F_{5}$ on a system of 3 generic homogeneous polynomials with degree $d$

### 2.5.4. Complexity of FGLM

As before, let $I$ be a zero-dimensional ideal and let $D$ be its degree.
Algorithm FGLM has two main steps: computing the multiplication matrices, and computing the Gröbner basis with respect to the new monomial order.

The cost of computing the multiplication matrices with Algorithm 2.4 is concentrated in the computation of the normal forms of the monomials of the form $X_{i} \mu^{\prime}, \mu^{\prime} \in F$. Each of these normal forms is computed as the product of a $K \times K$ matrix and a $1 \times K$ vector, where $K \leq D$.We need to compute at most $\# F \leq n D$ such normal forms, and overall, computing the multiplication matrices can be done in time

$$
O\left(n D \times D^{2}\right)=O\left(n D^{3}\right) .
$$

Computing the Gröbner basis is done by incrementally building a linear basis of the staircase of $I$ with respect to the second order. If we maintain a matrix of the basis throughout the algorithm, each step of verification of linear independence and recovery of the coefficients $\lambda_{j}$ can be done by reducing one vector by this matrix. Overall, it is equivalent to computing the normal form of the whole matrix using the Gauss algorithm. This matrix has $D$ columns and at most $n D$ rows, so the complexity is again $O\left(n D^{3}\right)$.

Theorem 2.55 ([Fau+93]). The complexity of algorithm FGLM for computing the Gröbner basis of a zero-dimensional ideal with degree $D$ is

$$
O\left(n D^{3}\right) .
$$

Faster variants of FGLM have been proposed. In [FM13], the authors proposed algorithms taking advantage of the sparsity of the multiplication matrices. In particular, for a lexicographical basis in shape position, it suffices to compute the minimal polynomial of $M_{n}$, and this can be
done in time $O\left(n D^{\omega}\right)$ using fast linear algebra. In [Fau+14], the authors show that, up to a generic change of coordinates, the multiplication matrices can be read at no cost on a GRevLex basis of the ideal. Overall, up to a generic change of coordinates, computing a Lex Gröbner basis from a GRevLex one can be done in time

$$
O\left(n D^{\omega}\right) .
$$

## Part II

## Contributions

## Chapter 3

## Weighted homogeneous systems

The results presented in this chapter are extracted from a joint work with fean-Charles Faugère and Mohab Safey El Din. They have been published in [FSV13] and [FSV15].

In this chapter, we consider weighted-homogeneous systems, as defined in Section 1.1.3. We define some properties of these systems, and we rework some characterizations which were available for homogeneous systems. We exhibit a computational strategy making use of existing algorithms, and show, under genericity assumptions, that the complexity is divided by $\left(\Pi w_{i}\right)^{\omega}$ when compared with homogeneous systems with the same degree. Finally, we provide some experimental results showing that taking advantage of the weighted homogeneous structure can yield substantial speed-ups.

### 3.1. Introduction

Consider a vector of positive integers $W=\left(w_{1}, \ldots, w_{n}\right)$. As in Section 1.1.3, we define a gradation on $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$ by defining the $W$-degree of a monomial as:

$$
\operatorname{deg}_{W}\left(X_{1}^{\alpha_{1}} \cdots X_{n}^{\alpha_{n}}\right)=w_{1} \alpha_{1}+\cdots+w_{n} \alpha_{n}
$$

An equivalent definition is that a polynomial $f$ is $W$-homogeneous with $W$-degree $d$ if and only if $f\left(X_{1}^{w_{1}}, \ldots, X_{n}^{w_{n}}\right)$ is homogeneous with degree $d$.

There are two main strategies for computing Gröbner bases for weighted homogeneous systems:

- apply existing algorithms, disregarding the weighted homogeneous structure;
- use the above property to transform a weighted homogeneous system and apply existing algorithms to this system, exploiting its new homogeneous structure.

Experimentally, the second strategy appears to be faster than the first one for computing a GRevLex basis using $F_{4}$ or $F_{5}$. But when needed, the change of order using FGLM becomes a major bottleneck.

To the best of our knowledge, there is no way to evaluate the complexity of the first strategy, and complexity analyses for the second strategy only rely on the homogeneous structure of the system.

In this work, we justify that the second strategy is correct, and is actually a way to efficiently take advantage of the weighted homogeneous structure with pairs-based algorithms, such as $F_{4}$ or $F_{5}$. We also explain the bottleneck in the FGLM algorithm for zero-dimensional systems, and propose a computational strategy working around this obstacle, effectively dividing the complexity of FGLM by a factor $\left(\prod w_{i}\right)^{3}$.

We also give complexity results for the $F_{5}$ algorithm. For this purpose, we adapt algorithm Matrix- $\mathrm{F}_{5}$ to the weighted case, and show that the bounds of Proposition 1.20 on the number of monomials in a $W$-graded algebra allow to divide the complexity bounds by $\left(\Pi w_{i}\right)^{3}$. We also show that the complexity in practice is improved, because the degree of regularity of a weighted homogeneous system is lower than that of a homogeneous system with the same degrees.

As in the homogeneous case, complexity studies are obtained under some generic assumptions. More precisely, we give bounds on the degree of regularity for:

- zero-dimensional systems defined by a regular sequence;
- positive-dimensional systems defined by a sequence in Noether position with respect to the smallest variables;
- zero- or positive-dimensional systems defined by a sequence in simultaneous Noether position;

In the over-determined case, we define semi-regular weighted homogeneous sequences, and we show that under some hypotheses on the weights and the degrees, this property admits the same characterization as in the homogeneous case. This can be used to bound the degree of regularity of weighted semi-regular sequences. We give an example of such computations, by transposing an asymptotic analysis [Bar+05] of the degree of regularity of semi-regular sequences with $m=n+k$ to the weighted homogeneous case.

Prior work The special case $W=(1, \ldots, 1)$ is the usual homogeneous case. In this case, all properties and characterizations shown in this chapter specialize to known results that we described in Chapters 1 and 2.

Weighted homogeneous systems have been studied before, from the angle of singularity theory and commutative algebra. In particular, the results about the Hilbert series and the Hilbert function of weighted homogeneous ideals (see Chapter 1 ) can be found in most commutative algebra textbooks.

The computational strategy for systems with a weighted structure is not new either. For example, it is already implemented (partially: for weighted homogeneous or heuristically, systems containing some weighted homogeneous polynomials) in the computer algebra system Magma [Magma]. Additionally, the authors of [Tra96] proposed another way of taking into account the weighted structure, using the Hilbert series of the ideal. The authors of [CDR96] generalized this algorithm to systems homogeneous with respect to a multigraduation. Their definition of a system of weights is more general than the one we use in this thesis.

To the best of our knowledge, nobody presented a formal description of a computational strategy for systems with a weighted homogeneous structure (not necessarily weighted homogeneous), together with complexity estimates.

Some of the results presented in this chapter appeared in a conference paper [FSV13]: these results are the weak form of the weighted Macaulay bound (3.2) and the formal description of the algorithmic strategy for weighted homogeneous systems, with the complexity estimates (3.1) and (3.4). The remaining parts of this chapter were presented in the journal article [FSV 15 ]: the conference paper lacked a hypothesis (reverse chain-divisible systems of weights), and therefore lacked the precise description of Hilbert series required to obtain results for semi-regular sequences. The sharp variant of the weighted Macaulay bound (3.3), under the assumption of simultaneous Noether position, was also added in the extended paper. Finally, the benchmarks section of the journal paper contains additional systems, arising in polynomial inversion problems.
The conference paper was using quasi-homogeneous instead of weighted homogeneous to describe the studied structure. While both names exist in the literature, weighted homogeneous seems to be more common, and to better convey the notion that this structure is a generalization of homogenity, instead of an approximation. The same notion is sometimes also named simply homogeneous (in which case the weights are determined by the degree of the generators; see for example [Eis95]), or homogeneous for a nonstandard graduation [DSo6].

Main results By definition, weighted homogeneous polynomials can be made homogeneous by raising all variables to their weight. The resulting system can then be solved using algorithms for homogeneous systems. However, it appears experimentally that solving such systems is much faster than solving generic homogeneous systems. In this chapter, we show that the complexity estimates for homogeneous systems can be divided by $\left(\Pi w_{i}\right)^{\omega}$ in case the system was originally $W$-homogeneous, $\omega$ being the complexity exponent of linear algebra operations ( $\omega=3$ for naive algorithms, such as the Gauss algorithm).
These complexity estimates depend on two parameters of the system: its degree of regularity $d_{\text {reg }}$ and its degree $\operatorname{deg}(I)$. Both parameters can be obtained from the Hilbert series of the ideal, which can be precisely described under generic assumptions. To be more specific, we will consider systems defined by a regular sequence and systems which are in simultaneous Noether position.

Theorem. Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a system of weights, and $F=\left(f_{1}, \ldots, f_{m}\right)$ a zero-dimensional $W$-homogeneous system of polynomials in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$, with respective $W$-degree $d_{1}, \ldots, d_{m}$. The complexity (in terms of arithmetic operations in $\mathbb{K}$ ) of Algorithm $F_{5}$ to compute a $W$-GREVLEX Gröbner basis of $I:=\langle F\rangle$ is bounded by

$$
\begin{equation*}
C_{\mathrm{F}_{5}}=O\left(\frac{1}{\left(\prod w_{i}\right)^{3}} \cdot\binom{n+d_{\mathrm{reg}}-1}{d_{\mathrm{reg}}}^{3}\right) . \tag{3.1}
\end{equation*}
$$

If $F$ is a regular sequence (and in particular $m=n$ ), then $d_{\text {reg }}$ can be bounded by the weighted Macaulay bound:

$$
\begin{equation*}
d_{\mathrm{reg}} \leq \sum_{i=1}^{n}\left(d_{i}-w_{i}\right)+\max \left\{w_{j}\right\} . \tag{3.2}
\end{equation*}
$$

If additionally $F$ is in simultaneous Noether position w.r.t the order $X_{1}>\cdots>X_{n}$, then the weighted Macaulay bound can be refined:

$$
\begin{equation*}
d_{\mathrm{reg}} \leq \sum_{i=1}^{n}\left(d_{i}-w_{i}\right)+w_{n} \tag{3.3}
\end{equation*}
$$

The complexity of Algorithm FGLM to perform a change of ordering is bounded by

$$
\begin{equation*}
C_{F G L M}=O\left(n(\operatorname{deg}(I))^{\omega}\right) . \tag{3.4}
\end{equation*}
$$

Recall that if $F$ forms a regular sequence, then $\operatorname{deg}(I)$ is given by the weighted Bézout bound

$$
\operatorname{deg}(I)=\frac{\prod_{i=1}^{n} d_{i}}{\prod_{i=1}^{n} w_{i}}
$$

In particular, the bound (3.3) indicates that in order to compute a Gröbner basis faster for a generic enough system, one should order the variables by decreasing weights whenever possible.

The hypotheses of the theorem are not too restrictive. In the homogeneous case, regularity and simultaneous Noether position are generic properties. However, in the weighted homogeneous case, there are systems of weights and systems of weighted degrees for which they are not generic. We shall identify large families of systems of weights and systems of weighted degrees for which they are (Proposition 3.19).

All sequences in simultaneous Noether position are regular. Conversely, in the homogeneous case, all regular sequences are in simultaneous Noether position up to a generic linear change of coordinates. In the weighted homogeneous case, this is no longer true. In fact, there are even systems of weights for which there exists no non-trivial change of coordinates.

To avoid this limitation, we consider reverse chain-divisible systems of weights, that is systems of weights such that $w_{n}\left|w_{n-1}\right| \cdots \mid w_{1}$. This property ensures that there are non-trivial changes of coordinates of the form $X_{i} \leftarrow X_{i}+P_{i}\left(X_{i+1}, \ldots, X_{n}\right)$ for all $i$, with $P_{i}$ a $W$-homogeneous polynomial with $W$-degree $w_{i}$. Under this assumption, many properties from the homogeneous case remain valid in a weighted setting. In particular, any regular sequence is in simultaneous Noether position up to a $W$-homogeneous change of coordinates (Theorem 3.15).

For many systems from practical applications, the weights can be chosen to be reverse chain-divisible. We give a few examples in the last section of this chapter.

If $m>n$, there is no regular sequence. Instead, we will consider systems defined by a semi-regular sequence, that is systems for which no reduction to zero appear in a run of Algorithm $F_{5}$. This property has several equivalent definitions in the homogeneous case. While these definitions can be easily extended to the weighted case, their equivalence is not necessarily true. However, we prove that these definitions are equivalent in the special case where the weights form a reverse chain-divisible sequence.

In the homogeneous case, the property of being semi-regular is only conjectured to be generic, but this conjecture is proved in a handful of cases [Mor96, Thm. 1.5]. In this chapter, we adapt the proof of one of these cases, namely the case $m=n+1$ in a base field of characteristic 0 .

For semi-regular systems with $m=n+1$, we obtain a bound on the degree of regularity. More generally, in the homogeneous case, one can compute asymptotic estimates on the degree
of regularity of a semi-regular sequence [Bar+05; Baro4]. These estimates can be adapted to the weighted homogeneous case. As an example, we give an asymptotic bound on the degree of regularity for semi-regular systems with $m=n+k$ for a given integer $k$ :

Theorem. Let $n$ and $k$ be two positive integers, and let $m=n+k$. Let $w_{0}$ and $d_{0}$ be two positive integers such that $w_{0} \mid d_{0}$. Consider the system of $n$ weights $W=\left(w_{0}, \ldots, w_{0}, 1\right)$. Let $F$ be a semi-regular sequence in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$, made of $W$-homogeneous polynomials with $W$-degree $d_{0}$. Then the highest degree reached in the computation of a $W$-GREvLEX Gröbner basis of $\langle F\rangle$ is asymptotically bounded by

$$
d_{\mathrm{reg}}=n \frac{d_{0}-w_{0}}{2}-\alpha_{k} \sqrt{n \frac{d_{0}^{2}-w_{0}^{2}}{6}}+O\left(n^{1 / 4}\right)
$$

where $\alpha_{k}$ is the largest root of the $k$ 'th Hermite polynomial.
Experimentally, if we lift the assumption that the system of weights is reverse chain-divisible, the degree of regularity does not appear to rise too far beyond the bound. Future work on the topic could include characterizing the Hilbert series of $W$-homogeneous semi-regular sequences in full generality, in order to obtain bounds on the $W$-degree of regularity.

In practice, taking advantage of the weighted structure when applicable yields significant speed-ups. Some instance of a weighted structure has already been successfully exploited for an application in cryptography [Fau+13]. We also present timings obtained with several polynomial inversion problems, with speed-ups ranging from $1-2$ to almost 100 . In particular, we use these techniques in order to compute the relations between fundamental invariants of several groups [Stuo8]. For some groups such as the Cyclic-5 group or the dihedral group $D_{5}$, computing these relations is intractable without considering the weighted structure of the system, while it takes only a few seconds or minutes when the weighted structure is exploited. All these systems are examples of applications where the weights giving the appropriate $W$ homogeneous structure are naturally reverse chain-divisible. These experimentations have been carried using $F_{5}$ and FGLM with the Gröbner basis library FGb [Fau10] and $F_{4}$ with the computer algebra system Magma [Magma].

There are other applications where Gröbner bases are computed for polynomial systems with a weighted-homogeneous structure, for example in coding theory, both for generating codes [BP99, sec. 5], [Leoo9] and for decoding through Guruswami-Sudan's algorithm (see [GRo9] for an overview).

Organisation of the chapter In Section 3.2, we describe some algebraic properties of $W$ homogeneous systems. In particular, we show that regular sequences and sequences in Noether position are generic, and we give a characterization of the Hilbert series of a $W$-homogeneous regular sequence. In Section 3.3, we give bounds for the degree of regularity of a generic zerodimensional system, and in Section 3.3.4, we give some bounds for positive-dimensional systems as well. In Section 3.4, we consider semi-regular systems. We give some equivalent definitions of this property, and we show how asymptotic estimates of the degree of regularity can be adapted from the homogeneous case to the weighted case. Additionally, we prove that Fröberg's conjecture in the case $m=n+1$ is true in the weighted case, as in the homogeous case,


Figure 3.1: Counting monomials in the image of hom $_{W}$ (here $W=(2,3)$ )
provided that the base field is large enough. In Section 3.5, we describe strategies for computing Gröbner bases for weighted homogeneous systems, and we give complexity estimates for these strategies. Finally, in Section 3.6, we show how weighted structures can appear in applications, and we give some benchmarks for each example.

### 3.2. Properties

### 3.2.1. Definitions and properties

As described in the introduction, one can make a $W$-homogeneous polynomial homogeneous using the following proposition.
Proposition 3.1. Let $\left(\mathbb{K}\left[X_{1}, \ldots, X_{n}\right], W\right)$ be a graded polynomial algebra. Then the application

$$
\begin{array}{ccc}
\operatorname{hom}_{W}:\left(\mathbb{K}\left[X_{1}, \ldots, X_{n}\right], W\right) & \longrightarrow\left(\mathbb{K}\left[Y_{1}, \ldots, Y_{n}\right], 1\right) \\
f & \longmapsto f\left(Y_{1}^{w_{1}}, \ldots, Y_{n}^{w_{n}}\right)
\end{array}
$$

is an injective graded morphism, and in particular the image of a weighted homogeneous polynomial is a homogeneous polynomial.
Remark 3.2. This morphism is not surjective. More precisely, the bounds and asymptotics for Sylvester denumerants (Proposition 1.20 and 1.19) apply, and show that on average, there are $\prod w_{i}$ times more monomials in $\mathbb{K}[\mathbf{X}]$ at degree $d$ than at $W$-degree $d$.

The morphism hom ${ }_{W}$ gives an intuition for these asymptotic estimates (Figure 3.1). For example, consider $\mu=Y_{1}^{w_{1}-1} \cdots Y_{n}^{w_{n}-1}$. There is only 1 monomial in the divisors of $\mu$ which also lies in the image of $\operatorname{hom}_{W}$ (namely 1 ); there are $\prod w_{i}$ monomials dividing $\mu$; and all non-trivial multiples of $\mu$ are divisible by more monomials in the image of hom ${ }_{W}$.

The $W$-GRevLex monomial ordering is the pullback of the GRevLex monomial ordering by this morphism:

$$
u<_{W \text {-grevlex }} v \Longleftrightarrow \operatorname{hom}_{W}(u)<_{\text {grevlex }} \operatorname{hom}_{W}(v)
$$

Given a $W$-homogeneous system $F$, one can build the homogeneous system $\operatorname{hom}_{W}(F)$, and then apply classic algorithms (Section 2.4) to that system to compute a GRevLex or Lex Gröbner basis of the ideal generated by $\operatorname{hom}_{W}(F)$.

Definition 3.3. The $W$-degree of regularity of the system $F$ is the highest degree $d_{\mathrm{reg}, W}(F)$ reached in a run of $\mathrm{F}_{5}$ to compute a GRevLex Gröbner basis of $\operatorname{hom}_{W}(F)$. When the graduation is clear in the context, we may call it degree of regularity, and denote it $d_{\text {reg }}$.

Remark 3.4. Unlike what we could observe in the homogeneous case, this definition depends on the order of the variables (we shall give an example in Table 3.1 in Section 3.3.2, and another, with timings, in Table 3.4 in Section 3.6.1).
As before, we only consider the affine varieties associated with the ideals we consider. In particular, for polynomials in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$, the dimension of $V(0)$ is $n$, and a zero-dimensional variety is defined by at least $n$ polynomials if the base field is algebraically closed.

### 3.2.2. Degree and Bézout's bound

The degree of a weighted homogeneous ideal can be defined in the same way as the homogeneous case, using the Hilbert series:

Definition 3.5. Let $I$ be a weighted homogeneous ideal with dimension $d$. Its degree is $\operatorname{deg}(I)=$ $Q(1)$ where $Q=\mathrm{HS}_{A / I}(T) \cdot(1-T)^{d}$.

Remark 3.6. With this definition, the degree has a geometrical interpretation in the zerodimensional case. Then $d=0$, so $\operatorname{deg}(I)=\mathrm{HS}_{A / I}(1)$ is an integer, which corresponds to the dimension of $A / I$ as a $\mathbb{K}$-vector space, that is the cardinality of the staircase of the ideal, or the number of solutions counted with multiplicity.
Remark 3.7. This is the definition used by the software Macaulay 2 for the degree of a positivedimensional variety (function degree (Module) in [Macaulay2]).
Remark 3.8. In the positive-dimensional case, even if the sequence $F$ is in Noether position w.r.t. the variables $X_{1}, \ldots, X_{m}$, the degree of $\langle F\rangle$ is not necessarily the number of solutions of $F_{\text {ext }}=\left(f_{1}, \ldots, f_{m}, X_{m+1}, \ldots, X_{n}\right)=0$. For example, consider $\mathbb{K}[X, Y]$ with the weights $W=(2,1)$, and let $f=X+Y^{2}$. The ideal $\langle f\rangle$ has degree 2 , and indeed, a generic hypersurface of $W$-degree 2 is another parabola with horizontal axis, so that there are 2 intersection points (see Figure 3.2). However, $(f)$ is in Noether position with respect to the variable $Y$, and $\langle f, Y\rangle=\langle X, Y\rangle$ has only 1 zero.
Remark 3.9. It may even happen that positive-dimensional varieties have a rational degree. For example, consider the ideal $\langle X Y\rangle$ in $\mathbb{K}[X, Y]$ graded with respect to $W=(2,3)$. Then

$$
\mathrm{HS}_{A / I}(T)=\frac{\left(1-T^{5}\right)}{\left(1-T^{2}\right)\left(1-T^{3}\right)}=\frac{1+T+T^{2}+T^{3}+T^{4}}{(1-T)(1+T)\left(1+T+T^{2}\right)}
$$

and

$$
\operatorname{deg}(I)=\frac{5}{6}
$$

Intuitively, computing the degree of a positive-dimensional variety can be done by cutting it by generic hypersurfaces. In the total degree case, one can use hyperplanes, that is degree 1 hypersurfaces, so that cutting does not change the degree. However, with $W=(2,3)$, the


Figure 3.2: Common zeroes of $X^{2}+Y$ (in blue) and a generic polynomial of [1, 2]-degree 2 (in red)
first $W$-degree at which one can really form generic hypersurfaces is 6 , and intersecting $X Y$ with a degree 6 hypersurface generically yields 5 solutions: going through hom $_{W}$, a generic intersection will have 30 points, and hom $_{W}$ multiplies this number of solutions by 6 .
The following proposition formalizes the intuitive notion that " going through hom $_{W}$ multiplies the number of solutions by the product of the weights":

Proposition 3.10. Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a system of weights, and $F=\left(f_{1}, \ldots, f_{m}\right)$ a sequence of $W$-homogeneous polynomials, with respective $W$-degrees $d_{1}, \ldots, d_{m}$. Then

$$
\operatorname{deg}(F)=\frac{\operatorname{deg}\left(\operatorname{hom}_{W}(F)\right)}{\prod_{i=1}^{n} w_{i}}
$$

Proof. Let $I=\langle F\rangle$ and $\operatorname{hom}_{W}(I)=\left\langle\operatorname{hom}_{W}(F)\right\rangle$. Write the Hilbert series of $A / I$ (as an algebra graded in $W$-degree) and $A / \operatorname{hom}_{W}(I)$ (as an algebra graded in total degree) using Proposition 1.30:

$$
\begin{aligned}
& \mathrm{HS}_{A / I}(T)=\frac{P(T)}{\prod_{i=1}^{n}\left(1-T^{w}\right)} \\
& \mathrm{HS}_{A / \operatorname{hom}_{W}(I)}=\frac{Q(T)}{(1-T)^{n}} .
\end{aligned}
$$

The proof of Proposition 1.38 actually shows that $P=Q$, so

$$
\mathrm{HS}_{A / I}(T)=\frac{\mathrm{HS}_{A / \operatorname{hom}_{W}(I)}(T)}{\prod_{i=1}^{n}\left(1+T+\cdots+T^{w_{i-1}}\right)}
$$

and evaluating this equation at 1 yields the wanted result.
A classic consequence of this is the weighted version of Bézout's bound:

Theorem 3.11. Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a system of weights, and $F=\left(f_{1}, \ldots, f_{m}\right)$ a regular sequence of $W$-homogeneous polynomials, with respective $W$-degrees $\left(d_{1}, \ldots, d_{m}\right)$. Then

$$
\operatorname{deg}(I)=\operatorname{HS}_{\mathbb{K}[\mathrm{X}] / I}(1)=\frac{\prod_{i=1}^{m} d_{i}}{\prod_{i=1}^{n} w_{i}}
$$

### 3.2.3. Changes of variables and reverse chain-divisible systems of weights

Several properties from the homogeneous case turn out to be no longer true in the weighted case.

For example, many properties of polynomial systems are true up to a linear change of coordinates. Since linear changes of coordinates leave the 1-homogeneous components of the polynomial algebra stable, it is also true of 1-homogeneous polynomial systems. A good example of such a property is Noether's normalization theorem.

However, in a weighted setting, linear transformations on the variables do not necessarily preserve the $W$-degree of the monomials.

If we add some constraints on the system of weights, we may define non-trivial changes of variables. More precisely, we will consider reverse chain-divisible systems of weights, defined as follows.

Definition 3.12. A system of weights, $W$ is reverse chain-divisible if

$$
w_{n}\left|w_{n-1}\right| \cdots \mid w_{1}
$$

In this situation, the weights are coprime if and only if $w_{n}=1$.
Remark 3.13. The name "chain-divisible" can be found in [Alfo5], referring to a notion introduced in [Alf98].

In this setting, many results from the homogeneous case can now be adapted to the weighted homogeneous case. For example, Noether's normalization theorem (Theorem 1.88) is true in a weighted setting if we replace linear changes of coordinates by $W$-homogeneous changes of coordinates with degree $\left(w_{i}\right)$.

Lemma 3.14 (Noether normalization lemma, weighted case). Let $\mathbb{K}$ be an infinite field, $W$ be a reverse chain-divisible system of weights and $f \in R=\mathbb{K}\left[X_{1}, \ldots, X_{r}\right]$ be a non-constant polynomial, $W$-homogeneous with $W$-degree $d$. Then there are elements $X_{1}^{\prime}, \ldots, X_{r-1}^{\prime} \in R$ such that $R$ is a finitely generated module over $\mathbb{K}\left[X_{1}^{\prime}, \ldots, X_{r-1}^{\prime}, f\right]$. Furthermore, if the field has characteristic 0 or large enough, there exists a dense Zariski-open subset $U \subset \mathbb{K}^{r-1}$ such that for all $\left(a_{1}, \ldots, a_{r-1}\right) \in U$, one can choose $X_{i}^{\prime}=X_{i}-a_{i} X_{r}^{w_{i} / w_{r}}$.

Proof. We follow the proof of Lemma 1.89. For any $1 \leq i \leq r-1$, let $a_{i} \in \mathbb{K}$, and let $X_{i}^{\prime}=$ $X_{i}-a_{i} X_{r}^{w_{i} / w_{r}}$. We need to show that for generic $a_{i}$, under this change of variables, $f$ is monic in $X_{r}$. Under this change of variables, collecting the coefficients of $f$ in $X_{r}$ yields:

$$
\begin{aligned}
f\left(X_{1}, \ldots, X_{r}\right) & =f\left(X_{1}^{\prime}+a_{1} X_{r}^{w_{1} / w_{r}}, X_{2}^{\prime}+a_{2} X_{r}^{w_{2} / w_{r}}, \ldots, X_{r-1}+a_{r-1} X_{r}^{w_{r-1} / w_{r}}\right) \\
& =f\left(a_{1}, \ldots, a_{r-1}, 1\right) X_{r}^{d}+\cdots
\end{aligned}
$$

So the set of all $a_{i}$ 's such that $f$ is monic in $X_{r}$ is exactly the set of all $a_{i}$ 's such that

$$
f\left(a_{1}, \ldots, a_{r-1}, 1\right) \neq 0,
$$

and since $f$ is $W$-homogeneous non-constant, this is a non-empty open subset of $\mathbb{K}^{r-1}$.
Apart from the lemma, the proof of Noether's normalization theorem does not depend on the graduation, and as a consequence:

Theorem 3.15 (Noether's normalization theorem, weighted case). Let $W$ be a reverse chaindivisible system of weights, and let $F$ be a $W$-homogeneous zero-dimensional regular sequence in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$. Then, for a generic choice of $W$-homogeneous polynomials $P_{i}$ with $W$-degree $w_{i}$, the change of variable

$$
X_{i}=X_{i}^{\prime}+P_{i}\left(X_{i+1}, \ldots, X_{n}\right),
$$

is such that $F\left(X_{1}\left(\mathrm{X}^{\prime}\right), \ldots, X_{n}\left(\mathrm{X}^{\prime}\right)\right)$ is in simultaneous Noether position with respect to the order $X_{1}^{\prime}>X_{2}^{\prime}>\cdots>X_{n}^{\prime}$.

A central property of reverse chain-divisible weights is the following proposition. In the homogeneous case, if $d_{1} \leq d_{2}$ are two non-negative integers, then any monomial with degree $d_{2}$ is divisible by a monomial with degree $d_{1}$. When the system of weights is reverse chain-divisible, the following proposition states a similar result for the weighted case.

Proposition 3.16. Assume that $W=\left(w_{1}, \ldots, w_{n}\right)$ is a system of weights, such that $w_{1} \geq w_{2} \geq$ $\cdots \geq w_{n}$. The following statements are equivalent:

1. The system of weights $W$ is reverse chain-divisible;
2. Let $d_{1} \leq d_{2}$ positive integers, $i \in\{1, \ldots, n\}$, and $m_{2}$ a monomial of $W$-degree $d_{2}$. Assume that $w_{i}$ divides $d_{1}$, and that $m_{2}$ is not divisible by any of the variables $X_{1}, \ldots, X_{i-1}$. Then there exists a monomial $m_{1}$ with $W$-degree $d_{1}$ such that $m_{1} \mid m_{2}$.

Proof. $(1 \Longrightarrow 2)$. Fix $d_{1}$. We shall prove by induction over $d_{2}$ that for any monomial $m_{2}$ with $W$-degree $d_{2}$ satisfying the hypotheses of (2), there exists a monomial $m_{1}$ with $W$-degree $d_{1}$ dividing $m_{2}$. The case $d_{2}=d_{1}$ is immediate, since we can use $m_{1}:=m_{2}$.
Assume that $d_{2}>d_{1}$, and let $m_{2}$ be a monomial of $W$-degree $d_{2}$. Let $j$ be the greatest index of a variable dividing $m_{2}$, write $m_{2}=X_{j}^{\alpha} m_{2}^{\prime}$, where $m_{2}^{\prime}$ is a monomial in $\mathbb{K}\left[X_{i}, \ldots, X_{j-1}\right]$, with $W$-degree $d_{2}^{\prime}=d_{2}-w_{j} \alpha$. If $d_{2}^{\prime} \geq d_{1}$, the result follows by induction. If $i=j$, then $m_{2}=X_{i}^{\alpha}$, and $m_{1}:=X_{i}^{d_{1} / w_{i}}$ has $W$-degree $d_{1}$ and divides $m_{2}$. So we can assume that $d_{2}^{\prime}<d_{1}$ and that $i<j$.

Since $W$ is a reverse chain-divisible system of weights, $w_{j-1}$ divides $w_{k}$ for any $k$ in $\{1, \ldots, j-$ $1\}$. Hence, since $m_{2} \in \mathbb{K}\left[X_{i}, \ldots, X_{j}\right]$ and $m_{2}^{\prime} \in \mathbb{K}\left[X_{i}, \ldots, X_{j-1}\right], d_{2} \equiv 0\left[w_{j}\right]$ and $d_{2}^{\prime} \equiv 0\left[w_{j-1}\right]$. By hypothesis, $d_{1}$ is divisible by $w_{i}$, and in particular it is divisible by $w_{j-1}$. All in all, this shows that $d_{1}-d_{2}^{\prime}$ is divisible by $w_{j-1}$, and so it is divisible by $w_{j}$. Let

$$
m_{1}:=m_{2}^{\prime} \cdot X_{j}^{\left(d_{1}-d_{2}^{\prime}\right) / w_{j}} .
$$

Then the monomial $m_{1}$ has $W$-degree $d_{1}$ and divides $m_{2}$.
(2 $\Longrightarrow$ 1). Assume that $W$ is a system of weights which is not reverse chain-divisible. We shall find integers $d_{1} \leq d_{2}$ and a monomial $m_{2}$ with $W$-degree $d_{2}$ which is not divisible by any monomial of $W$-degree $d_{1}$.

Since $W$ is not reverse chain-divisible, there exists $i$ such that $w_{i+1}$ does not divide $w_{i}$. In particular, $\operatorname{gcd}\left(w_{i}, w_{i+1}\right)<w_{i}$ and $\operatorname{gcd}\left(w_{i}, w_{i+1}\right)<w_{i+1}$. Without loss of generality, we may consider only the variables $X_{i}, X_{i+1}$. Let $d_{1}=w_{i} w_{i+1}, d_{2}=d_{1}+\operatorname{gcd}\left(w_{i}, w_{i+1}\right)$. By Paoli's lemma (see for example [Luc91, chap. 264] or the discussion after [NZM91, th. 5.1]), there exists exactly

$$
\left\lfloor\frac{d_{2}}{w_{i} w_{i+1}}\right\rfloor=\left\lfloor 1+\frac{\operatorname{gcd}\left(w_{i}, w_{i+1}\right)}{w_{i} w_{i+1}}\right\rfloor=1
$$

couple of non-negative integers $a, b$ such that $a w_{i}+b w_{i+1}=d_{2}$. Let $m_{2}$ be the monomial $X_{i}^{a} X_{i+1}^{b}$. The $W$-degree $d_{1}$ is divisible by $w_{i}$, and $m_{2}$ is not divisible by $X_{1}, \ldots, X_{i-1}$. The maximal divisors of $m_{2}$ are

$$
\begin{aligned}
& \frac{m_{2}}{X_{i}}=X_{i}^{a-1} X_{i+1}^{b} \text { with } W \text {-degree } d_{2}-w_{i}=d_{1}+\operatorname{gcd}\left(w_{i}, w_{i+1}\right)-w_{i}<d_{1} \\
& \frac{m_{2}}{X_{i+1}}=X_{i}^{a} X_{i+1}^{b-1} \text { with } W \text {-degree } d_{2}-w_{i+1}=d_{1}+\operatorname{gcd}\left(w_{i}, w_{i+1}\right)-w_{i+1}<d_{1}
\end{aligned}
$$

As a consequence, $m_{2}$ is not divisible by any monomial of $W$-degree $d_{1}$.
This proposition essentially states that the staircase of a $W$-homogeneous ideal is reasonably shaped when $W$ is a reverse chain-divisible system of weights. For example, let $W$ be a reverse chain-divisible system of weights, and let $I$ be the ideal generated by all monomials of $W$-degree $w_{1}$ (that is, the least common multiple of the weights). Then the proposition proves that $I$ contains all monomials of $W$-degree greater than $w_{1}$.

If on the other hand the system of weights is not reverse chain-divisible, this property needs not hold. For example, consider the algebra $\mathbb{K}\left[X_{1}, X_{2}, X_{3}\right]$ graded w.r.t. the system of weights $W=(3,2,1)$, the least common multiple of the weights being 6 , and let $I$ be the ideal generated by all monomials of $W$-degree 6 . Consider the monomial $X_{1} X_{2}^{2}$ : it has $W$-degree 7 , yet it is not divisible by any monomial with $W$-degree 6 , and so it does not belong to the ideal $I$.

### 3.2.4. Genericity of regularity properties and $W$-compatibility

Proposition 3.17. Let $m \leq n$ be two integers, $W=\left(w_{1}, \ldots, w_{n}\right)$ a system of weights, and $D=$ $\left(d_{1}, \ldots, d_{m}\right)$ a system of $W$-degrees. Then

- the set of regular sequences,
- the set of sequences in Noether position with respect to the variables $X_{1}, \ldots, X_{m}$, and
- the set of sequences in simultaneous Noether position w.r.t. the order $X_{1}>\cdots>X_{m}$
are Zariski-open subsets of the affine space of $W$-homogeneous polynomials with $W$-degree $D$.

This states that the set of regular sequences, sequences in Noether position and sequences in simultaneous Noether position are Zariski-dense subsets if and only if they are not empty. Unfortunately, depending on the weights and the weighted degrees, there may exist no regular sequence, and thus no sequences in (simultaneous) Noether position either. For example, let $W=(2,5)$ and $D=(4,8)$, the only $W$-homogeneous sequence with $W$-degree $D$ in $\mathbb{K}[X, Y]$ is (up to scalar multiplication) $\left(X^{2}, X^{4}\right)$, and it is not regular. However, this is only the case for very specific systems of $W$-degrees for which there does not exist enough monomials to build non-trivial sequences.

Definition 3.18. Let $m \leq n$ be two integers, $W=\left(w_{1}, \ldots, w_{n}\right)$ a system of weights, and $D=\left(d_{1}, \ldots, d_{m}\right)$ a system of $W$-degrees. We say that $D$ is $W$-compatible if there exists a regular $W$-homogeneous sequence in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$ with $W$-degree $D$. We say that $D$ is strongly $W$ compatible if for any $1 \leq i \leq m, d_{i}$ is divisible by $w_{i}$.

Using these definitions, we can identify cases where the properties of being regular, in Noether position or in simultaneous Noether position are generic.

Proposition 3.19. Let $m \leq n$ be two integers, $W=\left(w_{1}, \ldots, w_{n}\right)$ a system of weights, and $D=$ $\left(d_{1}, \ldots, d_{m}\right)$ a system of $W$-degrees. For any $1 \leq i \leq m$, write $W_{i}:=\left(w_{1}, \ldots, w_{i}\right)$ and $D_{i}:=$ $\left(d_{1}, \ldots, d_{i}\right)$. Write $A_{W, D}$ the affine space of $W$-homogeneous sequences of $W$-degree $D$. Then the following statements are true:

1. if $D$ is $W$-compatible, then regular sequences form a Zariski-dense subset of $A_{W, D}$;
2. if $D$ is $W_{m}$-compatible, then sequences in Noether position with respect to the variables $X_{1}, \ldots, X_{m}$ form a Zariski-dense subset of $A_{W, D}$;
3. if $D$ is strongly $W$-compatible, then $D$ is $W$-compatible, $W_{m}$-compatible, and for any $i, D_{i}$ is $W_{i}$-compatible;
4. if $m=n, D$ is $W$-compatible and $W$ is reverse chain-divisible, then, up to some reordering of the degrees, $D$ is strongly $W$-compatible.

Proof. The proofs of statements 1 and 2 follow the same technique: by Theorem 3.17, we know that the sets we consider are Zariski-open in $A_{W, D}$. So in order to prove the density, we only need to prove that they are non empty. Statement 1 is exactly the definition of the $W$-compatibility.

For statement 2, by $W_{m}$-compatibility, we know that there exists a $W$-homogeneous sequence $F=\left(f_{1}, \ldots, f_{m}\right)$ with $W$-degree $D$ in $\mathbb{K}\left[X_{1}, \ldots, X_{m}\right]$, which is regular. As a consequence, the sequence $\left(f_{1}, \ldots, f_{m}, X_{m+1}, \ldots, X_{n}\right)$ is regular, and from the characterization 4 of Noether position (prop. 1.44), this means that $F$ is in Noether position with respect to the variables $X_{1}, \ldots, X_{m}$.

In order to prove statement 3, let $1 \leq i \leq m$, we need to exhibit a regular sequences of length $i$ in $\mathbb{K}\left[X_{1}, \ldots, X_{i}\right]$. We may choose $F_{i}:=\left(X_{1}^{d_{1} / w_{1}}, \ldots, X_{i}^{d_{i} / w_{i}}\right)$, it is regular and each polynomial lies in $\mathbb{K}\left[X_{1}, \ldots, X_{i}\right]$.

Finally, statement 4 is a consequence of Theorem 3.15. Let $W$ be a reverse chain-divisible system of weights, and $D$ a $W$-compatible system of $W$-degrees. Up to reordering, we can assume that the polynomials are ordered so that $d_{1} \geq d_{2} \geq \cdots \geq d_{n}$; this does not cancel the
$W$-compatibility. Let $F=\left(f_{1}, \ldots, f_{n}\right)$ be a regular sequence, $W$-homogeneous with $W$-degree $D$. By Theorem 3.15, there exist polynomials $P_{i}\left(X_{i+1}, \ldots, X_{n}\right)$ which are $W$-homogeneous with $W$-degree $w_{i}$, and such that $F$, under the change of variables $X_{i}=X_{i}^{\prime}+P_{i}\left(X_{i+1}, \ldots, X_{n}\right)$, is in simultaneous Noether position with respect to the order $X_{1}^{\prime}>X_{2}^{\prime}>\cdots>X_{n}^{\prime}$.
From the characterization 4 of Noether position, this means in particular that for any $i \in$ $\{1, \ldots, n\}, f_{i}\left(X_{1}\left(X_{1}^{\prime}, \ldots, X_{i}^{\prime}\right), \ldots, X_{n}\left(X_{1}^{\prime}, \ldots, X_{i}^{\prime}\right)\right)$ belongs to a regular sequence, and thus is not zero. By definition of reverse chain-divisible weights, its $W$-degree $d_{i}$ is a sum of multiples of $w_{i}$, and so it is itself a multiple of $w_{i}$.

Remark 3.20. The statement 4 is a converse of 3 in the reverse chain-divisible case. In the non-reverse chain-divisible case, that converse is false: let $W=(3,2), D=(6,5)$ and consider $F=\left(X^{2}+Y^{3}, X Y\right)$ in $\mathbb{K}[X, Y]$. The sequence $F$ is in simultaneous Noether position w.r.t. the order $X>Y$, yet 5 is neither divisible by 3 nor by 2 .
The weaker converse that if $D$ is $W$-compatible, then $D$ is $W_{m}$-compatible is also false: with the same weights and algebra, let $D=(5)$, the only polynomial with $W$-degree 5 is (up to scalar multiplication) $f=X Y$. It is non-zero, so that $(f)$ is a regular sequence, but $(f, Y)$ is not regular, hence $(f)$ is not in Noether position w.r.t $X$.
Remark 3.21. These examples lead to the following attempt at writing a general characterization of $W$-compatibility.

Let $n$ be a positive integer, $W=\left(w_{1}, \ldots, w_{n}\right)$ a system of weights, and $D=\left(d_{1}, \ldots, d_{n}\right)$ a system of $W$-degrees. Further assume that

- for all $i \in\{1, \ldots, n\}, \mathbb{K}[\mathrm{X}]_{d_{i}} \neq 0$
- the formal series

$$
S_{D, W}(T)=\frac{\prod_{i=1}^{n}\left(1-T^{d_{i}}\right)}{\prod_{i=1}^{n}\left(1-T^{w_{i}}\right)}
$$

is a polynomial.
Is $D$ necessarily $W$-compatible?
The answer is no: take the system of weights $W=(3,5,11)$, and the system of $W$-degrees $D=(165,19,19)$. Note that 165 is the product of the weights, and 19 the sum of the weights. The series

$$
S_{D, W}(T)=\frac{\left(1-T^{165}\right) \cdot\left(1-T^{19}\right) \cdot\left(1-T^{19}\right)}{\left(1-T^{3}\right) \cdot\left(1-T^{5}\right) \cdot\left(1-T^{11}\right)}=1+T^{3}+\cdots+T^{184}
$$

is a polynomial. But at $W$-degree 19 , there are only 2 monomials, namely $X_{1} X_{2} X_{3}$ and $X_{1}^{3} X_{2}^{2}$, and they are not coprime, so we cannot form a regular sequence of $W$-degrees $(165,19,19)$.
Remark 3.22. Bézout's bound gives another necessary condition for $D$ to be $W$-compatible in the zero-dimensional case. Indeed, if $n=m$ and $D$ is $W$-compatible, let $I$ be an ideal generated by a regular sequence of $W$-homogeneous polynomials with respective $W$-degree $D$. Then $\mathbb{K}[X] / I$ is a finite-dimensional $\mathbb{K}$-vector space, and its dimension is given by

$$
\operatorname{dim}_{\mathbb{K}}(\mathbb{K}[\mathbf{X}] / I)=\operatorname{deg}(I)=\frac{\prod_{i=1}^{n} d_{i}}{\prod_{i=1}^{n} w_{i}}
$$

Since this dimension is necessarily an integer, $\prod_{i=1}^{n} w_{i}$ divides $\prod_{i=1}^{n} d_{i}$.
In the homogeneous case, affine systems were studied by homogenization. In order to be able to study affine systems with a system of weights, we introduce the following definition:

Definition 3.23. Let $F=\left(f_{1}, \ldots, f_{m}\right) \subset \mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$ be a sequence of polynomials, and let $W$ be a system of weights. We say that $F$ has a $W$-homogeneous structure if the highest $W$-degree components of $f_{1}, \ldots, f_{n}$ form a regular sequence.

### 3.2.5. Characterization of the Hilbert series

Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a reverse chain-divisible system of weights such that $w_{n}=1$, and let $D=\left(d_{1}, \ldots, d_{n}\right)$ be a system of $W$-degrees, such that for any $i \in\{1, \ldots, n\}, d_{i}$ is divisible by all of the $w_{j}$ 's. Let $R=\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$ be a polynomial algebra graded with respect to $W$.

We use the notations below, following [Morg6]:

- $\delta_{j}=\sum_{i=1}^{j}\left(d_{i}-w_{i}\right) ;$
- $\delta=\delta_{n}, \delta^{*}=\delta_{n-1}$;
- $\sigma=\min \left(\delta^{*},\left\lfloor\frac{\delta}{2}\right\rfloor\right), \sigma^{*}=\min \left(\delta_{n-2},\left\lfloor\frac{\delta^{*}}{2}\right\rfloor\right)$;
- $\mu=\delta-2 \sigma, \mu^{*}=\delta^{*}-2 \sigma^{*}$.

Given a formal series $S(T)=\sum_{d=0}^{\infty} a_{d} T^{d}$, we also define

$$
\begin{aligned}
\Delta S(T) & =\sum_{d=0}^{\infty}\left(a_{d}-a_{d-1}\right) T^{d} \quad\left(\text { with the convention } a_{-1}=0\right) \\
& =(1-T) \cdot S(T)
\end{aligned}
$$

and

$$
\int S=\sum_{d=0}^{\infty}\left(a_{0}+\cdots+a_{d}\right) T^{d}=\frac{S(T)}{1-T} .
$$

Lemma 3.24. Under the above notations and assumptions, the following properties hold.

$$
\begin{align*}
& \left\{\begin{array}{l}
\delta^{*}>\left\lfloor\frac{\delta}{2}\right\rfloor \quad \Longleftrightarrow \\
\left\{\begin{array}{l}
\delta_{n}-\delta^{*} \leq 0 \\
\delta^{*}=\left\lfloor\frac{\delta}{2}\right\rfloor
\end{array} \quad \Longleftrightarrow \quad 1 \leq d_{n}-\delta^{*} \leq 2\right. \\
\delta^{*}<\left\lfloor\frac{\delta}{2}\right\rfloor \quad \Longleftrightarrow \quad 3 \leq d_{n}-\delta^{*}
\end{array}\right.  \tag{3.5}\\
& \sigma=\left\lfloor\frac{\delta}{2}\right\rfloor \Longrightarrow \mu=\delta[2] \in\{0,1\}  \tag{3.6}\\
& 0 \leq \mu<d_{n}  \tag{3.7}\\
& d_{n-1} \leq d_{n} \Longrightarrow \sigma^{*}+\mu^{*} \leq \sigma \tag{3.8}
\end{align*}
$$

Proof. The proof of statements (3.5) and (3.7) can be found in [Mor96, Lemma 2.1]. This proof depends only on the value of $w_{n}$, and since we assume it to be 1 , it is also valid in our setting. It also proves (3.6) as a side-result.

For the statement (3.8), we proceed by case disjunction on the values of $\sigma$.

- If $\sigma=\delta^{*}$ :

$$
\sigma^{*}+\mu^{*}=\delta^{*}-\sigma^{*} \leq \delta^{*}=\sigma
$$

- If $\sigma=\lfloor\delta / 2\rfloor$, then $\sigma=\left\lfloor\left(\delta^{*}+d_{n}-1\right) / 2\right\rfloor$ which implies $2 \sigma=\delta^{*}+d_{n}-1-\mu$ and $\mu=\delta[2] \in\{0,1\}$ (from statement (3.6)). Now consider the possible values of $\sigma^{*}$ :
- if $\sigma^{*}=\left\lfloor\delta^{*} / 2\right\rfloor$, then $\mu^{*}=\delta^{*}[2]$, and thus $2 \sigma=2 \sigma^{*}+\mu^{*}+d_{n}-1-\mu$. It implies that $d_{n}-1-\mu+\mu^{*}$ is even. We shall prove that it is greater than or equal to 0 .
From statement (3.7), $d_{n}-1-\mu \geq 0$, so if $\mu^{*}=0$, we are done. If $\mu^{*}=1$, by parity $d_{n}-1-\mu$ is odd, and thus $d_{n}-1-\mu \geq 1=\mu^{*}$.
It implies that:

$$
2 \sigma=2 \sigma^{*}+\mu^{*}+d_{n}-1-\mu \geq 2 \sigma^{*}+2 \mu^{*}
$$

- otherwise, $\sigma^{*}=\delta^{* *}$, and in that case

$$
\sigma^{*}+\mu^{*}=\delta^{*}-\sigma^{*}=\delta^{*}-\delta^{* *}=d_{n-1}-w_{n-1}
$$

which implies that:

$$
d_{n}-1 \geq \sigma^{*}+\mu^{*}\left(\text { since } w_{n-1} \geq w_{n} \text { and } d_{n-1} \leq d_{n}\right)
$$

and

$$
\delta^{*}=\delta^{* *}+d_{n-1}-w_{n-1} \geq \sigma^{*}+\mu^{*}
$$

So we have:

$$
\begin{aligned}
2 \sigma & =\delta^{*}+d_{n}-1-\mu \\
& \geq \sigma^{*}+\mu^{*}+\sigma^{*}+\mu^{*}-\mu
\end{aligned}
$$

Recall that $\mu \in\{0,1\}$, so by parity, $2 \sigma \geq 2 \sigma^{*}+2 \mu^{*}$, hence $\sigma \geq \sigma^{*}+\mu^{*}$.
The following theorem describes the shape of the Hilbert series of a zero-dimensional complete intersection. It states that it is a self-reciprocal (or palindromic) polynomial, that is a polynomial with symmetrical coefficients, and that these coefficients increase at small degrees, then station, then decrease again. Furthermore, between every strict increase, they reach a step, which has width $w_{n-1}$. For an example, see figure $3 \cdot 3$, where the width of the steps is 3 , and the width of the central plateau is 5 .

This is a generalization of a known result in the homogeneous case (Proposition 1.40) which has been proved for example in [Mor96, prop. 2.2] (we will follow that proof for the weighted case). The result is simpler in the homogeneous case: there is no such step in the growth of the coefficients, and they are strictly increasing, then stationary, then strictly decreasing.


Figure 3.3: Shape of the Hilbert series of a $W$-homogeneous complete intersection for $W=(3,3,1)$ and $D=(9,6,3)$

Theorem 3.25. Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a reverse chain-divisible system of weights, and $D=$ $\left(d_{1}, \ldots, d_{n}\right)$ a system of degrees such that for any $i \in\{1, \ldots, n\}, d_{i}$ is divisible by $w_{1}$. Consider the formal series

$$
S_{W, D}(T)=\frac{\prod_{i=1}^{n}\left(1-T^{d_{i}}\right)}{\prod_{i=1}^{n}\left(1-T^{w_{i}}\right)}=\sum_{d=0}^{\delta} a_{d} T^{d}
$$

The series $S_{W, D}$ is a self-reciprocal polynomial in $T$ (i.e. for any $d \leq \delta, a_{d}=a_{\delta-d}$ ) and its coefficients satisfy the inequalities:

$$
\begin{aligned}
\forall d \in\{0, \ldots, \sigma-1\}, & a_{d} \leq a_{d+1} \\
\forall d \in\{\sigma, \ldots, \sigma+\mu-1\}, & a_{d}=a_{d+1} \\
\forall d \in\{\sigma+\mu, \ldots, \delta\}, & a_{d} \geq a_{d+1}
\end{aligned}
$$

Furthermore, if $d<\sigma$ (resp. $d>\sigma+\mu$ ), the coefficients increase (resp. decrease) with steps, and these steps have width $w_{n-1}$ :

$$
\forall d \in\{0, \ldots, \sigma-1\}, a_{d}-a_{d-1} \begin{cases}>0 & \text { if } w_{n-1} \text { divides } d \\ =0 & \text { otherwise } .\end{cases}
$$

Proof. We adapt the proof from [Mor96, Prop. 2.2] for the homogeneous case to the weighted case. Up to permutation of the $d_{i}$ 's, we can assume that for any $i, d_{i} \geq d_{i-1}$. We proceed by induction on $n$. The result for the case $n=1$ is a consequence of the homogeneous case, since $w_{n}=1$.

Let $n>1$. Let $\bar{W}^{*}=\left(w_{1} / w_{n-1}, \ldots, w_{n-1} / w_{n-1}\right)$ and $\bar{D}^{*}=\left(d_{1} / w_{n-1}, \ldots, d_{n-1} / w_{n-1}\right)$, and consider the series

$$
\bar{S}^{*}:=S_{W^{*}, D^{*}}=\frac{\prod_{i=1}^{n-1}\left(1-T^{d_{i} / w_{n-1}}\right)}{\prod_{i=1}^{n-1}\left(1-T^{w_{i} / w_{n-1}}\right)}=\sum_{d=0}^{\delta} \bar{a}_{d}^{*} T^{d}
$$

Write $\bar{S}^{*}$ the Hilbert series of $\bar{F}^{*}$, with generic coefficient $\bar{a}_{d}^{*}$. The system $\bar{F}^{*}$ satisfies the hypotheses of the theorem (from the caracterization 3 of the Noether position 1.44), so by induction hypothesis, $\bar{S}^{*}$ has the shape predicted by the theorem.

The Hilbert series $S$ can be computed from $\bar{S}^{*}$ with

$$
S(T)=\frac{1-T^{d_{n}}}{1-T} \bar{S}^{*}\left(T^{w_{n-1}}\right)=\left(1-T^{d_{n}}\right) \cdot \int \bar{S}^{*}\left(T^{w_{n-1}}\right)
$$

and so for any $d$, we have:

$$
\begin{aligned}
& a_{d}=a_{d-d_{n}+1}^{*}+\cdots+a_{d}^{*} \\
& a_{d}^{\prime}:=a_{d}-a_{d-1}=a_{d}^{*}-a_{d-d_{n}}^{*}
\end{aligned}
$$

where

$$
a_{d}^{*}= \begin{cases}\bar{a}_{\bar{d}}^{*} & \text { if } d=\bar{d} w_{n-1} \\ 0 & \text { otherwise }\end{cases}
$$

This proves that the polynomial is self-reciprocal. Indeed:

$$
\begin{aligned}
a_{\delta-d} & =a_{\delta-d-d_{n}+1}^{*}+\cdots+a_{\delta-d}^{*} \\
& =a_{d-d_{n}+1}^{*}+\cdots+a_{d}^{*} \text { since, by induction hypothesis, } \bar{S}^{*} \text { is self-reciprocal } \\
& =a_{d}
\end{aligned}
$$

To prove the properties regarding the sign of $a_{d}^{\prime}=a_{d}-a_{d-1}$, we shall consider two cases, according to the value of $d_{n}$.

- If $d_{n} \geq \delta^{*}+1$, then from statement (3.5) in Lemma 3.24, and the definition of $\sigma$ and $\mu$, $\sigma=\delta^{*}$ and $\sigma+\mu=d_{n}-1$. Let $0 \leq d \leq \sigma$, then $d \leq \delta^{*}<d_{n}$, so:

$$
a_{d}^{\prime}=a_{d}^{*}= \begin{cases}\bar{a}_{d / w_{n-1}}^{*}>0 & \text { if } w_{n-1} \text { divides } d ; \\ 0 & \text { otherwise }\end{cases}
$$

Let $d \in\{\sigma+1, \ldots, \sigma+\mu\}$, that implies that $\delta^{*}<d \leq d_{n}-1$, so:

$$
a_{d}^{\prime}=a_{d}^{*}=0\left(\text { since } \delta^{*} \text { is the degree of } S^{*}\right) .
$$

- If $d_{n} \leq \delta^{*}$, then from statement (3.5) again, $\sigma=\lfloor\delta / 2\rfloor$ and $\mu=\delta$ [2]. Let $d \leq \sigma$. We want to prove that $a_{d}-a_{d-1}$ is greater or equal to zero, depending on whether $d$ is divisible by $w_{n-1}$. We shall consider two ranges of values for $d$ :
- if $d \leq \sigma^{*}+\mu^{*}$, then $d-d_{n} \leq \sigma^{*}+\mu^{*}-d_{n}<\sigma^{*}\left(\right.$ since $\left.\mu^{*}<d_{n}\right)$. Recall that $a_{d}^{\prime}=a_{d}^{*}-a_{d-d_{n}}^{*}$. By hypothesis, $d_{n}$ is divisible by $w_{n-1}$, and so, either both $d$ and $d-d_{n}$ are divisible by $w_{n-1}$, or both are not. Thus,
$a_{d}^{\prime} \begin{cases}>0 & \text { if both } d \text { and } d-d_{n} \text { are divisible by } w_{n-1} \\ =0 & \text { if neither } d \text { nor } d-d_{n} \text { is divisible by } w_{n-1} ;\end{cases}$
- if $\sigma^{*}+\mu^{*}<d \leq \sigma$, then $2 d \leq 2 \sigma \leq \delta$; by definition, $\delta=\delta^{*}+d_{n}-1$, so $d-d_{n}<\delta^{*}-d$; furthermore, $\delta^{*}-d<\delta^{*}-\left(\sigma^{*}+\mu^{*}\right)=\sigma^{*}$, so in the end:

$$
d-d_{n}<\delta^{*}-d<\sigma^{*}
$$

Since, by construction, $\delta^{*}$ is divisible by $w_{n-1}$, the same reasoning as before yields that

$$
a_{d}^{\prime}=a_{d}^{*}-a_{d-d_{n}}^{*}=a_{\delta^{*}-d}^{*}-a_{d-d_{n}}^{*}
$$

and

$$
a_{d}^{\prime} \begin{cases}>0 & \text { if both } \delta^{*}-d \text { and } d-d_{n} \text { are divisible by } w_{n-1} \\ =0 & \text { if neither } \delta^{*}-d \text { nor } d-d_{n} \text { is divisible by } w_{n-1}\end{cases}
$$

Still assuming that $d_{n} \leq \delta^{*}$, let now $d \in\{\sigma+1, \ldots, \sigma+\mu\}$. We want to prove that $a_{d}-a_{d-1}=0$. If $\mu=0$ there is nothing to prove, so assume that $\mu=1$ and $d=\sigma+1$. But then $\sigma+1-d_{n}=\delta-\sigma-d_{n}=\delta^{*}-\sigma$, and so by symmetry, $a_{d}^{\prime}=a_{\sigma+1}^{*}-a_{\sigma+1-d_{n}}^{*}=0$.

Remark 3.26. The hypothesis that the weights are reverse chain-divisible is necessary. As a counter-example, let $W=(3,2,2)$ and $D=(6,6,6)$. Then the Hilbert series of a complete intersection of $W$-degree $D$ is illustrated in Figure 3.4. It is self-reciprocal, but the coefficients do not vary as predicted by Theorem 3.25.

The hypothesis that each of the $W$-degrees should be divisible by $w_{1}$ is also necessary. As a counter-example, let $W=(4,2,1)$ and $D=(8,8,2)$. Then the Hilbert series of a complete intersection of $W$-degree $D$ is illustrated in Figure 3.5: the width of the steps is greater than $w_{n-1}$. Furthermore, following the proof, the parameters for this series should be defined by $\sigma=\lfloor\delta / 2\rfloor$ and $\mu=\delta[2]$, where $\delta=11$, so that $\sigma=5$ and $\mu=1$. However, we cannot reorder the degrees such that $d_{3} \geq d_{2} \geq d_{1}$, and we cannot deduce from statement (3.8) in Lemma 3.24 that $\sigma^{*}+\mu^{*} \leq \sigma$ : indeed, we have $\sigma=4$ but $\sigma^{*}+\mu^{*}=6$.

However, the fact that the Hilbert series is self-reciprocal for complete intersections is true even for general system of weights, and is a consequence of the Gorenstein property of complete intersections (see [Eis95, Chap. 21]; this property is also central to the proof of Theorem 3.42).

## 3.3. $W$-Degree of regularity of regular sequences

### 3.3.1. Macaulay's bound in dimension zero

Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a system of weights, and $F=\left(f_{1}, \ldots, f_{m}\right)$ a regular sequence of $W$ homogeneous polynomials, with respective $W$-degrees $d_{1}, \ldots, d_{m}$. Further assume that the set of solutions is zero-dimensional, that is $m=n$. We denote by $I$ the weighted homogeneous ideal generated by $F$.

One can read on the Hilbert series a weighted version of Macaulay's bound:


Figure 3.4: Hilbert series of a weighted homogeneous complete intersection with $W=(3,2,2)$ and $D=(6,6,6)-$ The weights are not reverse chain-divisible.


Figure 3.5: Hilbert series of a weighted homogeneous complete intersection with $W=(4,2,1)$ and $D=(8,8,2)-$ The $W$-degrees are not all divisible by $w_{1}$.

Theorem 3.27 (Macaulay's bound, weighted case). With the same notations and hypotheses, the $W$-degree of regularity of $F$ is bounded by

$$
\begin{equation*}
d_{\mathrm{reg}, W}(F) \leq i_{r e g}(I)+\max \left\{w_{i}\right\}=\sum_{i=1}^{n} d_{i}-w_{i}+\max \left\{w_{i}\right\} . \tag{3.9}
\end{equation*}
$$

Proof. As in the homogeneous case, since $I$ is zero-dimensional, the Hilbert series is a polynomial with degree $i_{\text {reg }}$. This means that all monomials of $W$-degree greater than $i_{\text {reg }}$ are in the ideal, and as such, that the leading terms of the $W$-GRevLex Gröbner basis of $F$ need to divide all the monomials of $W$-degree greater than $i_{\text {reg }}$.

Let $f \in \mathbb{K}[\mathbf{X}], W$-homogeneous with $W$-degree at least $i_{\text {reg }}+\max \left\{w_{i}\right\}+1$. If it belongs to a reduced Gröbner basis of $I$, then none of the proper divisors of his leading term are in $\mathrm{LT}(I)$, but they necessarily have $W$-degree at least $i_{\text {reg }}+1$, so we reached a contradiction. Hence all polynomials in a reduced $W$-GRevLex Gröbner basis of $I$ have $W$-degree at most $i_{\text {reg }}+\max \left\{w_{i}\right\}$.

In the homogeneous case, this bound is sharp, and is actually reached for generic system, but it is not true in the weighted case. In particular, it appears that the $W$-degree of regularity depends on the order set on the variables (see Table 3.1 for some examples).

### 3.3.2. Sharper bound on the degree of regularity

The following theorem is an improvement over the previous bound, under the additional assumption that the system is in simultaneous Noether position. Recall that this property is generic, and that for reverse chain-divisible systems of weights, it is always true for regular sequences, up to a weighted homogeneous change of coordinates.

Theorem 3.28. Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a (not necessarily reverse chain-divisible) system of weights and $D=\left(d_{1}, \ldots, d_{n}\right)$ be a strongly $W$-compatible system of $W$-degrees. Further assume that for any $j \in\{2, \ldots, n\}, d_{j} \geq w_{j-1}$. Let $F=\left(f_{1}, \ldots, f_{n}\right)$ be a system of $W$-homogeneous polynomials, with $W$-degree $D$, and assume that $F$ is in simultaneous Noether position for the variable ordering $X_{1}>X_{2}>\cdots>X_{n}$. Then the $W$-degree of regularity of $F$ is bounded by

$$
\begin{equation*}
d_{\mathrm{reg}, W}(F) \leq \sum_{i=1}^{n}\left(d_{i}-w_{i}\right)+w_{n} . \tag{3.10}
\end{equation*}
$$

Proof. We prove this by induction on $n$. If $n=1$, we simply have one $W$-homogeneous polynomial to consider, and so $d_{\text {reg }, W}=d_{1}$.
Assume now that $n>1$. We consider the system $F^{*}$ defined by:

$$
F^{*}=\left(f_{1}\left(X_{1}, \ldots, X_{n-1}, 0\right), \ldots, f_{n-1}\left(X_{n-1}, \ldots, X_{n-1}, 0\right)\right) .
$$

This system is $W^{*}$-homogeneous, for $W^{*}:=\left(w_{1}, \ldots, w_{n-1}\right)$. From the characterization 3 of Noether position, the sequence $F^{*}$ is in simultaneous Noether position. As a consequence, the induction hypothesis applies to $F^{*}$, and the $W^{*}$-degree of regularity of $F^{*}$ is bounded by

$$
d_{\mathrm{reg}, W^{*}}\left(F^{*}\right) \leq \sum_{i=1}^{n-1}\left(d_{i}-w_{i}\right)+w_{n-1} .
$$

Let $\delta$ be the degree of the Hilbert series of $F$, that is $\delta=\sum_{i=1}^{n}\left(d_{i}-w_{i}\right)$. We want to prove that $d_{\text {reg }} \leq \delta+w_{n}$, i.e. that the Gröbner basis of $F$ need not contain any polynomial with $W$-degree greater than $\delta+w_{n}$. Equivalently, let $\mu$ be a monomial with $W$-degree $d>\delta+w_{n}$. We will prove that $\mu$ is strictly divisible by a monomial in the initial ideal generated by $F$.

Write $\mu=X_{n}^{\alpha} \cdot \mu^{\prime}$, with $\mu^{\prime} \in \mathbb{K}\left[X_{1}, \ldots, X_{n-1}\right]$, and proceed by induction on $\alpha$ :

- if $\alpha=0$, then $\mu \in \mathbb{K}\left[X_{1}, \ldots, X_{n-1}\right]$. By assumption $d_{n} \geq w_{n-1}$, hence:

$$
\begin{equation*}
\delta+w_{n}=\delta^{*}+w_{n-1}-w_{n-1}+d_{n}-w_{n}+w_{n} \geq d_{\mathrm{reg}}^{*}+d_{n}-w_{n-1} \geq d_{\mathrm{reg}}^{*} ; \tag{3.11}
\end{equation*}
$$

and so $\mu$ has $W$-degree greater than $d_{\text {reg. }}^{*}$. By induction hypothesis, $\mu$ is strictly divisible by a monomial in the initial ideal generated by $F^{*}$.

- If $\alpha>0$, then consider $\mu^{\prime \prime}=X_{n}^{\alpha-1} \mu^{\prime}$ : it is a strict divisor of $\mu$. Furthermore, since $\operatorname{deg}(\mu)>\delta+w_{n}$, then $\operatorname{deg}\left(\mu^{\prime \prime}\right)=\operatorname{deg}(\mu)-w_{n}>\delta$. Recall that $\delta$ is by definition the degree of the Hilbert series of the ideal generated by $F$, so $\mu$ lies in that ideal.

Remark 3.29. The hypothesis stating that for any $i, d_{i} \geq w_{i-1}$ is necessary. For example, let $W=(2,1), D=(2,1)$ and the system $F=(X, Y)$ in $\mathbb{K}[X, Y]$, it is $W$-homogeneous with $W$ degree $D$ and in simultaneous Noether position. This system has Hilbert series 1 (the quotient vector span is generated by $\{1\}$ ), which has degree $\delta=0$. But the Gröbner basis of the system is given by $F$ itself, and contains $X$, with $W$-degree 2 .
More generally, without that hypothesis, we obtain the following bound for $d_{\mathrm{reg}, W}(F)$ :

$$
d_{\mathrm{reg}, W}(F) \leq \max \left\{\sum_{i=1}^{k}\left(d_{i}-w_{i}\right)+w_{k}: k \in\{1, \ldots, n\}\right\},
$$

and the proof is the same as that of Theorem 3.28, with the weaker induction hypothesis that $d_{\mathrm{reg}, W}(F) \leq \max \left(\delta+w_{n}, d_{\mathrm{reg}, W^{*}}\left(F^{*}\right)\right)$, which does not need inequality (3.11).
Remark 3.30. We give examples of the behavior of both bounds in Table 3.1: we give the degree of regularity of a generic $W$-homogeneous system of $W$-degree $D$, and show how this degree of regularity varies if we change the order of the weights $W$.
Remark 3.31. Theorem 3.28 gives an indication as to how to choose the order of the variables. Generically, in order to compute a $W$-GRevLex Gröbner basis of the system, the complexity estimates will be better if we set the variables in decreasing weight order.

### 3.3.3. Conjectured exact formula

While the new bound (3.10) is not sharp in full generality, it is sharp whenever $w_{n}=1$. We conjecture that the sharp formula is the following.

Conjecture 3.32. Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a system of weights, and $D=\left(d_{1}, \ldots, d_{n}\right)$ a strongly $W$-compatible system of $W$-degrees. Let $F \in \mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$ be a generic system of $W$-homogeneous polynomials. Let $\delta=\sum_{i=1}^{n}\left(d_{i}-w_{i}\right)$ be the degree of the Hilbert series of $\langle F\rangle$, and let $d_{0}$ be defined as

$$
d_{0}= \begin{cases}\delta+1 & \text { if there exists } i \text { such that } w_{i}=1 \\ \delta-g & \text { otherwise },\end{cases}
$$

where $g$ is the Frobenius number of $W$ (that is, the greatest $W$-degree at which the set of monomials is empty). In other words, $d_{0}$ is the degree of the first "unexpected" zero coefficient in the Hilbert series (by definition of the degree in the first case, and by self-reciprocality of the Hilbert series in the second case).

Then the degree of regularity of $F$ is the first multiple of $w_{n}$ greater than $d_{0}$ :

$$
\begin{equation*}
d_{\mathrm{reg}}=w_{n}\left\lceil\frac{d_{0}}{w_{n}}\right\rceil . \tag{3.12}
\end{equation*}
$$

Table 3.1: Macaulay's bound on the degree of regularity of generic weighted homogeneous systems

| $W$ | $D$ | $d_{\text {reg }}$ | Bound (3.9) | Bound (3.10) |
| :---: | :---: | :---: | :---: | :---: |
| $(3,2,1)$ | $(6,6,6)$ | 13 | 15 | 13 |
| $(3,1,2)$ | $(6,6,6)$ | 14 | 15 | 14 |
| $(1,2,3)$ | $(6,6,6)$ | 15 | 15 | 15 |

### 3.3.4. Positive-dimensional regular sequences

In the positive-dimensional case, the main complexity parameter that we study is the degree of regularity of the system. Under some Noether position assumptions, the bounds that we obtained for zero-dimensional systems apply.

As before, given $F=\left(f_{1}, \ldots, f_{m}\right)$ we define $F_{\text {ext }}=\left(f_{1}, \ldots, f_{m}, X_{m+1}, \ldots, X_{n}\right)$ and $\theta_{m}(F)=$ $F\left(X_{1}, \ldots, X_{m}, 0, \ldots, 0\right)$.

Theorem 3.33. Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a $W$-homogeneous polynomial system with respective $W$-degree $\left(d_{1}, \ldots, d_{m}\right)$. If $F$ is in Noether position with respect to $X_{m+1}, \ldots, X_{n}$, then

$$
d_{\mathrm{reg}} \leq \sum_{i=1}^{m}\left(d_{i}-w_{i}\right)+\max _{1 \leq j \leq m}\left\{w_{j}\right\}
$$

If $F$ is in simultaneous Noether position with respect to the order $X_{1}>\cdots>X_{m}$, then

$$
d_{\mathrm{reg}} \leq \sum_{i=1}^{m}\left(d_{i}-w_{i}\right)+w_{m}
$$

Proof. The proof is analogous to that of Theorem 2.47.

### 3.4. Overdetermined systems

As in the homogeneous case, we want to be able to study systems with more equations than unknowns. For this purpose, we want to generalize the notion of semi-regularity to the weighted case.

We will use the same algebraic definition as in the homogeneous case (Definition 1.45). However, the characterization of Proposition 1.46 does not hold in full generality.

We prove that it does hold in the case of reverse chain-divisible systems of weights however. We then give some consequences on the degree of regularity of the ideal, and show that Fröberg's conjecture is true if $m=n+1$, as in the homogeneous case.

### 3.4.1. Definition of semi-regularity

Let $n$ and $m$ be two integers, $m \geq n, W=\left(w_{1}, \ldots, w_{n}\right)$ a system of weights, and $D=\left(d_{1}, \ldots, d_{m}\right)$ a system of $W$-degrees. Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a system of $W$-homogeneous polynomials with $W$-degree $D$. For any $i \in\{i, \ldots, n\}$, write $F_{i}=\left(f_{1}, \ldots, f_{i}\right)$.

Definition 3.34 (Semi-regularity). We say that $F$ is semi-regular if, for any $i \in\{1, \ldots, m\}$ and for any $d \in \mathbb{N}$, the linear map given by the multiplication by $f_{i}$ :

$$
s_{i, d}:\left(\mathbb{K}\left[X_{1}, \ldots, X_{n}\right] /\left\langle F_{i-1}\right\rangle\right)_{d} \xrightarrow{\cdot f_{i}}\left(\mathbb{K}\left[X_{1}, \ldots, X_{n}\right] /\left\langle F_{i-1}\right\rangle\right)_{d+d_{i}}
$$

is full-rank (either injective or surjective).
Furthermore, let

$$
S_{D, W}(T)=\frac{\prod_{i=1}^{m}\left(1-T^{d_{i}}\right)}{\prod_{i=1}^{n}\left(1-T^{w_{i}}\right)}=\sum_{d=0}^{\infty} a_{d} T^{d} .
$$

We say that $F$ has a semi-regular Hilbert series if the Hilbert series of $F$ is equal to $\left\lfloor S_{D, W}(T)\right\rfloor$, that is the series truncated at the first coefficient less than or equal to zero.

The motivation behind these definitions is that in the homogeneous case, they are equivalent (Proposition 1.46):
Proposition 3.35. If $W=(1, \ldots, 1)$, the following conditions are equivalent:

1. the system $F$ is semi-regular;
2. for any $1 \leq i \leq n$, the system $F_{i}$ has a semi-regular Hilbert series.

For weighted homogeneous systems, the converse implication $(2 \Longrightarrow 1)$ is still true:
Proposition 3.36. Let $F$ be a $W$-homogeneous system such that, for any $1 \leq i \leq n$, the system $F_{i}$ has a semi-regular Hilbert series. Then $F$ is semi-regular.
Proof. We prove this by induction on the number $m$ of polynomials. The initial case is $m=n$, and it is a direct consequence of the characterization of a regular sequence.
Assume $m>n$. Write $R^{*}=\mathbb{K}\left[X_{1}, \ldots, X_{n}\right] /\left\langle f_{1}, \ldots, f_{m-1}\right\rangle$, and for any $d \in \mathbb{N}$, consider the multiplication map

$$
s_{m, d}=R_{d}^{*} \xrightarrow{\cdot f_{m}} R_{d+d_{m}}^{*}
$$

Let $K_{m, d}=\operatorname{ker}\left(s_{m, d}\right)$. Write $S(T)$ the Hilbert series of $F, a_{d}$ its coefficient at degree $d, \delta$ its degree, $H(T)=\left(\prod_{i=1}^{m}\left(1-T^{d_{i}}\right)\right) /\left(\prod_{i=1}^{n}\left(1-T^{w_{i}}\right)\right), b_{d}$ its coefficient at degree $d$, and $S^{*}(T), a_{d}^{*}$, $\delta^{*}, H^{*}(T)$ and $b_{d}^{*}$ their counterparts with $m-1$ polynomials. From the exact sequence

$$
0 \longrightarrow K_{m, d} \longrightarrow R_{d}^{*} \xrightarrow{s_{m, d}} R_{d+d_{m}}^{*} \longrightarrow R_{d+d_{m}} \longrightarrow 0,
$$

we know that the following identity holds

$$
a_{d+d_{m}}=a_{d+d_{m}}^{*}-a_{d}^{*}+\operatorname{dim}\left(K_{m, d}\right) .
$$

We want to prove that either $a_{d+d_{m}}=0$ or $\operatorname{dim}\left(K_{m, d}\right)=0$. Assume that $a_{d+d_{m}}>0$. This means that $d+d_{m} \leq \delta$ and $a_{d+d_{m}}=b_{d+d_{m}}$, so:

$$
\begin{aligned}
a_{d+d_{m}} & =a_{d+d_{m}}-a_{d}^{*}+\operatorname{dim}\left(K_{m, d}\right) \\
& =b_{d+d_{m}} \\
& =b_{d+d_{m}}^{*}-b_{d}^{*} \text { by definition of } H(T) \\
& =a_{d+d_{m}}^{*}-a_{d}^{*} \text { since } \delta^{*} \geq \delta .
\end{aligned}
$$

Thus $K_{m, d}=0$.

### 3.4.2. Characterization with the Hilbert series

In this section, we prove that for reverse chain-divisible systems of weights, semi-regular sequences have a semi-regular Hilbert series. First, we characterize the shape of semi-regular Hilbert series, by extending Theorem 3.25 to the overdetermined case.

Theorem 3.37. Let $m \geq n \geq 0$ be two integers. Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a reverse chain-divisible system of weights, and let $D=\left(d_{1}, \ldots, d_{m}\right)$ be a system of $W$-degrees such that $d_{1}, \ldots, d_{m}$ are all divisible by $w_{1}$. Write

$$
S_{D, W}(T)=\frac{\prod_{i=1}^{m}\left(1-T^{d_{i}}\right)}{\prod_{i=1}^{n}\left(1-T^{w_{i}}\right)}=\sum_{d=0}^{\infty} a_{d} T^{d} .
$$

Then there exist $W$-degrees $\sigma, \delta$ such that

$$
\begin{align*}
& \forall d \in\{1, \ldots, \sigma\}, a_{d} \geq a_{d-1} \\
& a_{\sigma}>a_{\sigma-1} \\
& \forall d \in\{\sigma+1, \ldots, \delta\}, a_{d} \leq a_{d-1} \\
& a_{\delta}>0, a_{\delta+1} \leq 0
\end{align*}
$$

Furthermore, if $m>n$, let $D^{*}=\left(d_{1}, \ldots, d_{m-1}\right)$ and define $\delta^{*}$ as above for the series $S_{D^{*}, W}$. Then the following statements hold:

$$
\begin{cases}\forall d \in\left\{\delta+1, \ldots, \delta^{*}\right\}, a_{d} \leq 0 & \text { ifn }=0 \\ \forall d \in\left\{\delta+1, \ldots, \delta^{*}+d_{m}\right\}, a_{d} \leq 0 & \text { ifn }>0\end{cases}
$$

If $n>0$, let $W^{*}=\left(w_{1}, \ldots, w_{n-1}\right)$, and let $\delta^{\prime}$ be the degree of $\left\lfloor S_{D, W^{*}}(T)\right\rfloor$. If $n=0$, let $\delta^{\prime}=0$. Then the following equality holds

$$
\sigma=\delta^{\prime}
$$

Proof. We prove the theorem by induction on $n$, and for any given $n$, by induction over $m$. The base cases are:

- $n=0, m \geq 0$ : then $S_{D, W}(T)=1-a_{k} T^{k}+O\left(T^{k+1}\right)$ with $a_{k}>0$, and we can conclude, taking $\delta=0$ and $\sigma=0$.
- $n=m>0$ : then this is a consequence of Theorem 3.25 (shape of the Hilbert series of a complete intersection).

Assume that $m>n>0$. Let $D^{*}=\left(d_{1}, \ldots, d_{m-1}\right), W^{*}=\left(w_{1}, \ldots, w_{n-1}\right)$, and write:

$$
\begin{aligned}
& S(T):=S_{D, W}(T)=\sum_{d=0}^{\infty} a_{d} T^{d} \\
& S^{*}(T):=S_{D^{*}, W}(T)=\sum_{d=0}^{\infty} a_{d}^{*} T^{d} .
\end{aligned}
$$

The derivatives of these series are

$$
\begin{aligned}
& \Delta S(T)=S_{D, W^{*}}(T)=\sum_{d=0}^{\infty} a_{d}^{\prime} T^{d} ; \\
& \Delta S^{*}(T)=S_{D^{*}, W^{*}}(T)=\sum_{d=0}^{\infty} a_{d}^{\prime *} T^{d} .
\end{aligned}
$$

Furthermore, let $w=w_{n-1}, \overline{W^{*}}=\left(w_{1} / w, \ldots, w_{n-1} / w\right)$ and $\overline{D^{*}}=\left(d_{1} / w, \ldots, d_{n-1} / w\right)$, and consider the series

$$
\begin{aligned}
& \overline{\Delta S}(T)=S_{\bar{D}, \overline{W^{*}}}(T)=\sum_{d=0}^{\infty} \overline{a^{\prime}} d T^{d} ; \\
& \overline{\Delta S^{*}}(T)=S_{\overline{D^{*}, W^{*}}}(T) \sum_{d=0}^{\infty} \overline{a^{\prime *}} d T^{d} .
\end{aligned}
$$

In particular,

$$
\Delta S(T)=\overline{\Delta S}\left(T^{w}\right) \text { and } \Delta S^{*}(T)=\overline{\Delta S^{*}}\left(T^{w}\right)
$$

All the series $S^{*}, \overline{\Delta S}$ and $\overline{\Delta S^{*}}$ satisfy the induction hypothesis. The $W$-degrees for which the coefficients of $S^{*}$ satisfy properties $(\sigma 1)-(\sigma 4)$ and $(\delta 1)-(\delta 2)$ are denoted by $\sigma^{*}$ and $\delta^{*}$. We write $\overline{\sigma^{\prime}}, \overline{\delta^{\prime}}, \overline{\sigma^{\prime *}}, \overline{\delta^{* *}}$ the respective values of the $W$-degrees for which these properties apply to $\overline{\Delta S}$ and $\overline{\Delta S^{*}}$.

From $S(T)=\left(1-T^{d_{m}}\right) S^{*}(T)$, we deduce the recurrence relation

$$
a_{d}=a_{d}^{*}-a_{d-d_{m}}^{*} .
$$

Since $S^{*}$ satisfies the induction hypothesis, we know that there exists a degree $\delta$ such that

$$
\left\{\begin{array}{ll}
\forall d \in\{0, \ldots, \delta\} & a_{d}^{*}>a_{d-d_{m}}^{*} \\
\forall d \in\left\{\delta+1, \ldots, \delta^{*}+d_{m}\right\} & a_{d}^{*} \leq a_{d-d_{m}}^{*}
\end{array} .\right.
$$

This proves statements $(\delta 1)$ and $(\delta 2)$. As a side result, since $a_{\delta}^{*}>a_{\delta-d_{m}}^{*}$, we also deduce that

$$
\begin{equation*}
\delta-d_{m}<\sigma^{*} \tag{3.13}
\end{equation*}
$$

Let $\sigma=\delta^{\prime}$, we prove that it satisfies equations $(\sigma 1),(\sigma 2)$ and $(\sigma 3)$. We need to evaluate the sign of $a_{d}-a_{d-1}$, depending on $d$. The generating series of $a_{d}-a_{d-1}$ is:

$$
(1-T) S(T)=(1-T) \cdot \frac{\prod_{i=1}^{m}\left(1-T^{d_{i}}\right)}{\prod_{i=1}^{n}\left(1-T^{w_{i}}\right)}=\Delta S(T) \text { since } w_{n}=1
$$

In other words, $a_{d} \geq a_{d-1}$ if and only if $a_{d}^{\prime} \geq 0$, which proves statements ( $\sigma 1$ ) and ( $\sigma 2$ ), by definition of $\delta^{\prime}$ :

$$
\begin{aligned}
& \forall d \in\{0, \ldots, \sigma\}, a_{d}-a_{d-1}=a_{d}^{\prime} \geq 0 \\
& a_{\sigma}-a_{\sigma-1}=a_{\sigma}^{\prime}=a_{\delta^{\prime}}^{\prime}>0
\end{aligned}
$$



Figure 3.6: Shape of the Hilbert series of a semi-regular $W$-homogeneous sequence with $W=(3,3,1)$ and $D=(12,9,6,6,3)$

To finish the proof, we need to prove that for any $d \in\left\{\delta^{\prime}+1, \ldots, \delta\right\}, a_{d}^{\prime} \leq 0$.
From the induction hypothesis (statement $(\sigma 4))$ applied to $S^{*}$, we know that $\delta^{\prime *}=\sigma^{*}$. Moreover, $\overline{\Delta S}$ satisfies the induction hypothesis, and statement ( $\delta 2$ ) yields that:

$$
\forall \bar{d} \in\left\{\overline{\delta^{\prime}}+1, \ldots, \overline{\delta^{* *}}\right\}, \overline{a_{\bar{d}}^{\prime}} \leq 0
$$

As a consequence, since $\sigma^{*}=\delta^{*}=w \overline{\delta^{\prime *}}$ :

$$
\forall d \in\left\{\delta^{\prime}+1, \ldots, \sigma^{*}\right\}, a_{d}^{\prime}= \begin{cases}\overline{a_{\bar{d}}^{\prime}} \leq 0 & \text { if } d=w \bar{d}  \tag{i}\\ 0 & \text { otherwise }\end{cases}
$$

Now assume that $\sigma^{*}<d \leq \delta$. We can write $a_{d}^{\prime}$ as

$$
\begin{aligned}
a_{d}^{\prime}=a_{d}-a_{d-1} & =a_{d}^{*}-a_{d-d_{m}}^{*}-a_{d-1}^{*}+a_{d-d_{m}-1}^{*} \\
& =\left(a_{d}^{*}-a_{d-1}^{*}\right)-\left(a_{d-d_{m}}^{*}-a_{d-d_{m}-1}^{*}\right)=a_{d}^{*}-a_{d-d_{m}}^{* *}
\end{aligned}
$$

Since $a_{d} \leq a_{d}^{*}$ for any $d$, we necessarily have $\delta \leq \delta^{*}$, hence $\sigma^{*}<d \leq \delta^{*}$. Then, by induction hypothesis (statement ( $\sigma 3$ )), we know that $a_{d}^{*}-a_{d-1}^{*} \leq 0$. Additionally, equation (3.13) and induction hypothesis (statement $(\sigma 1)$ ) together yield that $a_{d-d_{m}}^{*}-a_{d-d_{m}-1}^{*} \geq 0$, so we conclude that

$$
\begin{equation*}
\forall d \in\left\{\sigma^{*}+1, \ldots, \delta\right\}, a_{d}^{\prime} \leq 0 \tag{ii}
\end{equation*}
$$

Sticking the ranges of statements (i) and (ii) together, we prove statement ( $\sigma 3$ ) which completes the proof.

Using this description of semi-regular Hilbert series, we now prove that for reverse chaindivisible systems of weights, semi-regular sequences have a semi-regular Hilbert series. As an illustration, Figure 3.6 shows the coefficient of a semi-regular Hilbert series. The black dots correspond to the actual coefficients, and the gray dots are the coefficients which were truncated away.

Theorem 3.38. Let $m \geq n \geq 0$ be two integers, $W=\left(w_{1}, \ldots, w_{n}\right)$ be a reverse chain-divisible system of weights and $D=\left(d_{1}, \ldots, d_{m}\right)$ be a system of $W$-degrees such that $d_{1}, \ldots, d_{n}$ are all
divisible by $w_{1}$. Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a system of $W$-homogeneous polynomials, with respective $W$-degree $D$.

If $F$ is a semi-regular sequence, then $F$ has a semi-regular Hilbert series.
Proof. We proceed by induction on $m$. If $m=n$, then the result is a consequence of the characterization of regular sequences.

Assume that $m>n$. We consider the series $S(T)=S_{D, W}(T)$ with generic coefficient $a_{d}$, $S^{*}(T)=S_{D^{*}, W}(T)$ with generic coefficient $a_{d}^{*}, H(T)$ the Hilbert series of $F$ with generic coefficient $b_{d}$, and $H^{*}(T)$ the Hilbert series of $F^{*}:=\left(f_{1}, \ldots, f_{m-1}\right)$ with generic coefficient $b_{d}^{*}$.

By induction hypothesis, $H^{*}(T)=\left\lfloor S^{*}(T)\right\rfloor$. Furthermore, since $F$ is semi-regular, we have the exact sequence

$$
0 \longrightarrow K_{m, d} \longrightarrow R_{d}^{*} \xrightarrow{s_{m, d}} R_{d+d_{m}}^{*} \longrightarrow R_{d+d_{m}} \longrightarrow 0
$$

where $R=\mathbb{K}\left[X_{1}, \ldots, X_{n}\right] /\langle F\rangle$ and $R^{*}=\mathbb{K}\left[X_{1}, \ldots, X_{n}\right] /\left\langle F^{*}\right\rangle$. As a consequence, for any $d \geq 0$, the coefficient $b_{d}$ satisfies the recurrence relation:

$$
b_{d+d_{m}}=b_{d+d_{m}}^{*}-b_{d}^{*}+\operatorname{dim}\left(K_{m, d}\right)
$$

where either $K_{m, d}=0$ or $b_{d+d_{m}}=0$. Since $s_{m, d}$ is defined from a space of dimension $b_{d}^{*}$ to a space of dimension $b_{d+d_{m}}^{*}$, this can be rephrased as

$$
b_{d}=\max \left(b_{d}^{*}-b_{d-d_{m}}^{*}, 0\right) .
$$

From Theorem 3.37 applied to $S(T)$, there exists $\delta$ such that

$$
\forall d \in\{0, \ldots, \delta\}, a_{d}=a_{d}^{*}-a_{d-d_{m}}^{*}>0
$$

Furthermore, the induction hypothesis shows that there exists a degree $\delta^{*}$ such that
$\forall d \in\left\{0, \ldots, \delta^{*}\right\}, a_{d}^{*}=b_{d}^{*}>0$
$\forall d>\delta^{*}, b_{d}^{*}=0$,
and that $\delta^{*}$ is defined as in Theorem 3.37. In particular, it implies that $\delta^{*} \geq \delta$.
We shall prove that the Hilbert series $H$ of $F$ is equal to $S$, truncated at degree $\delta$. Let $d \geq 0$ :

- if $0 \leq d \leq \delta \leq \delta^{*}$ :

$$
\begin{aligned}
b_{d} & =b_{d}^{*}-b_{d_{m}}^{*} \text { since } d \leq \delta \\
& =a_{d}^{*}-a_{d_{m}}^{*} \text { since } d \leq \delta^{*} \\
& =a_{d}
\end{aligned}
$$

- if $\delta<d$ :

$$
\begin{aligned}
b_{d} & =\max \left(b_{d}^{*}-b_{d-d_{m}}^{*}, 0\right) \\
& =0 \text { since } b_{d}^{*}=0 \text { and } b_{d-d_{m}}^{*} \geq 0
\end{aligned}
$$

And since $a_{\delta+1} \leq 0$, this proves that

$$
H(T)=\lfloor S(T)\rfloor .
$$

### 3.4.3. $W$-degree of regularity of semi-regular sequences

Another consequence of Theorem 3.37 is an explicit value for the degree $\delta$ of the Hilbert series of an ideal defined by a semi-regular sequence with $m=n+1$ polynomials in $n$ variables. In the homogeneous case, it is known that this degree is bounded by

$$
\delta=\min \left(\sum_{i=1}^{n} d_{i}-n,\left\lfloor\frac{\sum_{i=1}^{n+1} d_{i}-n}{2}\right\rfloor\right) .
$$

Proposition 3.39. Let $n$ be a positive integer, and $m=n+1$. Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a system of weights, and $F=\left(f_{1}, \ldots, f_{m}\right)$ a system of $W$-homogeneous polynomials, and assume that the hypotheses of Theorem 3.38 are satisfied. For all $i \in\{1, \ldots, m\}$, let $d_{i}:=\operatorname{deg}_{W}\left(f_{i}\right)$. Then the degree $\delta$ of the Hilbert series of $\langle F\rangle$ is given by:

$$
\delta=\min \left(\sum_{i=1}^{n} d_{i}-\sum_{i=1}^{n} w_{i},\left[\frac{\sum_{i=1}^{n+1} d_{i}-\sum_{i=1}^{n} w_{i}}{2}\right]\right) .
$$

Proof. Consider the system of weights $W^{+}=\left(w_{1}, \ldots, w_{n}, 1\right)$, and the series $S_{D, W^{+}}$as defined in Theorem 3.37. It satisfies the hypotheses of Theorem 3.25, which implies that its coefficients are increasing up to degree

$$
\sigma^{+}=\min \left(\sum_{i=1}^{n} d_{i}-\sum_{i=1}^{n} w_{i},\left[\frac{\sum_{i=1}^{n+1} d_{i}-\sum_{i=1}^{n} w_{i}}{2}\right]\right) .
$$

Theorem 3.37 (statement ( $\sigma 4$ )) states that the degree $\delta$ of the Hilbert series of $\langle F\rangle$ satisfies

$$
\delta=\sigma^{+}
$$

hence the result.
In the end of this section, we show how the results from [Bar+05] and [Baro4, Chap. 4] about the degree of regularity of semi-regular homogeneous sequences can be adapted to the weighted case.
Theorem 3.40. Let $k$ and $n$ be non-negative integers and let $m:=n+k$. Let $w_{0}$ and $d_{0}$ be nonnegative integers such that $w_{0} \mid d_{0}$. Consider the system of $n$ weights $W=\left(w_{0}, \ldots, w_{0}, 1\right)$ and the system of $m W$-degrees $D=\left(d_{0}, \ldots, d_{0}\right)$. Let $F=\left(f_{1}, \ldots, f_{m}\right) \subset A=\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$ be a semiregular sequence of weighted homogeneous polynomials with $W$-degree $D$. Then the asymptotic expansion of the $W$-degree of regularity $d_{\mathrm{reg}}$ of $F$ as $n$ tends to infinity is given by

$$
d_{\mathrm{reg}}=n\left(\frac{d_{0}-w_{0}}{2}\right)-\alpha_{k} \sqrt{n\left(\frac{d_{0}^{2}-w_{0}^{2}}{6}\right)}+O\left(n^{1 / 4}\right) .
$$

Remark 3.41. In the non-weighted case, this asymptotic expansion is

$$
d_{\mathrm{reg}}=n\left(\frac{d_{0}-1}{2}\right)-\alpha_{k} \sqrt{n\left(\frac{d_{0}^{2}-1}{6}\right)}+O\left(n^{1 / 4}\right) .
$$

Overall, the bound is improved by $O\left(n w_{0}\right)=O\left(\sum w_{i}\right)$.

Proof. Let $I:=\langle F\rangle$. The Hilbert series of $A / I$ is given by

$$
\operatorname{HS}_{A / I}(T)=\left\lfloor\frac{\left(1-T^{d_{0}}\right)^{m}}{\left(1-T^{w_{0}}\right)^{n-1}(1-T)}\right\rfloor .
$$

Write

$$
\begin{aligned}
& H(T)=\frac{\left(1-T^{d_{0}}\right)^{m}}{\left(1-T^{w_{0}}\right)^{n-1}(1-T)}=\sum_{d=0}^{\delta} a_{d} T^{d} ; \\
& H^{*}(T)=\frac{\left(1-T^{d_{0} / w_{0}}\right)^{m-1}}{(1-T)^{n-1}}=\sum_{d=0}^{\delta} a_{d}^{*} T^{d},
\end{aligned}
$$

these series are related through

$$
H(T)=H^{*}\left(T^{w_{0}}\right) \cdot \frac{1-T^{d_{0}}}{1-T}=H^{*}\left(T^{w_{0}}\right) \cdot\left(1+T+\cdots+T^{d_{0}-1}\right) .
$$

It means for the coefficients that for any $d$ in $\mathbb{N}$ :

$$
a_{d}=a_{\left\lfloor d / w_{0}\right\rfloor}^{*}+\cdots+a_{\left\lceil\left(d-d_{0}+1\right) / w_{0}\right\rceil}^{*}
$$

If truncated before its first non-positive coefficient, the series $H^{*}$ is the Hilbert series of a semi-regular 1 -homogeneous sequence of $m-1$ polynomials in $n-1$ variables, with degree $d_{0} / w_{0}$. Let $\delta^{*}$ be the degree of this truncated series, so that $\delta^{*}+1$ is an upper bound for its degree of regularity. Statement ( $\delta 2$ ) of Theorem 3.37 states that:

$$
\forall d \in\left\{\delta^{*}+1, \ldots, \delta^{*}+d_{0} / w_{0}\right\}, a_{d}^{*} \leq 0
$$

Let $\delta_{0}:=w_{0} \delta^{*}+d_{0}$, we have

$$
\delta^{*}<\frac{\delta_{0}-d_{0}+1}{w_{0}} \leq \delta^{*}+1
$$

and

$$
\frac{\delta_{0}}{w_{0}}=\delta^{*}+\frac{d_{0}}{w_{0}}
$$

and as a consequence

$$
a_{\delta^{0}}=a_{\left\lfloor\delta_{0} / w_{0}\right\rfloor}^{*}+\cdots+a_{\left\lceil\left(\delta_{0}-d_{0}+1\right) / w_{0}\right\rceil}^{*} \leq 0 .
$$

In other words, the degree of regularity $d_{\mathrm{reg}}$ of $F$ is bounded by

$$
w_{0} \delta^{*}<d_{\mathrm{reg}} \leq \delta_{0}=w_{0} \delta^{*}+d_{0}
$$

The degree $\delta^{*}$ is the degree of the Hilbert series of a homogeneous semi-regular sequence, and as such it follows the asymptotic estimates proved in [Baro4, Chap. 4]. For example in our
setting where $k$ is an integer and $m=n+k$, the asymptotic expansion of $\delta^{*}$ when $n$ tends to infinity is given by

$$
\delta^{*}+1=n \frac{d_{0} / w_{0}-1}{2}-\alpha_{k} \sqrt{n \frac{\left(d_{0} / w_{0}\right)^{2}-1}{6}}+O\left(n^{1 / 4}\right)
$$

where $\alpha_{k}$ is the largest root of the $k$ 'th Hermite polynomial. ${ }^{1}$
As a consequence,

$$
\begin{aligned}
d_{\mathrm{reg}} & =w_{0} \delta^{*}+O(1) \\
& =w_{0}\left(n \frac{d_{0} / w_{0}-1}{2}-\alpha_{k} \sqrt{n \frac{\left(d_{0} / w_{0}\right)^{2}-1}{6}}+O\left(n^{1 / 4}\right)\right)+O(1) \\
& =n \frac{d_{0}-w_{0}}{2}-\alpha_{k} \sqrt{n \frac{d_{0}^{2}-w_{0}^{2}}{6}}+O\left(n^{1 / 4}\right)
\end{aligned}
$$

### 3.4.4. Fröberg's conjecture

Fröberg's conjecture states that homogeneous semi-regular sequences are generic among sequences of fixed degree. The fact that semi-regularity is a Zariski-open condition is a known fact (the proof is the same as for regularity), so the conjecture states that for any system of degrees, there exists a semi-regular homogeneous sequence with these degrees.

This conjecture extends naturally to the weighted case. In this case, semi-regularity is still a Zariski-open condition.

We extend here one known result from the homogeneous case (see for example [RRR91]), stating that Fröberg's conjecture is true in characteristic 0 if $m=n+1$. We follow the proof given in [RRR91].

Proposition 3.42. Let $m=n+1, W=\left(w_{1}, \ldots, w_{n}\right)$ a reverse chain-divisible system of weights, $D=\left(d_{1}, \ldots, d_{n}\right)$ a strongly $W$-compatible system of degrees, and $d_{n+1}$ an integer divisible by $w_{1}$. Write $f_{n+1}=\left(X_{1}+X_{2}^{w_{1} / w_{2}}+\cdots+X_{n}^{w_{1}}\right)^{d_{n+1} / w_{1}}$.

Then the sequence $F:=\left(X_{1}^{d_{1} / w_{1}}, \ldots, X_{n}^{d_{n} / w_{n}}, f_{n+1}\right)$ is semi-regular.
Lemma 3.43. Let $f$ be a polynomial such that

$$
f \cdot f_{n+1}=0 \in A=\mathbb{K}\left[X_{1}, \ldots, X_{n}\right] /\left(X_{1}^{d_{1} / w_{1}}, \ldots, X_{n}^{d_{n} / w_{n}}\right) .
$$

Let $\delta=\sum_{i=1}^{n}\left(d_{i}-w_{i}\right)$, then we have

$$
\operatorname{deg}_{W}(f) \geq \frac{\delta-d_{n+1}+1}{2}
$$

[^0]Proof. If the $W$-degree of $f$ is 0 , this means that $\left(X_{1}+X_{2}^{w_{1} / w_{2}}+\cdots+X_{n}^{w_{1}}\right)^{d_{n+1} / w_{1}}=0$ in $A$. Assume that $\operatorname{deg}_{W}(f)<\left(\delta-d_{n+1}+1\right) / 2$, that means that $\delta-d_{n+1}+1 \geq 1$, so $\delta \geq d_{n+1}$. Consider the expansion of $f_{n+1}$. All its coefficients are nonzero since the base field has characteristic 0 . Its support is the set of monomials of degree $d_{n+1}$. Since $d_{n+1} \leq \delta, \operatorname{dim}\left(\mathbb{K}[\mathrm{X}] /\left\langle f_{1}, \ldots, f_{n}\right\rangle\right)_{d_{n+1}}>0$, which means that there exists at least one monomial with $W$-degree $d_{n+1}$ which does not lie in the initial ideal of $\left\langle f_{1}, \ldots, f_{n}\right\rangle$. As a consequence, $f$ is non-zero in the quotient.
Now assume that $\operatorname{deg}_{W}(f)>0$. Write $B=\mathbb{K}\left[X_{2}, \ldots, X_{n}\right] /\left(X_{2}^{d_{2} / w_{2}}, \ldots, X_{n}^{d_{n} / w_{n}}\right), X=X_{1}$, $R=B[X], d=d_{1} / w_{1}$, such that $A=R / X^{d}$. Let $S=\left(X+X_{2}^{w_{1} / w_{2}}+\cdots+X_{n}^{w_{1}}\right)$, and let $F$ be a weighted homogeneous polynomial in $R$ with image $f$ in $A$. By assumption, there exists $G \in R$ such that $S^{d_{n+1} / w_{1}} \cdot F=G \cdot X^{d}$. Differentiate this equality along $X$ to obtain:

$$
m S^{\prime} S^{d_{n+1} / w_{1}-1} F+S^{\left(d_{n+1}\right) / w_{1}} F=d G^{\prime} X^{d-1}+G^{\prime} X^{d}
$$

which gives, modulo $X^{d-1}$

$$
\begin{aligned}
S^{d_{n+1}-1}\left(m F+S F^{\prime}\right) \equiv 0[X]^{d-1} & \Longrightarrow S^{d_{n+1} / w_{1}}\left(m F+S F^{\prime}\right) \equiv 0[X]^{d-1} \\
& \Longrightarrow S^{d_{n+1} / w_{1}+1} F^{\prime} \equiv m F S^{d_{n+1} / w_{1}} \equiv 0[X]^{d-1}
\end{aligned}
$$

Since $X=X_{1}$ has weight $w_{1}, F^{\prime}$ is $W$-homogeneous with $W$-degree $\operatorname{deg}_{W}(f)-w_{1}$, and we can use the induction hypothesis on $F^{\prime}[X] \in A$ and $\operatorname{deg}(F)=d_{n+1}+w_{1}$ to deduce:

$$
\begin{aligned}
\operatorname{deg}_{W}(f) & =\operatorname{deg}_{W}(F)=\operatorname{deg}_{W}\left(F^{\prime}\right)+1 \\
& \geq \frac{(\delta-1)-\left(d_{n+1}+1\right)+1}{2}+1 \\
& \geq \frac{\delta-d_{n+1}+1}{2} .
\end{aligned}
$$

Proof of the proposition. The proof given in [RRR91, before prop. 7] still holds in the weighted case.

### 3.5. Algorithms and complexity

### 3.5.1. Critical pairs algorithms

As said previously, we can compute a Gröbner basis of a weighted homogeneous ideal generated by a system $F$ by running usual algorithms on $\operatorname{hom}_{W}(F)$. The following proposition shows that this is correct:

Proposition 3.44. Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a family of polynomials in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$, assumed to be weighted homogeneous for a system of weights $W=\left(w_{1}, \ldots, w_{n}\right)$. Let $<_{1}$ be a monomial order, $G$ the reduced Gröbner basis of $\operatorname{hom}_{W}(F)$ for this order, and $<_{2}$ the pullback of $<_{1}$ through hom $_{W}$. Then

1. all elements of $G$ are in the image of $\mathrm{hom}_{W}$;
2. the family $G^{\prime}:=\operatorname{hom}_{W}^{-1}(G)$ is a reduced Gröbner basis of the system $F$ for the order $<_{2}$.

Proof. The morphism hom $_{W}$ preserves $S$-polynomials, in the sense that

$$
S-\operatorname{pol}\left(\operatorname{hom}_{W}(f), \operatorname{hom}_{W}(g)\right)=\operatorname{hom}_{W}(S-\operatorname{pol}(f, g)) .
$$

Recall that we can compute a reduced Gröbner basis by running the Buchberger algorithm, which involves only multiplications, additions, tests of divisibility and $S$-polynomials. Since all these operations are compatible with hom $_{W}$, if we run the Buchberger algorithm on both $F$ and hom $_{W}(F)$ simultaneously, they will follow exactly the same computations up to application of $\operatorname{hom}_{W}$. The consequences on the final reduced Gröbner basis follow.

This shows that all pairs-based algorithms (for example Buchberger, $F_{4}$ and $F_{5}$ ) can be made to take into account a weighted homogeneous structure, simply by wrapping hom ${ }_{W}$ around the algorithm.

### 3.5.2. Adapting the Matrix- $F_{5}$ algorithm

In order to study the complexity of algorithm $F_{5}$, we want to run algorithm Matrix- $F_{5}$ on a weighted homogeneous system. Recall that this algorithm works by constructing a matrix of linear bases of $I_{d}$ incrementally, and does not use critical pairs to guide the reductions.

If we wrap the algorithm with $\mathrm{hom}_{W}$, Proposition 3.44 applies and shows that the result is correct. However, it also states that we can compute a Gröbner basis of $\left\langle\operatorname{hom}_{W}(F)\right\rangle$ by working with $S$-polynomials in the image of $\operatorname{hom}_{W}$, so those rows whose label monomial does not lie in this image will be effectively useless.

There are two equivalent ways around this problem. We can run the algorithm on hom ${ }_{W}(F)$, only considering monomials in the image of $\operatorname{hom}_{W}$ for the labels. Or we can adapt Matrix- $\mathrm{F}_{5}$ to use the weighted degree instead of the total degree. We present the second solution.

The only real difference is at line 14: if $\mu$ has $W$-degree $d-d_{i}, \mu / x_{j}$ has $W$-degree $d-d_{i}-w_{j}$ and so the polynomial $f$ comes from a row of the matrix at $W$-degree $d-d_{i}-w_{j}+d_{i}=d-w_{j}$.

Theorem 3.45. Let $F=\left(f_{1}, \ldots, f_{m}\right)$ be a system of $W$-homogeneous polynomials, and let $P=$ $\prod_{i=1}^{n} w_{i}$. Then Matrix- $\mathrm{F}_{5}\left(F, d_{\max }\right)$ computes a $d_{\max }-G r o ̈ b n e r ~ b a s i s ~ o f ~\langle F\rangle$ in time

$$
\begin{equation*}
O\left(d_{\mathrm{reg}} \frac{1}{P^{3}}\binom{n+d_{\max }-1}{d_{\max }}^{3}\right) \tag{3.14}
\end{equation*}
$$

If $d_{\max } \geq d_{\mathrm{reg}}$ and the monomial order is $W$-GREVLEX, then Matrix- $\mathrm{F}_{5}\left(F, d_{\max }\right)$ computes a Gröbner basis of $F$. If $F$ is in Noether position with respect to the variables $X_{m+1}, \ldots, X_{n}$, then $d_{\mathrm{reg}}$ is bounded by

$$
d_{\mathrm{reg}} \leq \sum_{i=1}^{m}\left(d_{i}-w_{i}\right)+\max \left\{w_{j}\right\}
$$

If $F$ is in simultaneous Noether position with respect to the order $X_{1}>X_{2}>\cdots>X_{m}$, then $d_{\text {reg }}$ is bounded by

$$
d_{\mathrm{reg}} \leq \sum_{i=1}^{m}\left(d_{i}-w_{i}\right)+w_{m}
$$

```
Algorithm 3.1 Matrix- \(\mathrm{F}_{5}\) (weighted case)
Input: \(G \subset \mathbb{K}[\mathbf{X}]\) a truncated Gröbner basis of \(\langle F\rangle\) at \(W\)-degree \(d_{\text {max }}\)
Output: \(F \subset \mathbb{K}[\mathbf{X}], d_{\max } \in \mathbb{N}\) an implicit monomial ordering, \(W \in \mathbb{N}^{n}\)
    \(G \leftarrow \varnothing\)
    for \(d=0\) to \(d_{\text {max }}\) do
        \(N_{d} \leftarrow\) number of monomials of \(W\)-degree \(d\)
        \(M_{d, 0}, \tilde{M}_{d, 0} \leftarrow\) matrix with o rows and \(N_{d}\) columns
        for \(i=1\) to \(m\) do
            \(M_{d, i} \leftarrow M_{d, i-1}\)
            if \(d=d_{i}\) then
                \(M_{d, i} \leftarrow M_{d, i} \cup\) row \(f_{i}\) with signature \((i, 1)\)
            else if \(d>d_{i}\) then
                for all \(\mu\) monomial of \(W\)-degree \(d-d_{i}\) do
                        if \(\mu\) is not the leading term of a row in \(\tilde{M}_{d-d_{i}, i-1}\) then
                        \(j \leftarrow\) largest \(j\) such that \(X_{j}\) divides \(m u\)
                                    \(\mu_{\tilde{f}}^{\prime} \leftarrow \mu / X_{j}\)
                                    \(\tilde{f} \leftarrow\) row of \(\tilde{M}_{d-w_{j}, i}\) with signature \(\left(i, \mu^{\prime}\right)\)
                                    \(M_{d, i} \leftarrow M_{d, i} \cup\) row \(X_{j} \tilde{f}\) with signature \((i, \mu)\)
                    end if
                end for
            end if
        end for
        \(\tilde{M}_{d, i} \leftarrow \operatorname{EchelonForm}\left(M_{d, i}\right)\)
        \(R \leftarrow\) number of rows in \(M_{d, i}\)
        \(G \leftarrow G \cup\left\{\operatorname{row}\left(\tilde{M}_{d, i}, j\right) \mid j \in\{1, \ldots, R\}, \operatorname{row}\left(\tilde{M}_{d, i}, j\right) \neq \operatorname{row}\left(M_{d, i}, j\right)\right\}\)
    end for
    return \(G\)
```

Proof. The complexity is computed as in the homogeneous case, in terms of the maximal degree $d_{\text {max }}$ and the size of the matrices. The size of the matrices depends on the number of monomials at a given $W$-degree, and Proposition 1.19 shows that asymptotically, this number of monomials is divided by $P$. The bounds on $d_{\mathrm{reg}}$ come respectively from Theorem 3.27 and Theorem 3.28.

### 3.5.3. Thin-grained complexity of Matrix- $F_{5}$

As in the homogeneous case, we wish to obtain sharper complexity bounds. We will show that the computations of Section 2.5 .3 can be carried in a weighted setting.

Let $W=\left(w_{1}, \ldots, w_{n}\right)$ be a system of weights, and $\left(f_{1}, \ldots, f_{m}\right) \in \mathbb{K}[\mathbf{X}]$ a system of weighted homogeneous polynomials in simultaneous Noether position with respect to the order $X_{1}>$ $\cdots>X_{n}$. Let $d_{1}, \ldots, d_{m}$ be the respective $W$-degrees of the polynomials $f_{1}, \ldots, f_{m}$.

We also denote by:

- $A_{i}=\mathbb{K}\left[X_{1}, \ldots, X_{i}\right]$, and $A=A_{n}$;
- $S_{i}=\sum_{j=2}^{i} w_{i} \frac{\operatorname{gcd}\left(w_{1}, \ldots, w_{j}\right)}{\operatorname{gcd}\left(w_{1}, \ldots, w_{j-1}\right)}$ (see the definition of $A$ in Prop. 1.20), and $S=S_{n}$;
- $P_{i}=\prod_{j=1}^{i} w_{j}$, and $P=P_{n}$;
- $I_{i}=\left\langle f_{1}, \ldots, f_{i}\right\rangle$, and $I=I_{m}$;
- $\operatorname{hom}_{W}\left(f_{j}\right)=\operatorname{hom}_{W}\left(f_{j}\right)$;
- $\operatorname{hom}_{W}\left(I_{i}\right)=\left\langle\operatorname{hom}_{W}\left(f_{1}\right), \ldots, \operatorname{hom}_{W}\left(f_{i}\right)\right\rangle$, and $\operatorname{hom}_{W}(I)=\operatorname{hom}_{W}\left(I_{m}\right)$;
- $D_{i}=\prod_{j=1}^{i} \frac{d_{j}}{w_{j}}\left(\right.$ with the convention $\left.D_{0}=0\right)$;
- $\tilde{D}_{i}=\operatorname{deg}\left(\operatorname{hom}_{W}\left(I_{i}\right)\right)=\prod_{j=1}^{i} d_{j}$;
- $d_{\mathrm{reg}}^{(i)}$ the degree of regularity of $F_{i}\left(\right.$ or of $\left.\operatorname{hom}_{W}\left(F_{i}\right)\right)$;
- $G_{i}$ the $W$-GRevLex Gröbner basis of $I_{i}$ as given by Matrix- $\mathrm{F}_{5}$.

With these notation and hypotheses, we are going to prove the following theorem:
Theorem 3.46. The complexity of weighted homogeneous Matrix- $F_{5}$ algorithm (Algorithm 3.1) is:

$$
\begin{aligned}
C_{\mathrm{F}_{5}} & =O\left(\sum_{i=2}^{m}\left(D_{i-1}-D_{i-2}\right) M_{d_{\mathrm{reg}}, W}^{(i)}(i) M_{d_{\mathrm{reg}}, W}^{(i)}(n)\right) \\
& =O\left(\sum_{i=2}^{m} \frac{D_{i-1}-D_{i-2}}{P_{i} P_{n}}\binom{i+d_{\mathrm{reg}}^{(i)}-1}{d_{\mathrm{reg}}^{(i)}}\binom{n+d_{\mathrm{reg}}^{(i)}-1}{d_{\mathrm{reg}}^{(i)}}\right)
\end{aligned}
$$

As in the homogeneous case, this theorem relies on the structure lemma:
Lemma 3.47 (Structure lemma, weighted case). Let $i \in\{1, \ldots, m\}$, and $g \in G_{i}$ with signature $(j, \mu)$. Then:

- $j \leq i$
- $\operatorname{LT}(g) \in \mathbb{K}\left[X_{1}, \ldots, X_{i}\right]$
- $\mu \in \mathbb{K}\left[X_{1}, \ldots, X_{i-1}\right]$

Proof. The proof of the homogeneous version of the lemma (Lemma 2.51) does not depend on the graduation.

With the structure lemma, we can count the number of polynomials in each intermediate basis, $W$-degree by $W$-degree:

Proposition 3.48. Let $\left(f_{1}, \ldots, f_{m}\right)$ be a $W$-homogeneous system in simultaneous Noether position with respect to the order $X_{1}>\cdots>X_{n}$. Let $G_{i}$ be a reduced Gröbner basis of $\left(f_{1}, \ldots, f_{i}\right)$ for the $W$-GREVLEX monomial ordering, for $1 \leq i \leq m$. Then the number of polynomials of $W$ degree $d$ in $G_{i}$ whose leading term does not belong to $\operatorname{LT}\left(G_{i-1}\right)$ is bounded by $b_{d, i}$, defined by the generating series

$$
\begin{equation*}
B_{W, i}(T)=\sum_{d=0}^{\infty} b_{d}^{(W, i)} T^{d}=T^{d_{i}} \prod_{k=1}^{i-1} \frac{1-T^{d_{k}}}{1-T^{w_{k}}} . \tag{3.15}
\end{equation*}
$$

Proof. The proof of Proposition 2.52 still holds in a weighted setting, substituting formula (3.18) for the Hilbert series of a weighted homogeneous regular sequence.

As in Section 2.1.3, we define for any $l \leq n$ the endomorphism $\theta_{l}$ of $\mathbb{K}[\mathbf{X}]$, sending $X_{i}$ on itself for $i \leq l$ and on 0 otherwise. Note that by Proposition 1.44, $\theta_{l}(F)$ is still in Noether position for all $l \in\{1, \ldots, m\}$.

Theorem 3.49. The complexity of weighted homogeneous Matrix- $F_{5}$ algorithm (Algorithm 3.1) is:

$$
\begin{align*}
C_{\mathrm{F}_{5}} & =O\left(\sum_{i=2}^{m} \sum_{d=0}^{\infty} b_{d+d_{i}}^{(W, i)} N_{W_{i}, d+d_{i}} N_{W, d+d_{i}}\right) \\
& =O\left(\sum_{i=2}^{m} \sum_{d=0}^{\infty} \frac{b_{d+d_{i}}^{(W, i)}}{P_{i} P_{n}}\binom{d+d_{i}+S_{i}}{d+d_{i}+S_{i}-i+1}\binom{d+d_{i}+S_{i}}{d+d_{i}+S_{i}-n+1}\right) \tag{3.16}
\end{align*}
$$

Proof of Theorem 3.49. The proof of Theorem 2.54 adapts to a weighted setting: $b_{d+d_{i}}^{(W, i)}$ bounds the number of polynomials reduced with signature $(i, \bullet)$ and $W$-degree $d, N_{W_{i}, d}$ bounds the number of rows used for each reduction at degree $d$ for a polynomial with signature $(i, \bullet)$, and $N_{W, d}$ bounds the length of the rows. The second equality follows from the bounds of Proposition 1.20 on $N_{W_{i}, d}$.

Proof of Theorem 3.46. With the same notations, for any $i \in\{1, \ldots, m\}, B_{W, i}(1)$ represents the number of reduced polynomials during the computation of a Gröbner basis of $\left(f_{1}, \ldots, f_{i}\right)$. Equation (3.15) states that $B_{W, i}(T) / T^{d_{i}}$ is the Hilbert series of $\left(\theta_{i}\left(f_{1}\right), \ldots, \theta_{i}\left(f_{i-1}\right)\right)$, and so its


Figure 3.7: Number of operations needed by Matrix- $F_{5}$ on a generic [4, 2, 1]-homogeneous system with degree $[d, d, d]$
value at $T=1$ is the degree of that ideal, or $D_{i-1}$. Therefore we know that the number of reduced polynomials with label $\left(\mu, f_{i}\right)$ is bounded by $D_{i-1}-D_{i-2}$.

The rest of the factors are obtained in the same way as in the proof of Theorem 3.49, using the fact that all polynomials with label $(i, \bullet)$ are reduced at a $W$-degree bounded by $d_{\text {reg }}^{(i)}$.

The behavior of both bounds is shown on Figure 3.7.

### 3.5.4. FGLM and computational strategy for zero-dimensional systems

Let $F=\left(f_{1}, \ldots, f_{n}\right) \subset \mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$ be a zero-dimensional affine system. In this situation, we usually want to enumerate all solutions of the system. As described in Section 2.3.1, one can generically obtain a triangular description of the solutions from a lexicographical basis of the ideal, and it can be computed from another Gröbner basis using algorithm FGLM.

Recall from Section 2.5 .4 that the complexity of algorithm FGLM mainly depends on the degree of the ideal. The weighted version of Bézout's bound implies that at least generically, applying hom $_{W}$ to a zero-dimensional ideal multiplies the number of solutions by the product of the weights. So wrapping $\operatorname{hom}_{W}$ around calls to FGLM incurs an additional complexity factor of $\prod w_{i}$.

On the other hand, algorithm FGLM does not need adapting in order to work with systems with a weighted homogeneous structure: it takes as input a Gröbner basis for any order, in particular one can use a $W$-GRevLex basis.

All in all, Bézout's bound gives the following complexity for FGLM for a weighted system:
Theorem 3.50. The complexity of algorithm FGLM for computing the Gröbner basis of a zerodimensional ideal defined by $\left(f_{1}, \ldots, f_{n}\right)$, polynomials with a $W$-homogeneous structure, and with
respective $W$-degrees $\left(d_{1}, \ldots, d_{n}\right)$.

$$
\begin{equation*}
O\left(n\left(\frac{\prod_{i=1}^{n} d_{i}}{\prod_{i=1}^{n} w_{i}}\right)^{3}\right) . \tag{3.17}
\end{equation*}
$$

Experimentally, the improvement of $[\mathrm{Fau}+14]$ appears to work for bases for the $W$-GRevLex order, so that $\omega$ can be taken as the complexity exponent.

### 3.5.5. Other algorithms

The gain from the reduced number of monomials applies to other algorithms as well, provided they are run on hom $_{W}(F)$ if they are only using critical pairs, or use the $W$-GRevLex order otherwise.
For example, when computing elimination bases for positive dimensional systems, depending on the situation, it may be interesting to perform a two-steps computation, or to simply use one of the direct algorithms with the desired order. In the former case, FGLM cannot be used for the change of ordering. A common choice in this situation is the Gröbner walk. This algorithm is much more complex and to the best of our knowledge, does not have good complexity estimates. However, it involves computing successive Gröbner bases using algorithm $F_{4}$ or $F_{5}$ as a blackbox. As such, assigning weights to a polynomial system will yield similar improvements for the computing time.

### 3.6. Experiments

In this section, we present some applications where taking into account the weighted structure of the system yields speed-ups. For each system, we compare two strategies: the "standard" strategy consists of computing a Gröbner basis without considering the weighted structure; the "weighted" strategy is the strategy we described in Section 3.5. For all these examples, we use a more compact notation for degrees and weights, so that for example, $\left(2^{3}, 1\right)$ is equivalent to (2, 2, 2, 1).

### 3.6.1. Generic systems

First, we present some timings obtained with generic systems, in the complete intersection ( $m=n$ ), the positive-dimensional $(m<n)$ and the over-determined $(m>n)$ cases. In all cases, we fix a system of weights $W=\left(w_{1}, \ldots, w_{n}\right)$ and a system of $W$-degrees $D=\left(d_{1}, \ldots, d_{m}\right)$, and we pick at random $m$ polynomials $\left(f_{i}\right)_{i=1 \ldots m}$, such that for any $i \in\{1, \ldots, m\}, f_{i}$ has dense support in the set of monomials with $W$-degree less than or equal to $d_{i}$.
For zero-dimensional regular sequences, we compute a lexicographical Gröbner basis using a two-steps strategy in Magma, with algorithm $\mathrm{F}_{4}$ as a first step (first block of lines in Table 3.2a) and algorithm FGLM for the change of ordering (Table 3.2b).
For over-determined systems, we compute a Gröbner basis for the GRevLex ordering, using algorithm $\mathrm{F}_{4}$ from Magma (second block of lines in Table 3.2a).

For positive-dimensional systems, we compute a basis for an elimination order, using a two-steps strategy with $\mathrm{FGb}^{2}$ : first we compute a GRevLex basis with algorithm $\mathrm{F}_{4}$ (third block of lines in Table 3.2a), and then we compute a basis for the wanted elimination order, again with $\mathrm{F}_{4}$ (Table 3.2c). In this table, the second column describes what variables we eliminate: for example, 3 means that we eliminate the first 3 variables, while $1 \rightarrow 3$ means that we first eliminate the first variable, then the next 2 variables, again resulting in a basis eliminating the first 3 variables.

For algorithm $F_{4}$ with the GRevLex ordering, the behavior we observe is coherent with the previous complexity studies: we observe some speed-ups when taking into account the weighted structure of the system, and these speed-ups seem to increase with the weights. However, the speed-ups cannot be expected to match rigorously the ones predicted by the complexity bounds, because the systems are usually not regular for the standard strategy. Experiments also confirm that it is more effective to order the variables with highest weight first.

For the lexicographical ordering with FGLM, we also observe some speed-ups when applying the weights (we will observe this behavior again in Section 3.6.2). These differences are not explained by the theoretical complexity bounds, since both ideals have the same degree in each case. However, it appears that the slower FGLM runs are those where the FGLM matrix is denser, and that this difference in density seems to match quantitatively the speed-ups we observe.

Finally, for elimination bases, the results are similar to what we observed with the GREvLEX ordering: when possible, one should take into account the weights, and order the variables such that the smallest weights are also the smallest variables. However, when eliminating variables, the largest variables need to be the ones that should be eliminated. If the variables need to be ordered such that those with the smallest weights are first, in most cases, taking into account the weighted structure is still profitable. However, if the smallest weight is on the largest variable and there is only one such variable, this is no longer true (see for example the second line in Table 3.2 c ). Experiments suggest that these systems naturally possess a good weighted structure for the weights $(1, \ldots, 1)$ : their construction ensures that every such polynomial of total degree $d$ will have a large homogeneous component at degree $d / 2$, and the higher degree components will be small, and divisible by large powers of $X_{1}$. On the other hand, with weights $(1,2, \ldots, 2)$, the same polynomial will have a large $W$-homogeneous component at $W$-degree $d$, overall leading to reductions at higher degree (an example is given in Table 3.3).

We conclude this section with timings illustrating the consequences of the estimates of the degree of regularity of a system, depending on the order of the variables (Section 3.3.2). For this purpose, we generate a generic system of $W$-degree ( $60^{4}$ ) with weights ( $20,5,5,1$ ). Then we compute a $W$-GRevLex Gröbner basis for the orders $X_{1}>\cdots>X_{4}$ (smallest weights last) and for the reverse order $X_{4}<\cdots<X_{1}$. We give the degree of regularity, the value predicted by the previous bound (3.9), by the new bound (3.10) and by the conjectured bound (3.12), as well

[^1]Table 3.2: Benchmarks with Magma for generic systems

| Parameters | Without <br> weights (s) | With <br> weights (s) | Speed-up |
| :--- | ---: | ---: | ---: |
| $n=8, W=\left(2^{6}, 1^{2}\right), D=\left(4^{8}\right)$ | 8.0 | 2.5 | 3.2 |
| $n=9, W=\left(2^{7}, 1^{2}\right), D=\left(4^{9}\right)$ | 101.2 | 12.5 | 8.1 |
| $n=7, W=\left(2^{5}, 1^{2}\right), D=\left(8^{15}\right)$ | 31.6 | 7.5 | 4.2 |
| $n=7, W=\left(2^{5}, 1^{2}\right), D=\left(8^{14}\right)$ | 29.0 | 9.4 | 3.1 |
| $n=7, W=\left(2^{5}, 1^{2}\right), D=\left(8^{13}\right)$ | 40.0 | 12.0 | 3.3 |
| $n=5, m=4, W=\left(2^{4}, 1\right), D=\left(8^{4}\right)$ | 2.6 | 0.2 | 13.0 |
| $n=5, m=4, W=\left(1,2^{4}\right), D=\left(8^{4}\right)$ | 2.5 | 0.3 | 8.3 |
| $n=5, m=4, W=\left(1^{3}, 2^{2}\right), D=\left(4^{4}\right)$ | 23.6 | 0.0 | 2360.0 |
| $n=5, m=4, W=\left(2^{2}, 1^{3}\right), D=\left(4^{4}\right)$ | 407.5 | 0.0 | 40750.0 |

(a) Benchmarks for the $\mathrm{F}_{4}$ algorithm for the GRevLex ordering

| Parameters | Degree | Without <br> weights (s) | With <br> weights (s) | Speed-up |
| :--- | :---: | :---: | :---: | :---: |
| $n=8, W=\left(2^{6}, 1^{2}\right), D=\left(4^{8}\right)$ | 1024.0 | 500.4 | 495.0 | 1.0 |
| $n=9, W=\left(2^{7}, 1^{2}\right), D=\left(4^{9}\right)$ | 2048.0 | 11995.8 | 7462.1 | 1.6 |

(b) Benchmarks for the FGLM algorithm (lexicographical ordering)

| Parameters | Elim. vars. | Without <br> weights (s) | With <br> weights (s) | Speed-up |
| :--- | :--- | :---: | :---: | :---: |
| $n=5, m=4, W=\left(2^{4}, 1\right), D=\left(8^{4}\right)$ | 1 | 120.3 | 12.0 | 10.0 |
| $n=5, m=4, W=\left(1,2^{4}\right), D=\left(8^{4}\right)$ | 1 | 27.6 | 30.4 | 0.9 |
| $n=5, m=4, W=\left(1^{3}, 2^{2}\right), D=\left(4^{4}\right)$ | 2 | 146.3 | 6.9 | 21.2 |
| $n=5, m=4, W=\left(1^{3}, 2^{2}\right), D=\left(4^{4}\right)$ | $1 \longrightarrow 2$ | 162.0 | 3.3 | 49.1 |
| $n=5, m=4, W=\left(2^{2}, 1^{3}\right), D=\left(4^{4}\right)$ | 1 | $>750$ | 0.1 | $>7500$ |
| $n=5, m=4, W=\left(2^{2}, 1^{3}\right), D=\left(4^{4}\right)$ | $1 \longrightarrow 2$ | NA | 0.1 | NA |
| $n=5, m=4, W=\left(2^{2}, 1^{3}\right), D=\left(4^{4}\right)$ | $1 \longrightarrow 2 \longrightarrow 3$ | NA | 7.9 | NA |

(c) Benchmarks for the $F_{4}$ algorithm for elimination

Table 3.3: Size of the $W$-homogeneous components for a generic polynomial with $W_{0}$-degree 4 for $W_{0}=(1,2,2,2)$

| $W$-degree | $W=(1,2,2,2)$ | $W=(1,1,2,2)$ | $W=(1,1,1,1)$ |
| :---: | ---: | ---: | ---: |
| o | 1 | 1 | 1 |
| 1 | 1 | 2 | 4 |
| 2 | 4 | 5 | 10 |
| 3 | 4 | 6 | 4 |
| 4 | 10 | 6 | 1 |

Table 3.4: Impact of the order of the variables on the degree of regularity and the computation times (generic weighted homogeneous system with $W$-degree $\left(60^{4}\right)$ w.r.t. $W=(20,5,5,1))$

| Order | $d_{\text {reg }}$ | Macaulay's <br> bound (3.9) | Bound (3.15) | Conjectured <br> bound (3.12) | $\mathrm{F}_{5}$ time |
| :---: | ---: | ---: | ---: | ---: | ---: |
| $X_{1}>X_{2}>X_{3}>X_{4}$ | 210 | 229 | 210 | 210 | 101.9 |
| $X_{4}>X_{3}>X_{2}>X_{1}$ | 220 | 229 | 229 | 220 | 255.5 |

as the timings. This experiment was run using algorithm $F_{5}$ from the FGb library, the results are in Table 3.4.

### 3.6.2. Discrete Logarithm Problem

Taking advantage of a weighted homogeneous structure has allowed the authors of the article $[\mathrm{Fau}+13]$ to obtain significant speed-ups for solving a system arising from the DLP on Edwards elliptic curves [Gauo9]. They observed that the system of equations they had to solve has symmetries, and rewrote it in terms of the invariants of the symmetry group. For a system in $n$ equations, these invariants are

$$
\begin{aligned}
E_{1} & =e_{1}\left(X_{1}^{2}, \ldots, X_{n}^{2}\right) \\
E_{2} & =e_{2}\left(X_{1}^{2}, \ldots, X_{n}^{2}\right) \\
& \vdots \\
E_{n-1} & =e_{n-1}\left(X_{1}^{2}, \ldots, X_{n}^{2}\right) \\
E_{n} & =e_{n}\left(X_{1}, \ldots, X_{n}\right) .
\end{aligned}
$$

The system they obtained is sparser, but does not have a good homogeneous structure. In particular, the highest total degree components of the system do not form a regular sequence, and in practice, a Gröbner basis computation will follow many degree falls.

However, the system had a weighted homogeneous structure for the weights ( $2, \ldots, 2,1$ ) (only $E_{n}$ has weight 1 ), with respective $W$-degree $\left(2^{n}, \ldots, 2^{n}\right)$. The highest $W$-degree components forming a sequence in simultaneous Noether position with respect to the order $E_{1}>E_{2}>\cdots>$ $E_{n}$, one could compute a Gröbner basis without any $W$-degree fall, with complexity bounded by the estimates (3.14) and (3.17).

This system also illustrates the role of the choice of the system of weights: recall that we want to avoid degree falls, that is reductions to zero of the highest $W$-degree components of the polynomials.

To compare the behavior of the algorithm with the total degree and the appropriate $W$-degree, we plotted the degree for each reduction step in a run of algorithm $F_{4}$ in Figure 3.6.

On the blue curve, any horizontal or downwards segment is a step where the degree is less than or equal to the degree of the previous step: a degree fall. Indeed, when using the total degree, the system was not regular in the affine sense, and degree falls are to be expected. On the other hand, the red curve shows the progress of the algorithm when using the weights. It is an example of a regular run: there is no degree fall.

Table 3.5: Benchmarks with FGb and Magma for DLP systems

| System | $\operatorname{deg}(I)$ | $\mathrm{F}_{5} \mathrm{w}(\mathrm{s})$ | $\mathrm{F}_{5} \mathrm{std}(\mathrm{s})$ | Speed-up for $\mathrm{F}_{5}$ | $\begin{array}{r} \text { FGLM } \\ \mathrm{w}(\mathrm{~s}) \end{array}$ | FGLM <br> std (s) | Speed-up <br> for FGLM |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| DLP Edwards: $n=4$, $W=\left(2^{3}, 1\right), D=\left(8^{4}\right)$ | 512 | 0.1 | 0.1 | 1.0 | 0.1 | 0.1 | 1.0 |
| DLP Edwards: $n=5$, $W=\left(2^{4}, 1\right), D=\left(16^{5}\right)$ | 65536 | 935.4 | 6461.2 | 6.9 | 2164.4 | 6935.6 | 3.2 |
| (a) Benchmarks with FGb |  |  |  |  |  |  |  |
| System | $\operatorname{deg}(I)$ | $\mathrm{F}_{4} \mathrm{w}(\mathrm{s})$ | $\mathrm{F}_{4} \mathrm{std}$ (s) | $\begin{aligned} & \text { Speed-up } \\ & \text { for } F_{4} \end{aligned}$ | FGLM <br> w (s) | FGLM <br> std (s) | Speed-up for FGLM |
| DLP Edwards: $n=4$, $W=\left(2^{3}, 1\right), D=\left(8^{4}\right)$ | 512 | 1 | 1 | 1.0 | 1 | 27 | 27 |
| DLP Edwards: $n=5$, $W=\left(2^{4}, 1\right), D=\left(16^{5}\right)$ | 65536 | 6044 | 56105 | 9.3 | $\infty$ | $\infty$ | NA |

(b) Benchmarks with Magma

In this example, the final $W$-degree reached with the weights is exactly twice the final degree reached without the weights, so that the matrices built and reduced have approximately the same size. But when there are degree falls, evaluating the highest degree reached is not sufficient for counting the number of steps.

### 3.6.3. Polynomial inversion

The polynomial inversion problem consists in finding polynomial relations between polynomials. More precisely, given a system of polynomial equations

$$
\left\{\begin{aligned}
f_{1}\left(X_{1}, \ldots, X_{n}\right) & =0 \\
f_{2}\left(X_{1}, \ldots, X_{n}\right) & =0 \\
& \vdots \\
f_{m}\left(X_{1}, \ldots, X_{n}\right) & =0
\end{aligned}\right.
$$

we want to compute all the relations of the form

$$
g_{i}\left(f_{1}, \ldots, f_{r}\right)=0
$$

One can compute these relations with Gröbner bases by computing an elimination ideal: consider the ideal generated by the polynomials

$$
\begin{aligned}
T_{1} & -f_{1}\left(X_{1}, \ldots, X_{n}\right) \\
T_{2} & -f_{1}\left(X_{1}, \ldots, X_{n}\right) \\
& \vdots \\
T_{m} & -f_{m}\left(X_{1}, \ldots, X_{n}\right)
\end{aligned}
$$



Figure 3.8: Degree at each reduction step in a run of Algorithm $\mathrm{F}_{5}$ on a DLP system ( $n=5$, $W=\left(2^{3}, 1\right), D=\left(16^{5}\right)$
in $R:=\mathbb{K}\left[X_{1}, \ldots, X_{n}, T_{1}, \ldots, T_{m}\right]$. Order $R$ with an elimination order on the variables $X_{1}, \ldots, X_{n}$, recall that it is an order such that

$$
\begin{align*}
& m_{X}\left(X_{1}, \ldots, X_{n}\right) m_{T}\left(T_{1}, \ldots, T_{m}\right)<_{\text {elim }} m_{X}^{\prime}\left(X_{1}, \ldots, X_{n}\right) m_{T}^{\prime}\left(T_{1}, \ldots, T_{m}\right)  \tag{3.18}\\
& \\
& \Longleftrightarrow\left\{\begin{array}{l}
m_{X}<_{X} m_{X}^{\prime} \\
\text { or } \\
m_{X}=m_{X}^{\prime} \text { and } m_{T}<_{T} m_{T}^{\prime}
\end{array}\right.
\end{align*}
$$

for some monomial orders $<_{X}$ and $<_{T}$. The usual choice is the block-GRevLex order Eima ${ }_{n}$.
This problem can benefit from being given a weighted structure [Tra96, sec. 6.1]. For any $i \in\{1, \ldots, m\}$, let $d_{i}$ be the degree of $f_{i}$. By setting the weight of $T_{i}$ to be $d_{i}$, the monomial $T_{i}$ becomes part of the highest $W$-degree component of $T_{i}-f_{i}\left(X_{1}, \ldots, X_{n}\right)$, giving this equation a weighted homogeneous structure.
More precisely:
Proposition 3.51. Let $f_{1}, \ldots, f_{m}$ be a system of polynomials with respective degree $d_{1}, \ldots, d_{m}$ in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$. Consider the algebra $R:=\mathbb{K}\left[X_{1}, \ldots, X_{n}, T_{1}, \ldots, T_{m}\right]$, graded with the weights $W=\left(1, \ldots, 1, d_{1}, \ldots, d_{m}\right)$, and consider the system $F=\left(T_{1}-f_{1}(\mathbf{X}), \ldots, T_{m}-f_{m}(\mathbf{X})\right)$ in $R$. Then the system $F^{h}$ formed with the highest $W$-degree components of $F$ is in Noether position with respect to the variables $T_{1}, \ldots, T_{m}$, and in particular it forms a regular sequence.

Proof. By the choice of the weights, the system $F^{h}$ is defined by

$$
F^{h}=\left(T_{1}-f_{1}^{h}(\mathbf{X}), \ldots, T_{m}-f_{m}^{h}(\mathbf{X})\right),
$$

where for any $i \in\{1, \ldots, m\}, f_{i}^{h}$ is the highest degree component of $f_{i}$. As a consequence, by the characterization 4 of the Noether position, the system $F^{h}$ is indeed in Noether position with respect to the variables $T_{1}, \ldots, T_{m}$.

In Tables 3.6, we present timings for a few systems with this kind of problem:

- group invariants [Stuo8]: given a group, compute its fundamental invariants, and then the relations between these invariants. Since these examples can lead to very long computations, in some cases, we only compute the relations between the $k$ first invariants;
- monomials: given $m$ monomials of degree $d$ in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$, compute the relations between them;
- matrix minors: given a $p \times q$ matrix of linear forms in $n$ indeterminates, compute all its minors of rank $r$ as polynomials in the $X_{i, j}$ 's, and compute the relations between them.

In each case, we compute an elimination basis using a two-steps strategy: first we compute a GRevLex basis (Table 3.6a), then we compute the elimination basis (Table 3.6b). In Table 3.6c, we show some timings for the computation of the elimination basis directly from the input system. All these experiments were run using algorithm $\mathrm{F}_{4}$ from Magma.

Table 3.6: Benchmarks with Magma on some polynomial inversion systems

| System | Without <br> weights (s) | With <br> weights (s) | Speed-up |
| :--- | ---: | ---: | ---: |
| Cyclic invariants, $n=4$ | 4.2 | 0.0 | 140.0 |
| Cyclic invariants, $n=5, k=12$ | 2612.6 | 54.7 | 47.8 |
| Cyclic invariants, $n=5$ | $75^{\text {a }}$ | 390.7 | NA |
| Cyclic invariants, $n=6, k=14$ | 32987.6 | 2787.7 | 11.8 |
| Cyclic invariants, $n=6, k=15$ | $>280000^{\text {a }}$ | 14535.4 | NA |
| Dihedral invariants, $n=5$ | $70000^{\text {a }}$ | 6.3 | NA |
| Generic monomials, $d=2, n=24, m=48$ | 216.1 | 0.2 | 1350.6 |
| Generic monomials, $d=2, n=25, m=50$ | 14034.7 | 0.1 | 116955.8 |
| Generic monomials, $d=2, n=26, m=52$ | 14630.6 | 0.2 | 66502.7 |
| Generic monomials, $d=2, n=27, m=54$ | 8887.6 | 0.2 | 55547.5 |
| Generic monomials, $d=3, n=11, m=22$ | 370.9 | 0.1 | 6181.7 |
| Generic monomials, $d=3, n=12, m=24$ | 4485.0 | 0.2 | 26382.4 |
| Matrix minors, $n=5,7 \times 7, r=3$ | 125.7 | 93.3 | 1.3 |
| Matrix minors, $n=6,7 \times 7, r=3$ | 1941.0 | 1029.1 | 1.9 |
| Matrix minors, $n=6,8 \times 8, r=3$ | 4115.8 | 2295.8 | 1.8 |
| Matrix minors, $n=4,6 \times 6, r=5$ | 612.6 | 159.2 | 3.8 |
| Matrix minors, $n=4,7 \times 7, r=6$ | 8043.3 | 2126.9 | 3.8 |
| Matrix minors, $n=4,7 \times 10, r=7$ | 69386.1 | 43910.1 | 1.6 |

a. Memory usage was over 120 GB

## (a) First step ( $\mathrm{F}_{4}$ for the GRevLex order)

| System | Without <br> weights (s) | With <br> weights (s) | Speed-up |
| :--- | ---: | ---: | ---: |
| Cyclic invariants, $n=4$ | 7.0 | 0.1 | 70.0 |
| Cyclic invariants, $n=5, k=12$ | 1683.2 | 70.7 | 23.8 |
| Cyclic invariants, $n=5$ | NA | 382.5 | NA |
| Cyclic invariants, $n=6, k=14$ | 9236.4 | 1456.0 | 6.3 |
| Cyclic invariants, $n=6, k=15$ | NA | 7179.7 | NA |
| Dihedral invariants, $n=5$ | NA | 20.3 | NA |
| Generic monomials, $d=2, n=24, m=48$ | 250.3 | 117.4 | 2.1 |
| Generic monomials, $d=2, n=25, m=50$ | 13471.2 | 15932.9 | 0.8 |
| Generic monomials, $d=2, n=26, m=52$ | 17599.5 | 8054.2 | 2.2 |
| Generic monomials, $d=2, n=27, m=54$ | 9681.0 | 3605.6 | 2.7 |
| Generic monomials, $d=3, n=11, m=22$ | 624.5 | 199.9 | 3.1 |
| Generic monomials, $d=3, n=12, m=24$ | 9751.6 | 3060.1 | 3.2 |
| Matrix minors, $n=5,7 \times 7, r=3$ | 52.6 | 66.6 | 0.8 |
| Matrix minors, $n=6,7 \times 7, r=3$ | 556.5 | 779.1 | 0.7 |
| Matrix minors, $n=6,8 \times 8, r=3$ | 1257.9 | 1714.0 | 0.7 |
| Matrix minors, $n=4,6 \times 6, r=5$ | 262.7 | 328.1 | 0.8 |
| Matrix minors, $n=4,7 \times 7, r=6$ | 2872.2 | 4299.8 | 0.7 |
| Matrix minors, $n=4,7 \times 10, r=7$ | 4728.4 | 5485.8 | 0.9 |

(b) Second step ( $\mathrm{F}_{4}$ for an elimination order)

| System | Without <br> weights (s) | With <br> weights (s) | Speed-up |
| :--- | ---: | ---: | :---: |
| Cyclic invariants, $n=4$ | 4.0 | 0.3 | 13.3 |
| Cyclic invariants, $n=5, k=12$ | 2705.8 | 73.4 | 36.9 |
| Cyclic invariants, $n=5$ | $90000^{\mathrm{b}}$ | 370.0 | $>243$ |
| Cyclic invariants, $n=6, k=14$ | 35922.4 | 2256.2 | 15.9 |
| Cyclic invariants, $n=6, k=15$ | $>300000^{\mathrm{b}}$ | 7426.7 | $>40$ |
| Dihedral invariants, $n=5$ | 40000 | 18.5 | $>2162$ |
| Generic monomials, $d=2, n=24, m=48$ | 216.5 | 110.9 | 2.0 |
| Generic monomials, $d=2, n=25, m=50$ | 31135.2 | 16352.2 | 1.9 |
| Generic monomials, $d=2, n=26, m=52$ | 14919.2 | 8142.8 | 1.8 |
| Generic monomials, $d=2, n=27, m=54$ | 5645.8 | 4619.0 | 1.2 |
| Generic monomials, $d=3, n=11, m=22$ | 370.1 | 193.1 | 1.9 |
| Generic monomials, $d=3, n=12, m=24$ | 4527.2 | 2904.6 | 1.6 |
| Matrix minors, $n=7,7 \times 7, r=3$ | 41220.0 | 26340.0 | 1.6 |
| Matrix minors, $n=7,8 \times 8, r=3$ | 48000.0 | 18060.0 | 2.7 |
| Matrix minors, $n=8,8 \times 8, r=3$ | 711690.0 | 390235.0 | 1.8 |
| Matrix minors, $n=4,6 \times 6, r=5$ | 613.9 | 325.4 | 1.9 |
| Matrix minors, $n=4,7 \times 7, r=6$ | 8059.4 | 3955.5 | 2.0 |
| Matrix minors, $n=4,7 \times 10, r=7$ | 71067.8 | 32721.5 | 2.2 |

b. Memory usage was over 120 GB .
(c) Direct strategy

As in the previous section, we plot in Figure 3.9 the degree of each reduction step on a run of Algorithm $\mathrm{F}_{4}$ for the GRevLex step in finding the relations between a random set of monomials. Once again, using the weighted degree instead of the total degree makes the system regular in the affine sense, and there is no degree falls with the weighted degree. In this case, this newly found regularity also makes the final degree reached much smaller, which justifies the ridiculous speed-ups that we observed on this class of systems (second block in Table 3.6a).


Figure 3.9: Degree at each reduction step in a run of Algorithm $F_{4}$ on a polynomial inversion systems (relations between 50 monomials of degree 2 in 25 variables)

## Chapter 4

# Real root classification for determinants Application to contrast optimization 

The results presented in this chapter are extracted from a joint work with Bernard Bonnard, JeanCharles Faugère, Alain facquemard and Mohab Safey El Din. They have been published in [Bon+16].

In this chapter, we consider the algorithmic problem of classifying the singularities of real parameterized determinantal systems. This problem has direct applications, for example in contrast optimization for Nuclear Magnetic Resonance imagery. In its most general setting, this example is out of reach of existing implementations of algorithms for real roots classification problems. We describe a computational strategy refining these classic tools in order to take advantage of the structure of determinantal systems. The new algorithm is able to find the classification for the general case of the contrast optimization problem.

### 4.1. Introduction

Real roots classification and determinantal systems In this chapter, we consider a real roots classification problem for determinantal ideals. A real roots classification problem can be stated as follows: let $V$ be an algebraic variety in $\mathbb{C}^{n} \times \mathbb{C}^{t}$, where the $n$ first coordinates are variables and the $t$ last are parameters, and let $B$ be a real semi-algebraic set in $\mathbb{R}^{n} \times \mathbb{R}^{t}$. Let $\pi: \mathbb{C}^{n} \times \mathbb{C}^{t} \longrightarrow \mathbb{C}^{t}$. The goal of roots classification is to subdivide $\mathbb{R}^{t}$ into areas where the cardinality of the fibers of the restriction of $\pi$ to $V \cap B \subset \mathbb{R}^{n+t}$ is constant. More precisely, we want to find connected open subsets $C_{1}, \ldots, C_{s} \subset \mathbb{R}^{t}$ for the Euclidean topology such that:

1. $C_{1} \cup \cdots \cup C_{s}$ is dense in $\mathbb{R}^{t}$
2. For any $i \in\{1, \ldots, s\}$, there exists $c_{i} \in \mathbb{N}$ such that for any $\mathbf{g} \in C_{i}, \#\left(V \cap B \cap \pi^{-1}(\mathbf{g})\right)=c_{i}$.

Obviously, this can only be done if the generic fiber of $\pi$ on $V \cap B$ is finite.
Here, the variety $V$ shall be a determinantal variety, as defined in Section 1.5. To fix the notations, let $M$ be a $k \times k$ matrix whose coefficients are polynomials in $\mathbb{Q}\left[X_{1}, \ldots, X_{n}, G_{1}, \ldots, G_{t}\right]$. As above, let $\pi: \mathbb{C}^{n} \times \mathbb{C}^{t} \longrightarrow \mathbb{C}^{t}$.

Fix $r \in\{1, \ldots, k-1\}$. Let $V_{r}$ be the $r$ 'th determinantal variety associated with $M$ (that is the set of points of $\mathbb{C}^{n+t}$ at which $M$ has rank at most $r$ ), and $V$ the union of the singular locus of $V_{r}$ and of the critical points of $\pi$ restricted to $V_{r}$.


Figure 4.1: Contrast optimization: both pictures were taken using the same setup of two tubes, the inner one containing deoxygenated blood and the outer one containing oxygenated blood. The left-hand picture is the picture obtained with no contrast optimization; the right-hand picture is contrast-optimized by saturating the inner sample.

Recall from Proposition 1.92 that generically, this variety is $V_{r-1}$, and it is equidimensional with dimension $n+t-(k-r+1)^{2}$. Further assume that

$$
n=(k-r+1)^{2}
$$

so that generically, $V$ is $t$-equidimensional. In this case, generically, there exists a Zariski-open set $O \subset \mathbb{C}^{t}$ such that $V \cap \pi^{-1}(\mathrm{~g})$ is a non-empty finite set for any $\mathbf{g} \in O$.

As before, let $B$ be a real semi-algebraic set in $\mathbb{R}^{n} \times \mathbb{R}^{t}$, and assume that:

1. $B$ has non-empty interior (for the euclidean topology)
2. the projection $\pi$ restricted to $V \cap B$ is proper.

We want to design an algorithm for the roots classification problem specialized to $V \cap B$.

Application to contrast optimization This problem arose from a control theory study of the dynamics of contrast optimization in Nuclear Magnetic Resonance (MNR) imaging.

Nuclear Magnetic Resonance (NMR) is a powerful tool in medical imagery. In order to distinguish two biological matters on a picture, it is required to optimize the contrast between the two matters. Because of its importance in medical sciences, this contrast imaging problem has received a lot of attention. The pioneering work of [Bon+12] has established geometric optimal control techniques as a major tool for designing optimal control strategies for the problem of improving the contrast.

These strategies depend on the biological matters under study. In NMR imagery the main physical parameters involved are the longitudinal and transversal relaxation times of each matter. This approach is formalized in [Bon+13], and it leads to a real root classification problem for a determinantal system, where the size of the matrix is $k=4$ and the number of variables is $n=4$.

The number of parameters is initially 4: the longitudinal and the transversal relaxation times, for 2 biological matters. Because the system is homogeneous and physical constraints ensure that the parameters are nonzero, we may set one of them to 1 , reducing the problem to $t=3$. An
important particular case is when one of the matters is water; then the corresponding relaxation times are equal, and the number of parameters is $t=2$.
We also mention that the optimal control problems in $[\mathrm{Chy}+15]$ lead to algebraic classification problems with similar structures.

State of the art The modeling through an optimal control problem is introduced in [Bon +13 ]. The so-called Bloch modeling and saturation method for tackling this problem is developed therein.
In [Bon +13 ], four experimental important cases are studied (all parameters of the classification problem are fixed). Among other properties, it has been observed there that the number of singularities is constant when water is involved. This led to the following questions:

1. Is this number of singularities preserved for any choice of second matter?
2. If not, how many different classes of pairs of matters can we distinguish through the analysis of those singularities?

Answering these questions leads to the real root classification problem described above. Symbolic computation techniques are good candidates to solve them.
Properties of Cylindrical Algebraic Decomposition (CAD) adapted to a given polynomial family allow to solve real root classification problems. Hence the CAD algorithm [Col ${ }_{75}$ ] can be used in our context. However, the complexity of computing a CAD is doubly exponential in the number of variables ([BDo7; DH88]); its implementations are usually limited to non-trivial problems involving 4 variables and cannot tackle our application.
The complexity of computing a CAD can be much improved when taking into account equational constraints (see e.g. [McC99]). In the context of real root classification problems, this leads us to take advantage of the presence of equations to compute closed sets in the parameter space ( $\mathbb{R}^{t}$ using our notation) containing the boundaries of the regions $C_{1}, \ldots, C_{\ell}$, hence substituting the recursive (doubly exponential) projection steps of CAD with more involved projection techniques. In the past ten years, several works have focused on this problem [LRo7; YHXor] using various computer algebra tools such as Gröbner bases, regular chains, etc. We also mention [ $\mathrm{RT}_{15}$ ] which uses evaluation/interpolation techniques to compute those closed sets in the parameters space.
While the implementation of [ YHXor$]$ is able to solve our classification problem for the case of water, none of the implementations were able to classify the singular locus of $V$ in the general case (the number of parameters is 3 ).
Our strategy is based on those of [LRo7] and [YHXo1], adapted to the determinantal case: by exploiting properties of sets defined by minors of matrices with polynomial entries, we can obtain a more regular behavior for the algorithms. As such, our strategy is less general than existing ones, but is able to compute the partition of the parameter space in the general case.
Such structures have been used for computing sample points in each connected component of the real trace of determinantal varieties [HNS15b; HNS15c; HNS16] or for solving linear matrix inequalities [ HNS 15 a ]. These works are based on dedicated strategies for computing critical loci of some projections restricted to determinantal varieties. Such computations are naturally related to real root classification problems and real quantifier elimination (see e.g. [HS12]).

Finally, the computations rely on Gröbner bases. Several works have focused on designing strategies improving the regularity of Gröbner basis computations on determinantal systems [FSS13], and in particular when computing singular and critical loci of determinantal varieties [FSS12; Spa14]. These strategies involve modelling the problem using incidence varieties instead of minors.

Main results The main results of this chapter are twofold:

- an algorithm solving the special real root classification problem described above and which exploits the determinantal structure of the input data arising in contrast imaging problems;
- its successful use for solving the challenging application to contrast problem in the general case.

We start by describing our algorithmic contribution. Recall that we are given a matrix, denoted by $M$, with polynomial entries. As in [HNS15a; HNS15b; HNS15c; HNS16], it is based on splitting computations according to the rank of $M$.

More precisely, in order to solve the real root classification problem, we need to identify where the number of real solutions inside $B$ of the determinantal system describing $V$ changes depending on the values of the parameters. In particular, this problem involves inequalities defining a semi-algebraic set with non-empty interior. In this context, we use standard tools from real geometry, such as Thom's first isotopy lemma, which reduce our classification problem to computing the singular points of $V$, the critical points of the projection of the parameter space restricted to $V$, and the intersection of $V$ with the boundary of the semi-algebraic set $B$.

This computation may be difficult because generically, the variety $V_{r}$ has singularities corresponding to points where $\operatorname{rank}(M)<r$ (Proposition 1.92). Hence, observe that the variety $V$ is naturally split according to the rank of $M$. This is the very basic idea on which our algorithm relies: we compute critical loci of the projection on the parameter space restricted to the variety $V$ by distinguishing those points at which $M$ has rank less than $r$ from those at which $M$ has rank exactly $r$.

Our algorithms need to compute projection of algebraic sets, which is done using elimination algorithms such as Gröbner bases or triangular sets for example. We have performed experiments for both these tools, using the package FGb ([Fau10]) in Maple and an implementation of $\mathrm{F}_{5}$ ([Fauo2]) for Gröbner bases, and using the package RegularChains ([LMXo5]) in Maple for triangular sets.

Regarding the contrast imaging problem, we illustrate the behaviour of our algorithm in the case of water $(t=2)$, giving the whole classification. Using Gröbner bases to perform the eliminations, the computation takes 10 s on an 2 GHz Intel Xeon CPU. The RealRootClassification command of the Maple RegularChain library needs 1600 s to find this classification.

We also ran our algorithm on the general case $(t=3)$. While none of the available implementations is able to tackle this classification problem directly, ours can find the polynomials separating the open sets $C_{i}$ within 4 h using FGb, or 2 min using $\mathrm{F}_{5}$, and the projection step of the CAD can be done in 4 h . We see similar speed-ups when using triangular sets to perform the elimination.

This illustrates how our dedicated algorithms take advantage of the special structure of the problem, to achieve speed-ups when compared with more general techniques.

We give an overview of the results at the end of the chapter ${ }^{1}$. The source code is also given in Appendix A.2.

Structure of the chapter In Section 4.2 , we present the mathematical background around NMR imagery and the contrast problem. Section 4.3 deals with the dedicated classification algorithm. Finally, in Section 4.4, we report on experimental results obtained when solving the application to the contrast imaging problem.

### 4.2. Modeling the dynamics of the contrast optimization

The model we describe below has been introduced in [Bon+12] in order to apply techniques from geometric optimal control theory to the control of the spin dynamics by NMR. Up to some normalization, each spin $1 / 2$ particle is governed by the Bloch equation

$$
\left\{\begin{array}{l}
\dot{x}=-\Gamma x+u_{y} z \\
\dot{y}=-\Gamma y-u_{x} z \\
\dot{z}=\gamma(1-z)+u_{x} y-u_{y} x
\end{array}\right.
$$

where the state variable $q=(x, y, z)$ represents the magnetization vector which must lie in the Bloch ball defined by $|q| \leq 1$, and the parameters $(\Gamma, \gamma)$ are related to the physical relaxation times. The parameters must also satisfy $2 \Gamma \geq \gamma>0$. The control $u=\left(u_{x}, u_{y}\right)$ represents the magnetic field whose magnitude is bounded by a maximum value $\mu$.

In the context of the contrast imaging problem, this leads to the simultaneous control of two non-interacting spins with different relaxation time parameters. The contrast by saturation method consists in bringing the magnetization vector of the first spin toward the center of the Bloch ball while maximizing the modulus of the magnetization vector of the other matter. The matter with a zero magnetization is black on the picture, while the other matter with a maximum modulus of the magnetization vector is bright.

Using the symmetry of revolution [Bon+13] which allows to eliminate one state variable for each matter, we obtain the system

$$
\left\{\begin{array}{r}
\dot{y_{1}}=-\Gamma_{1} y_{1}-u_{x} z_{1}  \tag{4.1}\\
\dot{z_{1}}=\gamma_{1}\left(1-z_{1}\right)+u_{x} y_{1} \\
\dot{y_{2}}=-\Gamma_{2} y_{1}-u_{x} z_{2} \\
\dot{z_{2}}=\gamma_{2}\left(1-z_{2}\right)+u_{x} y_{2}
\end{array} \quad|u| \leq \mu\right.
$$

and the optimal control problem is: starting from the equilibrium point $N=((0,1),(0,1))$, saturate the first spin, that is $q_{1}(T)=0$, where $T$ is the transfer time while maximizing $\left|q_{2}(T)\right|^{2}$, where $\left|q_{2}(T)\right|$ represents the final contrast. It is a standard Mayer problem in optimal control. It has been studied in [Bon+12] through the analysis of the Hamiltonian dynamics given by the Pontryagin Maximum Principle [Pon+62]. We summarize this analysis below.

[^2]Writing (4.1) as $\dot{q}=F(q)+u G(q), \quad|u| \leq \mu$, the optimality conditions associated with the Maximum Principle lead us to construct the optimal solution as a concatenation of bang-arcs and singular arcs. Bang-arcs are defined by setting the control to $u= \pm \mu$, and singular arcs are solutions $u_{s}$ of $X_{e}=F+u_{s} G$. In the latter case, the control $u_{s}$ is the rational fraction $-D^{\prime} / D$ with

$$
\begin{aligned}
D & =\operatorname{det}(F, G,[G, F],[[G, F], G]) \\
D^{\prime} & =\operatorname{det}(F, G,[G, F],[[G, F], F]),
\end{aligned}
$$

where [, ] denotes the Lie bracket of vector fields. Explicitly, with $d_{i}=\gamma_{i}-\Gamma_{i}(i \in\{1,2\})$,

$$
D=\operatorname{det}\left[\begin{array}{cccc}
-\Gamma_{1} y_{1} & -z_{1}-1 & d_{1} z_{1}-\Gamma_{1} & 2 d_{1} y_{1} \\
-\gamma_{1} z_{1} & y_{1} & d_{1} y_{1} & -2 d_{1} z_{1}+\Gamma_{1}-d_{1} \\
-\Gamma_{2} y_{2} & -z_{2}-1 & d_{2} z_{2}-\Gamma_{2} & 2 d_{2} y_{2} \\
-\gamma_{2} z_{2} & y_{2} & d_{2} y_{2} & -2 d_{2} z_{2}+\Gamma_{2}-d_{2}
\end{array}\right]
$$

The localization of the singularities of $\{D=0\}$ inside the Bloch ball is important to understand the geometry of the hypersurface, as well as the dynamics of the vector field $X_{e}$ which is closely linked to the presence of such singularities.

### 4.3. Algorithm

### 4.3.1. Classification strategy

We consider the polynomial algebra $\mathbb{Q}[\mathbf{X}, \mathbf{G}]$ with variables $\mathbf{X}=\left(X_{1}, \ldots, X_{n}\right)$ and parameters $\mathbf{G}=\left(G_{1}, \ldots, G_{t}\right)$. Let $F$ and $H$ be families of polynomials in $\mathbb{Q}[\mathbf{X}, \mathrm{G}]$. Let $V_{\mathbb{R}}=V_{\mathbb{R}}(F), V=V_{\mathbb{C}}(F)$ be the set of zeroes of $F$ in $\mathbb{R}^{n+t}$ and in $\mathbb{C}^{n+t}$ respectively. Let $B$ be the closed semi-algebraic set defined by $H$ :

$$
B=\left\{(\mathbf{x}, \mathbf{g}) \in \mathbb{R}^{n+t} \mid \forall h \in H, h(\mathbf{x}, \mathbf{g}) \leq 0\right\}
$$

and let $B_{0}=\bigcup_{h \in H} V_{\mathbb{C}}(h)$. Let $\pi: \mathbb{C}^{n+t} \longrightarrow \mathbb{C}^{t}$ be the projection onto the affine space with coordinates $G$. Let $\operatorname{sing}(V)$ be the singular locus of $V, \operatorname{crit}(\pi, V)$ be the set of critical points of $\pi$ restricted to $V$, and $K(\pi, V)=\pi(\operatorname{sing}(V) \cup \operatorname{crit}(\pi, V)) \cap \mathbb{R}^{t}$.

Given a subset $A$ of a real or complex affine space, $\bar{A}$ and $\partial A$ are used to denote the closure and the boundary of $A$ for the Euclidean topology respectively.

Assume that the following hypotheses are satisfied:
$\mathcal{H}$. There exists a nonempty Zariski-open subset $O_{1}$ of $\mathbb{C}^{t}$ such that for all $\mathbf{g} \in O_{1}$, the fiber $V \cap \pi^{-1}(\mathrm{~g})$ is a nonempty finite subset of $\mathbb{C}^{n+t}$;
$\mathcal{H}$ 2. The restriction of the projection $\pi$ to $B$ is proper (Definition 1.81);
$\mathcal{H}$ 3. The intersection $V \cap B_{0}$ has dimension at most $t-1$ in $\mathbb{C}^{n+t}$;
$\mathcal{H}$ 4. The variety $V$ is equidimensional of dimension $t$.

We want to find a non-zero polynomial $P \in \mathbb{Q}[\mathrm{G}]$ such that, on each connected component $U$ of $\mathbb{R}^{t} \backslash V_{\mathbb{R}}(P)$, for $\mathbf{g} \in U$, the cardinality of $V \cap B \cap \pi^{-1}(\mathbf{g})$ does not depend on $\mathbf{g}$.

In Lemma 4.1, we describe a well-known strategy for computing these objects (see for example [LRo7; YHXo1]).

Lemma 4.1. Let $F$ and $H$ be polynomial systems satisfying hypotheses $\mathcal{H} 1, \mathcal{H} 2, \mathcal{H} 3$ and $\mathcal{H} 4$. Let $C_{B}=\pi\left(V \cap B_{0}\right), U$ a non-empty connected open subset of $\mathbb{R}^{t}$ which does not meet $C_{B} \cup K(\pi, V)$, and $\mathbf{g} \in U$. Then $V \cap \pi^{-1}(\mathbf{g})$ is finite, and for any $\mathbf{g}^{\prime} \in U$, $\#\left(V \cap \pi^{-1}\left(\mathbf{g}^{\prime}\right)\right)=\#\left(V \cap \pi^{-1}(\mathbf{g})\right)$.

Proof. We will construct a Whitney stratification of $V \cap B$ ([BCR98, Def. 9.7.1]) with certain properties. First note that since $V$ is $t$-equidimensional by hypothesis $\mathcal{H} 4, V$ has real dimension at most $t$ (Proposition 1.6o). Let $\mathcal{S}_{=t}$ be the intersection of the points where $V_{\mathbb{R}}$ has local dimension $t$ and of the interior of $B$. There exists a Whitney stratification $\left(\mathcal{S}_{i}\right)$ of the semialgebraic set $V \cap B$ such that $\mathcal{S}_{=t}$ is the union of strata of dimension $t$ ([BCR98, Th. 9.7.11]). Let $\mathcal{S}_{<t}$ be the union of the other strata, they all have real dimension less than $t$. By construction, this is a semi-algebraic set which is the union of $(V \cap \partial B) \subset\left(V \cap B_{0}\right)$ and of the singular locus of $V \cap B$, and it has dimension less than $t$ (using hypothesis $\mathcal{H} 3$ for $V \cap B_{0}$ ). Its image through $\pi$ has dimension less than $t$, and so it has codimension at least 1 .

Now consider $\mathcal{S}_{=t}$. By hypothesis $\mathcal{H} 1$, for any $\mathbf{g} \in \mathcal{O}_{1}, \pi^{-1}(g) \cap V$ is non-empty, hence $\pi(V)$ contains the non-empty Zariski-open set $O_{1} \subset \mathbb{C}^{t}$. The intersection $O_{1} \cap \mathbb{R}^{t}$ is a non-empty Zariski-open set of $\mathbb{R}^{t}$, contained in $\pi(V)$, hence $\pi(V) \cap \mathbb{R}^{t}$ has real dimension $t$. Let $U_{0}$ be its interior. The subset $\mathcal{S}_{=t} \cap \pi^{-1}\left(U_{0}\right)$ is a locally closed semi-algebraic set. If it is empty, then there is nothing to prove. Otherwise, by construction it has dimension $t$; and the projection $\pi$ restricted to this subspace is proper, by hypothesis $\mathcal{H} 2$. Thom's isotopy lemma (Theorem 1.82) states that for any nonempty connected open set $U$ of $\mathbb{R}^{t}$ not meeting $K(\pi, V)$, and for any $\mathrm{g} \in U$, there exists a semi-algebraic diffeomorphism

$$
h=\left(h_{0}, \pi\right): V \cap B \cap \pi^{-1}(U) \xrightarrow{\sim} \pi^{-1}(\mathbf{g}) \times U .
$$

By hypothesis $\mathcal{H} 1$, if $U$ is nonempty, $\pi^{-1}(\mathrm{~g})$ is finite, and the cardinality of the fibers is constant on $U$.

So in order to compute the wanted decomposition of the parameter space, it suffices to compute a polynomial $P \in \mathbb{Q}[\mathbf{G}]$ such that $V(P)$ covers $\pi\left(V \cap B_{0}\right)$ and $K(\pi, V)$.

### 4.3.2. The determinantal problem

Let $k$ be an integer greater than $1, r_{0} \in\{1, \ldots, k-1\}$, and

$$
\begin{equation*}
n=\left(k-r_{0}+1\right)^{2} . \tag{4.2}
\end{equation*}
$$

Let $t \in \mathbb{N}$, and let $M(\mathbf{X}, \mathrm{G})$ be a $k \times k$ matrix with polynomial entries in $n$ variables $\mathbf{X}=$ $\left(X_{1}, \ldots, X_{n}\right)$ and $t$ parameters $\mathrm{G}=\left(G_{1}, \ldots, G_{t}\right)$. As before, let $\pi: \mathbb{C}^{n+t} \longrightarrow \mathbb{C}^{t}$ be the projection onto the affine space with coordinates G .


Figure 4.2: Roots classification using Thom's isotopy lemma: we classify the points of the $g$ axis according to the cardinality of the fibers of $\pi$ restricted to $V \cap B$. The points where this cardinality may change are the projections of $P_{1}, P_{2}$ and $P_{3}: P_{1}$ is in the intersection of $V$ and the boundary of $B, P_{2}$ is a critical point of $\pi$ restricted to $V$ and $P_{3}$ is a singular point of $V$.

Let $\{h(\mathbf{X}, \mathbf{G}) \leq 0 \mid h \in H\}$ be a system of inequalities, with $H \subset \mathbb{Q}[\mathbf{X}, \mathbf{G}]$. We will consider determinantal varieties associated with $M$, as defined in Section 1.5: for any $r \in\{0, \ldots, k\}$, we define the variety

$$
V_{r}=\left\{(\mathbf{x}, \mathbf{g}) \in \mathbb{C}^{n+t} \mid \operatorname{rank}(M(\mathbf{x}, \mathbf{g})) \leq r\right\}
$$

and let $V_{-1}=\varnothing$ by convention. Furthermore, we define the constructible set $V_{=r}=V_{r} \backslash V_{r-1}$, that is the set of points at which the matrix $M$ has rank exactly $r$.

Let $V$ be the union of the singular locus of $V_{r_{0}}$ and of the set of critical points of $\pi$ restricted to $V_{r_{0}}$. In other words,

$$
\pi(V)=K\left(\pi, V_{r_{0}}\right)
$$

We want to classify the cardinality of the real fibers by $\pi$ of the semi-algebraic set

$$
V \cap\{(\mathbf{x}, \mathbf{g}) \mid \forall h \in H, h(\mathbf{x}, \mathbf{g}) \leq 0\}
$$

Assume that $V$ and $H$ satisfy hypotheses $\mathcal{H} 1, \mathcal{H} 2, \mathcal{H} 3$ and $\mathcal{H} 4$. Further assume that:
$\mathcal{H} 5$. There exists a non-empty Zariski-open subset $O_{2} \subset \mathbb{C}^{t}$ such that

$$
V \cap \pi^{-1}\left(O_{2}\right)=V_{r_{0}-1} \cap \pi^{-1}\left(O_{2}\right)
$$

$\mathcal{H} 6$. For any $r \in\{0, \ldots, k-1\}$, the ideal defined by the $(r+1)$-minors of $M$ is radical;
$\mathcal{H} 7$. For any $r \in\{0, \ldots, k-1\}$, the variety $V_{r}$ is equidimensional with dimension $n+t-(k-$ $r+1)^{2}$.

By Proposition 1.92, these properties are generic. Furthermore, hypothesis $\mathcal{H} 7$ implies hypothesis $\mathcal{H} 4$ (by definition (4.2) of $n$ ), and generically it implies hypotheses $\mathcal{H} 1$.
In following subsections, we will describe two algorithms DeterminantCritVals and DeterminantBoundary, which, given such a matrix $M$, a target rank $r_{0}$ and inequalities $H$, compute respectively a polynomial whose zeroes cover $K(\pi, V)$, and a polynomial whose zeroes cover $\pi\left(V \cap B_{0}\right)$. By Lemma 4.1, the zeroes of the product of these polynomials will subdivide the parameter space into connected components where the cardinality of real fibers is constant. These algorithms are probabilistic, because they will rely on the choice of generic linear forms to ensure linear independence. However, the algorithms could be made deterministic by testing that these linear forms are generic enough for our purpose, and repeating the random draw otherwise.
The algorithms will also need to compute the projection of algebraic sets onto coordinate subspaces. For this purpose, we assume that we are given a routine Elimination, which, given a system of polynomials $F \subset \mathbb{Q}\left[V_{1}, \ldots, V_{N}\right]$ and a set of variables $\mathbf{V}^{\prime} \subset\left\{V_{1}, \ldots, V_{N}\right\}$, computes a system of generators of $\langle F\rangle \cap \mathbb{Q}\left[V^{\prime}\right]$. Such a routine can be implemented using Gröbner bases or regular chains, for example.

### 4.3.3. Incidence varieties

We decompose the problem depending on the rank of the matrix. For that purpose, we use incidence varieties, as defined in Section 1.5.3.

Recall that the incidence variety of rank $r$ associated with $M$ is the variety $\mathcal{V}_{r} \subset \mathbb{C}^{n+t} \times$ $\left(\mathbb{P}^{k-1}(\mathbb{C})\right)^{k-r}$ defined by:

$$
M \cdot\left[\begin{array}{ccc}
y_{1,1} & \ldots & y_{1, k-r}  \tag{4.3}\\
\vdots & & \vdots \\
y_{k, 1} & \ldots & y_{k, k-r}
\end{array}\right]=\left[\begin{array}{ccc}
0 & \ldots & 0 \\
\vdots & & \vdots \\
0 & \ldots & 0
\end{array}\right]
$$

with the additional condition that the matrix $\left(y_{i, j}\right)$ has rank $k-r$.
Let $\left(u_{1,1}, \ldots, u_{k-r, k}\right) \in \mathbb{C}^{k(k-r)}$, we define the variety $\mathcal{V}_{r, \mathrm{u}}^{\prime}$ as the intersection of $\mathcal{V}_{r}$ and the complex solutions of

$$
\left[\begin{array}{ccc}
u_{1,1} & \ldots & u_{1, k}  \tag{4.4}\\
\vdots & & \vdots \\
u_{k-r, 1} & \ldots & u_{k-r, k}
\end{array}\right] \cdot\left[\begin{array}{ccc}
y_{1,1} & \ldots & y_{1, k-r} \\
\vdots & & \vdots \\
y_{k, 1} & \ldots & y_{k, k-r}
\end{array}\right]=\operatorname{Id}_{k-r},
$$

this encodes that the matrix $\left(y_{i, j}\right)$ has rank $k-r$ and it fixes an affine chart for the homogeneous coordinates of the columns of $\left(y_{i, j}\right)$.
In the rest of Section 4.3, $\mathcal{F}_{r, \mathrm{u}}$ denotes the union of Equations (4.3) and (4.4).

Proposition 4.2. Let $r_{1}=r_{0}-1$. Let $\varphi: \mathcal{V}_{r_{1}, \mathrm{u}}^{\prime} \longrightarrow \mathbb{C}^{t}$ be the projection onto the affine space with coordinates G . Assuming that hypotheses $\mathcal{H} 1$ to $\mathcal{H} 7$ hold, there exists a Zariski-open subset $\mathcal{U} \subset \mathbb{C}^{k\left(k-r_{1}\right)}$ such that if $\mathbf{u} \in \mathcal{U} \cap \mathbb{Q}^{k\left(k-r_{1}\right)}, \overline{K\left(\pi, V_{r_{1}}\right)}=\overline{K\left(\varphi, \mathcal{V}_{r_{1}, \mathbf{u}}^{\prime}\right)}$.

Proof. Let $P=(\mathbf{x}, \mathbf{g}, \mathbf{y}) \in \mathcal{V}_{r_{1}, \mathbf{u}}^{\prime}$.
If $M(\mathbf{x}, \mathbf{g})$ has rank less than $r_{1}$, then by Proposition $1.91,(\mathbf{x}, \mathbf{g}) \in \operatorname{sing}\left(V_{r_{1}}\right)$, hence $\mathbf{g} \in K\left(\pi, V_{r_{1}}\right)$.
Since $M(\mathbf{x}, \mathbf{g})$ has rank less than $r_{1}$, its kernel $L_{1}$ has dimension at least $k-r_{1}+1$. Equations (4.4) encode that the vectors $\mathrm{y}_{i}$ given by the columns of matrix $\left(y_{i, j}\right)$ generate a $r_{1}$-dimensional vector space $L_{2}$. So there exists $\mathbf{y}_{0} \in L_{1} \cap L_{2}$, and for all $a \in \mathbb{C},\left(\mathbf{x}, \mathbf{g}, \mathbf{y}_{1}+a \mathbf{y}_{0}, \mathbf{y}_{2}\right)$ belongs to the fiber above ( $\mathrm{x}, \mathrm{g}$ ) in $\mathcal{V}_{r_{1}}$. So this fiber has dimension at least 1 , while the generic fiber has dimension 0 by hypothesis $\mathcal{H} 1$. So $(\mathrm{x}, \mathrm{g})$ is a critical value of the projection of $\mathcal{V}_{r_{1}}$ onto $\mathbb{R}^{n+t}$, hence $(\mathrm{g}) \in K\left(\varphi, \mathcal{V}_{r_{1}, \mathrm{u}}^{\prime}\right)$.

So we may assume that $M(\mathbf{x}, \mathrm{~g})$ has rank exactly $r_{1}$. There is a $r_{1} \times r_{1}$ submatrix $A$ of $M(\mathbf{x}, \mathbf{g})$ which is invertible, without loss of generality we may assume that it is the top-left $r_{1} \times r_{1}$ submatrix. In an open neighborhood of $(\mathbf{x}, \mathbf{g}), V_{=r_{1}}$ is described by the vanishing of the entries of $M / A$, that is the determinants of the $\left(r_{1}+1\right) \times\left(r_{1}+1\right)$ submatrices containing $A$. The same computations as in the proof of Lemma 1.100 give the following equations describing $\mathcal{V}_{r_{1}, \mathrm{u}}^{\prime}$ in the open neighborhood of $(\mathbf{x}, \mathbf{g}, \mathbf{y})$ where $\Delta=\operatorname{det}(A)$ does not vanish:

$$
\left\{\begin{array}{l}
M / A=0  \tag{4.5}\\
Y_{(2)}=\left(U_{(2)}-U_{(1)} A^{-1} B\right)^{-1} \\
Y_{(1)}=\frac{-1}{\Delta} A^{-1} B Y_{(2)}
\end{array}\right.
$$

and the truncated Jacobian matrix in ( $\mathrm{X}, \mathrm{Y}$ ) of this system can be written

$$
\left[\begin{array}{ccc}
\operatorname{Jac}_{\mathrm{X}}(M / A) & 0 & 0 \\
\star & \operatorname{Id}_{\left(k-r_{1}\right)^{2}} & \star \\
\star & 0 & \operatorname{Id}_{r_{1}\left(k-r_{1}\right)}
\end{array}\right]
$$

where $\operatorname{Jac}_{\mathrm{X}}(M / A)$ is the truncated Jacobian matrix in X of the $\left(k-r_{1}\right)^{2}$ entries of $M / A$, which define $V_{r_{1}} \backslash\{(\mathrm{x}, \mathrm{g} \mid \Delta=0)\}$ in $\mathbb{C}^{n+t}$. By hypothesis $\mathcal{H} 6$, the ideal defined by the entries of $M / A$, which is a subideal of the ideal of all $\left(r_{1}+1\right)$-minors of $M$, is radical. Since the Schur complement appears by multiplication with invertible matrices with entries in the localized ring $\mathbb{Q}[\mathrm{X}, \mathrm{g}]_{\Delta}$ (using the same notations as in the proof of Lemma 1.100):

$$
\left[\begin{array}{cc}
\mathrm{Id}_{r_{1}} & 0 \\
-C & \operatorname{Id}_{k-r_{1}}
\end{array}\right] \cdot\left[\begin{array}{cc}
A^{-1} & 0 \\
0 & \mathrm{Id}_{k-r_{1}}
\end{array}\right] \cdot M=\left[\begin{array}{cc}
\operatorname{Id}_{r_{1}} & A^{-1} B \\
0 & M / A
\end{array}\right],
$$

Equations (4.5) describe the localization of $\left\langle\mathcal{F}_{r_{1}, \mathbf{u}}\right\rangle$ in $\mathbb{Q}[\mathbf{X}, \mathbf{g}]_{\Delta}$, so this ideal is radical as well. So we can use the Jacobian criterion on $V_{r_{1}}$ near $(\mathbf{x}, \mathrm{g})$ and on $\mathcal{V}_{r_{1}, \mathbf{u}}^{\prime}$ near $(\mathbf{x}, \mathrm{g}, \mathrm{y})$. Both Jacobians matrices have the same rank and both varieties have the same local codimension $\left(k-r_{1}\right)^{2}$ (by Lemma 1.100 and hypothesis $\mathcal{H}($ ), so

$$
K\left(\pi, V_{=r_{1}}\right) \cap \varphi\left(\mathcal{V}_{r_{1}, \mathbf{u}}^{\prime}\right)=K\left(\varphi, \mathcal{V}_{r_{1}, \mathbf{u}}^{\prime}\right) \cap \pi\left(V_{=r_{1}}\right)
$$

The image $\varphi\left(\mathcal{V}_{r_{1}, \mathrm{u}}^{\prime}\right) \cap \pi\left(V_{=r_{1}}\right)$ is a Zariski-open subset $O_{\mathrm{u}}$ of $\pi\left(V_{=r_{1}}\right)$. It remains to prove that if $\mathbf{u}$ is sufficiently generic, then all irreducible components of $K\left(\pi, V_{=r_{1}}\right)$ meet this open subset.

Let $C_{1}, \ldots, C_{a}$ be these irreducible components, and let

$$
\left(\mathbf{x}_{1}, \mathbf{g}_{1}\right) \in \pi^{-1}\left(C_{1}\right), \ldots,\left(\mathbf{x}_{a}, \mathbf{g}_{a}\right) \in \pi^{-1}\left(C_{a}\right)
$$

For any $(\mathbf{x}, \mathbf{g}) \in V_{=r_{1}}$, the proof of $\left[H_{N S} 15\right.$ b, Prop. 4, Sec. 6] shows that there exists a non-empty Zariski-open subset $\mathcal{U}_{(\mathrm{x}, \mathrm{g})} \subset \mathbb{C}^{k\left(k-r_{1}\right)}$ such that if $\mathbf{u} \in \mathcal{U}_{(\mathrm{x}, \mathrm{g})} \cap \mathbb{Q}^{k\left(k-r_{1}\right)}$, then $(\mathbf{x}, \mathbf{g}) \in \mathcal{O}_{\mathbf{u}}$; namely, $\mathcal{U}_{(\mathrm{x}, \mathrm{g})}$ is the set of $\mathbf{u}$ such that

$$
\operatorname{rank}\left[\begin{array}{c}
M(\mathbf{x}, \mathbf{g}) \\
\left(u_{i, j}\right)
\end{array}\right]=k
$$

Taking the finite intersection of the non-empty Zariski-open subsets $\mathcal{U}_{\left(\mathrm{x}_{i}, \mathrm{~g}_{i}\right)}$ for $i \in\{1, \ldots, a\}$ yields the wanted subset $\mathcal{U}$.

### 4.3.4. Locus of rank exactly $r_{0}$

Recall that by $\mathcal{H} 5, \pi\left(V \cap V_{=r_{0}}\right)$ has codimension at least 1 , and that we want to compute a polynomial whose zeroes cover $K(\pi, V)$ and $\pi\left(V \cap B_{0}\right)$. So we may multiply the result by the equation of one hypersurface covering $\pi\left(V \cap V_{=r_{0}}\right)$, it will naturally cover $\pi\left(V \cap V_{=r_{0}}\right) \cap K(\pi, V)$ and $\pi\left(V \cap V_{=r_{0}}\right) \cap \pi\left(V \cap B_{0}\right)$.

```
Algorithm 4.1 RankExactly
Input: \(M \subset \mathbb{Q}[\mathbf{X}, \mathrm{G}]^{k \times k}, r_{0} \in\{1, \ldots, k-1\}\)
Output: \(P_{1} \in \mathbb{Q}[\mathrm{G}] \backslash\{0\}\) s.t. \(\pi\left(V \cap V_{=r_{0}}\right) \subset V\left(P_{1}\right)\)
    res \(\leftarrow 1\)
    \(F_{V, 0} \leftarrow\left\{\left(r_{0}+1\right)\right.\)-minors of \(\left.M\right\}\)
    \(J \leftarrow \mathrm{Jac}_{\mathbf{X}}\left(F_{V, 0}\right)\)
    \(F_{V, 1} \leftarrow F_{V, 0} \cup\left\{\left(k-r_{0}\right)^{2}\right.\)-minors of \(\left.J\right\}\)
    Pick at random \(u_{1}, \ldots, u_{k\left(k-r_{0}\right)}=\mathbf{u} \in \mathbb{Q}^{k\left(k-r_{0}\right)}\)
    \(F_{0} \leftarrow \mathcal{F}_{k-r_{0}, \mathbf{u}}\)
    \(\left\{\mathcal{M}_{1}, \ldots, \mathcal{M}_{N}\right\} \leftarrow\left\{r_{0}\right.\)-minors of \(\left.M\right\}\)
    for \(i\) in \(\{1, \ldots, N\}\) do
        \(F_{1} \leftarrow F_{0} \cup F_{V, 1} \cup\left\{\mathcal{M}_{1}, \ldots, \mathcal{M}_{i-1}, u \cdot \mathcal{M}_{i}-1\right\}\)
        \(G \leftarrow\) Elimination \(\left(F_{1},\{u, \mathrm{X}, \mathrm{Y}\}\right)\)
        Multiply res by 1 polynomial from \(G\)
    end for
    return res
```


### 4.3.5. Singularities

```
Algorithm 4.2 DeterminantCritVals
Input: \(M \subset \mathbb{Q}[\mathrm{X}, \mathrm{G}]^{k \times k}, r_{0} \in\{1, \ldots, k-1\}\)
Output: \(P_{c} \in \mathbb{Q}[\mathrm{G}] \backslash\{0\}\) s.t. \(K(\pi, V) \subset V\left(P_{c}\right)\)
    res \(\leftarrow \operatorname{RankExactly}\left(M, r_{0}\right)\)
    Pick at random \(u_{1}, \ldots, u_{k\left(k-r_{0}+1\right)}=\mathbf{u} \in \mathbb{Q}^{k\left(k-r_{0}+1\right)}\)
    \(F_{0} \leftarrow \mathcal{F}_{r_{0}-1, \mathbf{u}}\)
    \(J \leftarrow \mathrm{Jac}_{\mathrm{X}, \mathrm{Y}}\left(F_{0}\right)\)
    \(F_{1} \leftarrow F_{0} \cup\left\{k\left(k-r_{0}+1\right)+\left(k-r_{0}+1\right)^{2}\right.\)-minors of \(\left.J\right\}\)
    \(G \leftarrow\) Elimination \(\left(F_{1},\{\mathrm{X}, \mathrm{Y}\}\right)\)
    Multiply res by 1 polynomial from \(G\)
    return res
```

Proposition 4.3. Algorithm DeterminantCritVals is correct.
Proof. By definition, $V \subset V_{r_{0}}$. Using the decomposition $V_{r_{0}}=V_{=r_{0}} \cup V_{r_{0}-1}$, we decompose the variety $V$ as $V=\left(V \cap V_{=r_{0}}\right) \cup\left(V \cap V_{r_{0}-1}\right)$.

The subspace $\pi\left(V \cap V_{=r_{0}}\right)$ is covered by the output of RankExactly, so we may restrict to $V \cap V_{r_{0}-1}$, which is the whole variety $V_{r_{0}-1}$ by Lemma 1.91.
By Proposition 4.2, in order to compute $K\left(\pi, V_{r_{0}-1}\right)$, we can compute polynomials whose zeroes cover $K\left(\varphi, \mathcal{V}_{r_{0}-1, \mathbf{u}}^{\prime}\right)$ with $\mathbf{u}$ sufficiently generic instead.

By hypotheses $\mathcal{H} 6, \mathcal{H} 7$ and the proof of Proposition $4.2, \mathcal{V}_{r_{0}-1, \mathrm{u}}^{\prime}$ is $t$-equidimensional and $\mathcal{F}_{r_{0}-1, \mathrm{u}}$ is a set of generators of its ideal, so we can use the Jacobian criterion to compute equations defining $K\left(\varphi, \mathcal{V}_{r_{0}-1, \mathbf{u}}^{\prime}\right)$.

### 4.3.6. Boundary

```
Algorithm 4.3 DeterminantBoundary
Input: \(M \subset \mathbb{Q}[\mathbf{X}, \mathrm{G}]^{k \times k}, r_{0} \in\{1, \ldots, k-1\}, H \subset \mathbb{Q}[\mathbf{X}, \mathrm{G}]\)
Output: \(P_{b} \in \mathbb{Q}[\mathrm{G}] \backslash\{0\}\) s.t. \(\pi\left(V \cap B_{0}\right) \subset V\left(P_{b}\right)\)
    res \(\leftarrow \operatorname{RankExactly}\left(M, r_{0}\right)\)
    Pick at random \(u_{1}, \ldots, u_{k\left(k-r_{0}+1\right)}=\mathbf{u} \in \mathbb{Q}^{k\left(k-r_{0}+1\right)}\)
    \(F_{0} \leftarrow \mathcal{F}_{r_{0}-1, \mathbf{u}}\)
    for \(h\) in \(H\) do
        \(F_{1} \leftarrow F_{0} \cup\{h\}\)
        \(G \leftarrow\) Elimination \((F,\{\mathbf{X}, \mathbf{Y}\})\)
        Multiply res by 1 polynomial from \(G\)
    end for
    return res
```

Proposition 4.4. Algorithm DeterminantBoundary is correct.

Proof. As in Section 4.3.5, we write: $V=\left(V \cap V_{=r_{0}}\right) \cup\left(V \cap V_{r_{0}-1}\right)$. Since $B_{0}=\bigcup_{h \in H} V(h)$, the intersection $V \cap \partial B$ is contained in the union of the varieties $V(\langle F\rangle+\langle h\rangle)$ for $h$ ranging over $H$, and the equation of the projections can be obtained with polynomial elimination.

Remark 4.5. For the real root classification problem, the subdivision is given by the product of the outputs of DeterminantCritVals and DeterminantBoundary. In order to avoid repeating computations, we may skip the call to RankExactly in either subroutine (but not both), and initialize res to 1 instead.

### 4.4. The contrast problem

### 4.4.1. The case of water

With the notations of Section 4.2 , the variety $V$ is the complex algebraic variety defined by

$$
D=\frac{\partial D}{\partial y_{1}}=\frac{\partial D}{\partial y_{2}}=\frac{\partial D}{\partial z_{1}}=\frac{\partial D}{\partial z_{2}}=0 .
$$

With the notations of Section $4 \cdot 3$, we want to classify the singularities of the set of points where $M$ has rank at most $r_{0}=3$. Our semi-algebraic constraints are that the solutions are within the Bloch ball, that is

$$
\mathcal{B}:\left\{\begin{array}{l}
h_{1}=y_{1}^{2}+\left(z_{1}+1\right)^{2} \leq 1 \\
h_{2}=y_{2}^{2}+\left(z_{2}+1\right)^{2} \leq 1
\end{array}\right.
$$

Since the equations are homogeneous in $\Gamma_{1}, \Gamma_{2}, \gamma_{1}, \gamma_{2}$, and the parameters are supposed to be nonzero, we may normalize by setting $\gamma_{1}=1$. In the case where the first matter is water, we further simplify by setting $\Gamma_{1}=\gamma_{1}=1$, leaving free the two parameters $\Gamma_{2}, \gamma_{2}$ corresponding to the second matter. We recall that we also assume that $2 \Gamma_{2} \geq \gamma_{2}$ and that $\left(\gamma_{2}, \Gamma_{2}\right) \neq(1,1)=\left(\gamma_{1}, \Gamma_{1}\right)$ (that is, the second matter is not water).

This system satisfies hypotheses $\mathcal{H} 1$ to $\mathcal{H} 7$.
Theorem 4.6. Consider the 9 polynomials:

$$
\begin{aligned}
& f_{1}=\Gamma_{2}-1 \\
& f_{2}=3 \Gamma_{2}-2 \gamma_{2}-1 \\
& f_{3}=3 \Gamma_{2}^{2}-5 \Gamma_{2} \gamma_{2}+\gamma_{2}^{2}+2 \Gamma_{2}-2 \gamma_{2}+1 \\
& f_{4}=2 \Gamma_{2}^{2}-5 \Gamma_{2} \gamma_{2}+2 \gamma_{2}^{2}-2 \Gamma_{2}+3 \gamma_{2} \\
& f_{5}=2 \gamma_{2}^{3}-\left(3 \Gamma_{2}+11\right) \gamma_{2}^{2}+\left(9 \Gamma_{2}+6-3 \Gamma_{2}^{2}\right) \gamma_{2}+2 \Gamma_{2}\left(\Gamma_{2}+2\right)\left(\Gamma_{2}-1\right) \\
& f_{6}=\Gamma_{2}-2 \gamma_{2}+1 \\
& f_{7}=2 \Gamma_{2}-\gamma_{2}-1 \\
& f_{8}=\gamma_{2}-2+\Gamma_{2} \\
& f_{9}=2 \Gamma_{2}^{2}-5 \Gamma_{2} \gamma_{2}+2 \gamma_{2}^{2}+1
\end{aligned}
$$

The zeroes of their product divide the subset of $\mathbb{R}^{2}$ defined by $2 \Gamma_{2}>\gamma_{2}>0$ into connected components where the cardinality of $V_{\mathbb{R}} \cap \pi^{-1}\left(\gamma_{2}, \Gamma_{2}\right)$ is constant.


Figure 4.3: Curves involved in the definition of the semi-algebraic set $\mathcal{G}$. Sample points marked with a blue diamond (resp. a green circle) are points in $\mathcal{G}_{1}^{-} \cup \mathcal{G}_{1}^{+}\left(\right.$resp. $\left.\mathcal{G}_{2}^{-} \cup \mathcal{G}_{2}^{+}\right)$. Parameters in the red area are physically irrelevant.

Let $\psi:\left(y_{1}, z_{1}, y_{2}, z_{2}\right) \longmapsto\left(-y_{1}, z_{1},-y_{2}, z_{2}\right)$ be the symmetry fixing $\Pi=\left\{y_{1}=y_{2}=0\right\}$, and let us consider the semi-algebraic sets (see Figure 4.3):

$$
\begin{aligned}
\mathcal{G}_{1}^{-} & =\left\{\gamma_{2}<2 \Gamma_{2}, \Gamma_{2}<1, f_{2}>0, f_{4}<0\right\}, \\
\mathcal{G}_{1}^{+} & =\left\{\gamma_{2}<2 \Gamma_{2}, \Gamma_{2}>1, f_{2}<0, f_{4}>0\right\}, \\
\mathcal{G}_{2}^{-} & =\left\{\Gamma_{2}<1, f_{6}>0, f_{3}<0\right\}, \\
\mathcal{G}_{2}^{+} & =\left\{\Gamma_{2}>1, f_{6}<0, f_{5}>0\right\}, \\
\mathcal{G} & =\mathcal{G}_{1}^{-} \cup \mathcal{G}_{1}^{+} \cup \mathcal{G}_{2}^{-} \cup \mathcal{G}_{2}^{+} .
\end{aligned}
$$

Theorem 4.7. For all $\left(\gamma_{2}, \Gamma_{2}\right)$ such that $2 \Gamma_{2}>\gamma_{2}>0$, the center $O$ of the Bloch ball $\mathcal{B}$ is a singularity of $\{D=0\}$. And provided $\left(\gamma_{2}, \Gamma_{2}\right) \in \mathcal{G}$, there exist at most two other singularities:

1. provided $\left(\gamma_{2}, \Gamma_{2}\right) \in \mathcal{G}_{1}^{-} \cup \mathcal{G}_{1}^{+}$there is one other singularity lying on $\Pi \cap \mathcal{B}$;
2. provided $\left(\gamma_{2}, \Gamma_{2}\right) \in \mathcal{G}_{2}^{-} \cup \mathcal{G}_{2}^{+}$there are two other singularities in $\mathcal{B}, \psi$-symmetric, outside $\Pi$.

The configuration is illustrated in Figures 4.3 and 4.4. Observe that the number of singularities inside $\mathcal{B}$ is an invariant of the contrast problem. Two of the pairs of biological matters studied in [Bon+13], water-cerebrospinal fluid (normalized parameters [ $\left.\gamma_{2}=\frac{5}{4}, \Gamma_{2}=\frac{25}{3}\right]$ ) and water-fat (normalized parameters $\left[\gamma_{2}=\frac{25}{2}, \Gamma_{2}=25\right]$ ) correspond to points outside $\mathcal{G}$, and their invariant is 1 in both cases (see Figure $4 \cdot 4$ ). But our results give answers to our guiding questions: there


Figure 4.4: Positions of the parameters corresponding to the pairs water-cerebrospinal fluid (red circle) and water-fat (red square) and the set $\mathcal{G}$ (with the same conventions as in Figure 4.3).
exist pairs of matters for which this algebraic invariant can differ; and any pair (water, matter) belongs to one of 3 classes, depending on whether the number of singularities inside $\mathcal{B}$ is 1,2 or 3 .

Proof of Theorem 4.6. Let $V=3=\left\{p \in \mathbb{C}^{4} \times \mathbb{R}^{2} \mid \operatorname{rank}(M)=3\right\}$ and $V_{2}=\left\{p \in \mathbb{C}^{4} \times \mathbb{R}^{2} \mid\right.$ $\operatorname{rank}(M)<3\}$, where $p=\left(y_{1}, y_{2}, z_{1}, z_{2}, \gamma_{2}, \Gamma_{2}\right)$. We apply the strategies described in Section 4.3.
We study the generic case $V_{2} \cap V$ first. This set does cover a dense subset of $\mathbb{R}^{2}$. Its intersection with the boundary of $\mathcal{B}$ is given by the vanishing of either $h_{1}$ or $h_{2}$. The projection on $\left(\Gamma_{2}, \gamma_{2}\right)$ of the set of points of $V_{2} \cap V$ such that $h_{1}=0$ is described by $0=\gamma_{2} f_{1}^{2} f_{2} f_{3}$ which gives us polynomials $f_{1}, f_{2}$ and $f_{3}$.
The projection on $\left(\Gamma_{2}, \gamma_{2}\right)$ of the set of points of $V_{2} \cap V$ such that $h_{2}=0$ is described by $0=\left(2 \Gamma_{2}-\gamma_{2}\right) f_{1}^{2} f_{4} f_{5}$ which gives us new polynomials $f_{4}$ and $f_{5}$.

| System | RegularChain <br> (direct) | Gröbner <br> (direct) | RegularChain <br> (new algo.) | FGb <br> (new algo.) | F $_{5}$ (new algo.) | CAD |
| :--- | ---: | ---: | ---: | :--- | ---: | ---: |
| Water | 1600 s | 100 s |  | 10 s | 1 s | 50 s |
| General | $>24 \mathrm{~h}$ | $>24 \mathrm{~h}$ | $90 \times 200 \mathrm{~s}$ | $46 \times 200 \mathrm{~s}$ | 110 s | 4 h (projection step) |

Table 4.1: Timings

Next, we consider the incidence variety $\mathcal{V}_{2}$ associated with the matrix $M$ :

$$
M \cdot\left[\begin{array}{ll}
\lambda_{1,1} & \lambda_{1,2} \\
\lambda_{2,1} & \lambda_{2,2} \\
\lambda_{3,1} & \lambda_{3,2} \\
\lambda_{4,1} & \lambda_{4,2}
\end{array}\right]=\left[\begin{array}{ll}
0 & 0 \\
0 & 0 \\
0 & 0 \\
0 & 0
\end{array}\right]
$$

with random linear equations ensuring that the matrix $\left(\lambda_{i, j}\right)$ has rank 2.
Out of the surface $\gamma_{2}=0$, this affine variety is a complete intersection (it has dimension 2 and it is given by 9 equations in 11 variables, including the saturation by $\gamma_{2}$ ). The set of critical values of $\pi$ is described by $0=\left(2 \Gamma_{2}-\gamma_{2}\right)\left(\Gamma_{2}+1\right) f_{1}^{2} f_{6}^{2} f_{7}^{2}$ which gives us new polynomials $f_{6}$ and $f_{7}\left(\Gamma_{2}+1\right.$ has no solutions within our constraint range).

This completes the study of $V \cap V_{2}$. We now move on to the study of $V \cap V_{=3}$. As described in the algorithm, we define the incidence variety of rank 3 of $M$, and we saturate successively by the 3-minors of $M$. Only the first of these subcases is nonempty, and it is described by $0=\left(2 \Gamma_{2}-\gamma_{2}\right) f_{8} f_{9}$ which gives us $f_{8}$ and $f_{9}$.

Proof of Theorem 4.7. Observe first by means of a trivial evaluation that $O$ is a singularity of $\{D=0\}$. We now focus on singularities in $\mathcal{B}^{*}=\mathcal{B} \backslash\{O\}$. Theorem 4.6 provides a list of 9 polynomials to which we add our constraints $2 \Gamma_{2} \geq \gamma_{2}>0$. Let $\xi=\gamma_{2} \Gamma_{2}\left(\gamma_{2}-2 \Gamma_{2}\right) \prod_{i=1}^{9} f_{i}$. The complementary of $\{\xi=0\}$ is the union of a sequence of connected open semi-algebraic sets where the number of singularities is constant. The routine CylindricalAlgebraicDecompose of the Maple package RegularChains[SemiAlgebraicSetTools] provides 1533 sample points. Excluding those at which $\xi$ vanishes and those outside our physical constraints domain, remains a set $K_{c}$ of 548 points. At each point of $K_{c}$ we locate the singularities by computing a Gröbner basis.

We get 165 points of $K_{c}$ such that there exists at least one singularity in $\mathcal{B}^{*}$. We have a set $K_{s}$ of 37 points, each of them corresponding to a couple of $\psi$-symmetric singularities outside the symmetry plane $\Pi$, and a set $K_{p}$ of 128 points corresponding to a unique singularity on $\Pi \cap \mathcal{B}^{*}$. For parameters at which $\xi$ does not vanish, the number of singularities in $\mathcal{B}^{*}$ is at most two.

Points of $K_{s}$ (resp. $K_{p}$ ) are represented in green (resp. blue) in Figures 4.3, 4.5 and 4.6. Let us evaluate on $K_{c}$ the condition $\left(\Gamma_{2}<1, f_{2}>0, f_{4}<0\right)$ or $\left(\Gamma_{2}>1, f_{2}<0, f_{4}>0\right)$. Indeed the set of points of $K_{c}$ satisfying this condition coincides with $K_{p}$. This proves item 1. The proof of item 2 is similar.


Figure 4.5: The curves involved in the decomposition of the region $\Gamma_{2}>0, \gamma_{2}>0,2 \Gamma_{2} \geq \gamma_{2}$ of the parameter space (with the same conventions as in Figure 4.3)


Figure 4.6: Decomposition of the parameter space near $(1,1)$ (with the same conventions as in Figure 4.3).

### 4.4.2. The general case

The variety $V$ and the semi-algebraic set $\mathcal{B}$ are defined as in the previous section. We normalize again by $\gamma_{1}=1$, we assume that $2 \Gamma_{1} \geq 1,2 \Gamma_{2} \geq \gamma_{2}>0,\left(\gamma_{2}, \Gamma_{2}\right) \neq\left(1, \Gamma_{1}\right)$, and that $\Gamma_{1} \neq 1$, $\Gamma_{2} \neq \gamma_{2}$ (case of water).

Theorem 4.8. Splitting the subset of $\mathbb{R}^{3}$ defined by $2 \Gamma_{2}>\gamma_{2}>0$ and $2 \Gamma_{1}>1$ into open subsets where the number of real singularities of $V$ in the fibers is constant, can be done by cutting out 12 irreducible surfaces, consisting of 5 planes, 3 quadrics, two surfaces of degree 9 and one of degree 14.

These polynomials were obtained by applying the algorithms from Section 4.3 to our system. The elimination steps were done using both Gröbner bases with FGb or with $\mathrm{F}_{5}$, and with triangular sets with RegularChains. Table 4.1 presents some timings for these methods (for computations done with interpolation, we give the results as $a \times b$ where $a$ is the interpolation degree and $b$ the time taken for each specialized computation).

## Part III

## Appendices

## Appendix A

## Source code

## A.1. Magma source code for Matrix- $F_{5}$ (Chapters 2 and 3)

## A.1.1. Algorithm Matrix- $F_{5}$

File F5Mtx.m
// load "F5_comptes_clean.m";
// *** Aux functions
// ****** General purpose
function Status (Calc)
// String representing the current status of the calculation, in a
// very short way. It is meant to be used in verbose output.
res := Sprintf("[\%0](n=%25o)", Cputime(), Calc‘n);
return res;
end function;
function SortedMonomials0fWeightedDegree(P,d)
// Returns the sequence of monomials of $W$-deg $d$, sorted with
// decreasing order
M := MonomialsOfWeightedDegree(P,d);
l := IndexedSetToSequence(M);
Sort(~l); // Inc. order
Reverse(~l); // Reversion for dec. order
return l;
end function;
function DReg (F)
// Computes the theoretical degree of regularity of $F$
D := [Degree(f) : f in F];
res := \&+[d-1 : d in D] +1;
return res;
end function;
function MinVariableDividing(mu)

```
    // Returns the smallest variable dividing monomial mu
    // /!\ Smallest variable = Highest index /!\
    P := Parent(mu);
    e := Exponents(mu);
    i := Max([i : i in [1..#e] | e[i] ne 0 ]);
    x := Parent(mu).i;
    return x;
end function;
function LeadingMonomialOrZero(f)
    // Same as LeadingMonomial, but returns 0 instead of an error
    // if f=0
    if f eq 0 then
        return 0;
    else
        return LeadingMonomial(f);
    end if;
end function;
function PolynomialToMatrixRow(f,Columns,N)
    P := Parent(f);
    R := BaseRing(P);
    Coefs,Mons := CoefficientsAndMonomials(f);
    M := ZeroMatrix(R,1,N);
    /* printf "Mons : %o\nColumns : %o\n\n", Mons, Columns; */
    for c in [1..#Coefs] do
        i := Columns[Mons[c]];
        M[1,i] := Coefs[c];
    end for;
    return M;
end function;
function MatrixRowToPolynomial(L,Monomials)
    P := Parent(Monomials[1]);
    res := P!0;
    for i in [1..NumberOfColumns(L)] do
        res +:= L[i]*Monomials[i];
    end for;
    return res;
end function;
function MtxInsertRow (M,L,l)
    M2 := VerticalJoin(M,L);
    return M2;
end function;
function MtxAddRow (M, c,i,j)
    // M[j] <- M[j] + c*M[i]
    M2 := AddRow(M, c,i,j);
```


## return M2;

end function;
function MtxMultiplyRow (M, c,i)
// M[i] <- c*M[i]
M2 := MultiplyRow(M, c,i);
return M2;
end function;
procedure PrintDebug (s)
if GetVerbose("Groebner") ge 3 then print s;
end if;
end procedure;
function RescalePoly(f)
if $f$ eq 0 then
return f;
else
c := LeadingCoefficient(f);
f2 := (1/c)*f;
return f2;
end if;
end function;
// ****** Declarations and functions for signatures
// Record format for polynomials and their signatures
Sig := recformat
<mu, // The signature itself
i : Integers(), // The base polynomial
row : Integers(), // The row in the matrix
valbefore : RngMPolElt,
valafter : RngMPolElt,
d : Integers(), // The degree of the associated GB poly
log_redlist, // The list of the indices reducing this one LTbefore, // Leading term of the associated polynomial
LTafter,
LTfromGB
>;
function SigCreate (i,f,row)
P := Parent(f);
s := rec <Sig |>;
s‘mu := P!1;
s‘i := i;
s'row := row;
s'valbefore := f;
s'valafter := s‘valbefore;

```
    s'd := Degree(f);
    s‘log_redlist := [];
    s‘LTbefore := LeadingMonomialOrZero(f);
    s`LTafter := s`LTbefore;
    s`LTfromGB := s`LTbefore;
    return s;
end function;
function SigEqual (s1,s2)
    return (s1'i eq s2'i) and (s1`mu eq s2`mu);
end function;
function SigMultiplyByVar (s,v,row)
    // INPUT :
    // s : signature
    // v : variable
    // row : row
    // OUTPUT :
    // Signature of v*s, with row row
    PrintDebug(Sprintf("Multiplying by %o the signature %o\n",v,s));
    res := s;
    res`mu := v*s`mu;
    res`valbefore := v*s`valafter;
    res`valafter := res`valbefore;
    res`log_redlist := [];
    res'row := row;
    res`LTbefore := v*s`LTafter;
    res`LTafter := res`LTbefore;
    res'd := 1 + s'd;
    PrintDebug(Sprintf("Result : %o\n",res));
    return res;
end function;
function SigUpdate (s,Calc)
    // Update the value of s anywhere it appears in Calc
    mu := s`mu;
    d := s'd;
    i := s`i;
    l := s'row;
    Calc'SigsFromMons[i][mu] := s;
    Calc`SigsFromRows[d][l] := s;
    return Calc;
end function;
// ****** Specific aux functions
function F5_InsertRow(s,Calc)
    // Do what is necessary to insert the row defined by
    // signature s into the calculation Calc.
```

```
    // It means :
    // - inserting the row in the matrix
    // - inserting the signature in the list of indices
    // - updating the current and last row
    vprintf Groebner,3 :
    "Inserting the row defined by %o\ninto the calculation%o\n\n",
    s, Calc;
    f := s`valbefore;
    i := Calc`i_current;
    d := Calc'd_current;
    row := Calc'row_current;
    mu := s`mu;
    MatrixRow := PolynomialToMatrixRow(f,Calc`Columns[d],Calc`N[d]);
    Calc`M[d] := MtxInsertRow(Calc`M[d],MatrixRow,Calc`row_Current);
    Calc`SigsFromMons[i][mu] := s;
    Calc`SigsFromRows[d][row] := s;
    Calc`lastrow[d][i] := Calc`lastrow[d][i] +1;
    Calc'row_current := Calc'row_current +1;
    col := Calc`Columns[d][s'LTbefore];
    if not IsDefined(Calc`Pivots[d],col) then
        Calc'Pivots[d][col] := row;
    end if;
    return Calc;
end function;
function F5_Criterion (s, Calc)
    // Returns true iff the signature s defines a polynomial
    // which should be added to the matrix.
    //
    // Criterion : true iff the monomial part of s is not a
    // leading term in the matrix at degree d-di and for the
    // i-1'th polynomial
    mu := s`mu;
    d_di := Degree(mu); // d_di = d_current - di
    i := Calc'i_current;
    try
        res := (mu notin Calc'Crit[d_di][i-1]);
    catch e
        res := true;
    end try;
    return res;
end function;
function F5_EchelonForm (Calc : Full := false)
    // Computes the row-echelon form as required by the F5 algorithm
    // (no row or column swapping)
    vprintf Groebner,3 : "Reduction to echelon form of matrix\n%o\n\n",
                                Calc`M[Calc`d_current];
```

```
Comp := (assigned Calc`OpCntRed);
d := Calc'd_current;
i := Calc`i_current;
M := Calc'M[d];
N := Calc`N[d];
if i eq 1 then firstrow := 1;
else firstrow := Calc`lastrow[d][i-1] +1;
end if;
lastrow := Calc`lastrow[d][i];
Pivots := Calc'Pivots[d];
for row in [firstrow..lastrow] do
    s := Calc`SigsFromRows[d][row];
    LTinit := s'LTbefore;
    colinit := Calc'Columns[d][LTinit];
    col := colinit;
    leadingcol := colinit;
    donetop := (row eq Pivots[colinit]);
    donefull := donetop;
    while (col le N) and (not donefull) do
        if (M[row][col] ne 0) then
            if (IsDefined(Pivots,col)) and Pivots[col] lt row then
                    Append(~s'log_redlist,Pivots[col]);
                M := MtxAddRow(M,-M[row][col],Pivots[col],row);
                if Comp then
                                    nop := NumberOfNonZeroEntries(M[Pivots[col]])
                                    -1;
                                    Calc'OpCntRed +:= nop;
                    end if;
                elif not donetop then
                    leadingcol := col;
                    Pivots[leadingcol] := row;
                    donetop := true;
                    if not Full then
                                    donefull := true;
                    end if;
                end if;
            end if;
            col +:= 1;
        end while;
        try
            M := MtxMultiplyRow(M,1/M[row][leadingcol],row);
        catch e;
        end try;
        s`valafter := MatrixRowToPolynomial(M[row],Calc`Monomials[d]);
        s'LTafter := Calc`Monomials[d][leadingcol];
        if s'LTafter ne s'LTbefore then
            s`LTfromGB := s`LTafter;
        end if;
        Calc := SigUpdate(s,Calc);
```

```
    end for;
    Calc`Pivots[d] := Pivots;
    Calc`M[d] := M;
    PrintDebug(Sprintf("Result :\n%o\n",M));
    return Calc;
end function;
function F5_GBCriterion (s,Calc)
    // INPUT :
    // - s : signature Calc : state of the computation
    // OUTPUT :
    // - res : boolean stating if we should add the
    // polynomial associated to that signature to the Groebner
    // basis
    res := (((assigned s'LTafter)
        and (s`LTafter ne s'LTbefore))
        or ((Calc`i_current eq 1)
            and (s'mu eq 1)));
    return res;
end function;
function F5_ReduceGBMatrix(Calc : Full := false)
    // Reduce the Groebner basis in the matrix
    Comp := (assigned Calc`OpCntRed);
    for d in [1..Calc'dreg] do
        M := Calc`M[d];
        firstrow := 1;
        lastrow := NumberOfRows(M);
        N := NumberOfColumns(M);
        Pivots := Calc'Pivots[d];
        for row in [firstrow..lastrow] do
            reduced_p := false;
            s := Calc`SigsFromRows[d][row];
            if s'LTafter ne s'LTbefore then // Element of the GB
                col := Calc'Columns[d][s'LTafter]+1;
                while (col le N) do
                if (M[row][col] ne 0)
                    and IsDefined(Pivots, col)
                    and (Full or Pivots[col] gt row)
                        then
                        reduced_p := true;
                M := MtxAddRow(M,-M[row][col],Pivots[col],row);
                if Comp then
                                    nop:=NumberOfNonZeroEntries(M[Pivots[col]])
                                    -1;
                                    Calc'OpCntInterRed +:= nop;
                                    end if;
                end if;
```

```
                    col +:= 1;
    end while;
    // Update the value:
    s`valafter := MatrixRowToPolynomial(M[row],
                                    Calc`Monomials[d]);
        Calc := SigUpdate(s,Calc);
        if reduced_p then
            vprintf Groebner,2: "%o -> %0 -> %0\n",
                                    s`LTbefore, s`LTafter,
                                    LeadingTerm(s'valafter);
        end if;
            else // See if the original one was reduced
        sig := s`LTafter div s`LTfromGB;
        d0 := Degree(s`LTfromGB);
        row0 := Calc`Pivots[d0][Calc`Columns[d0][s`LTfromGB]];
        M[row] := PolynomialToMatrixRow(
        MatrixRowToPolynomial(Calc`M[d0][row0],
                                    Calc`Monomials[d0])
                            *sig,
                                Calc`Columns[d],N)[1];
            end if;
        end for;
        Calc`M[d] := M;
    end for; // for d
    return Calc;
end function;
// *** Main function
// Record containing information about the F5 run. See the
// initialization step in function F5_Mtx for explanations of the
// fields.
Calculation := recformat<F,
    m,
    n,
    dreg,
    D,
    quadratic,
    LTGB,
    NonCrit,
    M,
    Mtilde,
    Monomials,
    Columns,
    Crit,
    SigGB,
```

```
SigsFromMons,
SigsFromRows,
ReducedSigs,
Pivots,
N,
allMonomials,
lastrow,
i_current,
d_current,
row_current,
SigFromLT,
LTFromSig,
OpCntRed, OpCntInterRed, OpCntTot>;
```

function F5_Mtx (F : dreg := DReg(F),
bound := func<s | "">,
Full := false,
InterRed := false,
FullinterRed := Full and InterRed,
Comp := false)
// bound : function taking the signature of a row and
// printing various info about it
//
// Full : false -> no full reduction (only top)
// true -> full reduction (for rows non-triv. top-reduced)
//
// InterRed : true iff we want to reduce the intermediate
// Groebner bases
//
// FullinterRed : true iff we want to do this as a full
// reduction
//
// Bases after completing step i
// \{\{\{ *** Initialization
// \{\{\{ ****** Preparation of input
if exists(i)\{i : i in [1..\#F] | not IsHomogeneous(F[i])\} then
error Sprintf("Polynomial \#\%o is not homogeneous", i);
end if;
P := Parent(F[1]);
R := BaseRing(P);
$\mathrm{n}:=\operatorname{Rank}(\mathrm{P})$;
m := \#F;
D := [Degree(f) : f in F];
// \}\}\}
// \{\{\{ ****** Initialization of output
G := AssociativeArray([0..m]); // Groebner bases
G[0] := [];

```
S := []; // Signatures
// }}}
// {{{ ****** Initialization of internal values
Calc := rec<Calculation|>;
// We will store the current state of the calculation there
Calc'F := F;
Calc`m := m
Calc'n := n
Calc`dreg := dreg;
Calc`D := D
Calc`N := AssociativeArray([0..dreg]);
// Calc`N[d] = number of monomials at degree d
Calc`Monomials := AssociativeArray([0..dreg]);
// Calc'Monomials[d][i] = i'th monomial of degree d
Calc`Columns := AssociativeArray([0..dreg]);
// Calc'Columns[d][m] = index of the column of monomial m at
// degree d
Calc'Pivots := AssociativeArray([0..dreg]);
Calc'Crit := AssociativeArray([0..dreg]);
// Calc'Crit[d] = sequence of monomials at degree d serving
// in the F5 criterion
Calc‘lastrow := AssociativeArray([0..dreg]);
// Calc`lastrow[d][i] = last row reached at the step i at degree d
Calc`row_current := 1;
Calc`M := AssociativeArray([0..dreg]);
// Calc`M[d] = matrix built at degree d
Calc`SigsFromMons := AssociativeArray([1..m]);
Calc`SigsFromRows := AssociativeArray([0..dreg]);
// Calc`SigsFromMons[i][mu] is the sig associated to the pair
// i,mu
//
// Calc`SigsFromRows[d][l] is the sig associated to the row l
// at degree d
for d in [0..dreg] do
Calc`Monomials[d] := SortedMonomialsOfWeightedDegree(P,d);
Calc`N[d] := #(Calc`Monomials[d]);
Calc'Pivots[d] := [];
Calc`Columns[d] := AssociativeArray(Calc`Monomials[d]);
```

```
    for i := 1 to Calc'N[d] do
        Calc`Columns[d][Calc'Monomials[d][i]] := i;
    end for;
    Calc'Crit[d] := AssociativeArray([1..m]);
    Calc`lastrow[d] := [];
    Calc`M[d] := ZeroMatrix(R,0,Calc`N[d]);
    Calc`SigsFromRows[d] := AssociativeArray();
end for;
Calc`allMonomials := [];
for d in [0..dreg] do
    Calc'allMonomials cat:= Calc'Monomials[d];
end for;
if Comp then
    Calc'OpCntRed := 0;
    Calc`OpCntInterRed := 0;
    Calc'OpCntTot := 0;
end if;
// }}}
// }}}
// {{{ *** Computation
F := [RescalePoly(f) : f in F];
for i in [1..m] do
    // {{{ ****** Preparation
    Calc`i_current := i;
    fi := F[i]
    di := D[i];
    Calc`SigsFromMons[i] := AssociativeArray(Calc'allMonomials);
    for d in [1..dreg] do
        if i eq l then
            Calc`lastrow[d][i] := 0;
            Calc'Crit[d][i] := [];
        else
            Calc`lastrow[d][i] := Calc`lastrow[d][i-1];
            Calc'Crit[d][i] := Calc'Crit[d][i-1];
        end if;
    end for;
    // }}}
    for d in [di..dreg] do
        vprintf User2 : "%o i=%o, d=%o\n", Status(Calc), i, d;
        // {{{ ****** Construction of the matrix
        Calc`d_current := d;
```

```
    Calc'row_current := Calc'lastrow[d][i]+1;
    if d eq di then
        // {{{ Case of a new polynomial in the matrix
        s := SigCreate(i,fi,Calc'row_current);
        Calc := F5_InsertRow(s, Calc);
    // }}}
elif d gt di then
    // {{{ Case of a polynomial from a lower degree
    for k := Calc'N[d-di] to 1 by -1 do // Decreasing order
                mu := Calc`Monomials[d-di][k];
                PrintDebug(Sprintf("Considering signature %o",
                                    <i,mu>));
                s := rec <Sig| i := i, mu := mu, d := d>;
                if F5_Criterion(s, Calc) then
                PrintDebug("F5 criterion : OK");
                x := MinVariableDividing(mu);
                mu2 := mu/x;
                if IsDefined(Calc'SigsFromMons[i],mu2) then
                    // {{{ Construction and addition of
                    // the new poly
                    PrintDebug("Presence at "
                                    cat "previous degree : OK");
                    s2 := Calc`SigsFromMons[i][mu2];
                    s := SigMultiplyByVar(s2,x,
                                    Calc'row_current);
                    Calc := F5_InsertRow(s, Calc);
                    // }}}
                else
                    PrintDebug("Presence at "
                                    cat "previous degree : No");
                end if;
                else
                        PrintDebug("F5 criterion : No");
                end if;
            end for;
            // }}}
        end if;
        // }}}
        // {{{ ****** Reduction of the matrix
        Calc := F5_EchelonForm(Calc : Full := Full);
        // }}}
end for; // for d
// {{{ 2.3. Update of the GB
if InterRed then
    // "Reduce the Groebner basis", in matrix words
    Calc := F5_ReduceGBMatrix(Calc : Full := FullInterRed);
    // And recompute the Groebner basis
```

```
        jmin := 1;
        G[i] := [];
    else
        jmin := i;
        G[i] := G[i-1];
        end if;
        for j in [jmin..i] do
        for m in Keys(Calc`SigsFromMons[j]) do
            s := Calc`SigsFromMons[j][m];
            d := s`d;
        Append(~Calc`Crit[d][j],s`LTafter);
        if F5_GBCriterion(s,Calc) then
            g := s`valafter;
            Append(~G[i],g);
        end if;
        // }}}
        end for; // for m
    end for; // for j
end for; // for i
// }}}
// {{{ *** Conclusion
if Comp then
    Calc`OpCntTot := Calc`OpCntRed + Calc`OpCntInterRed;
    end if;
    return G, Calc;
    // }}}
end function;
// }}}
// Set of functions to help extracting information out of the Calc
// output of F5_Mtx
function FoldAlongGB (L)
    // Given a list of <s,list of whatever>, returns a array whose
    // entry m is <l,ls> where l is the cat of all lists associated to
    // a signature of leading term m and ls is the list of these
    // signatures.
    res := AssociativeArray();
    for k in Keys(L) do
        cpl := L[k];
        s := cpl[1];
        l := cpl[2];
        m := s`LTfromGB;
        if not IsDefined(res,m) then
            res[m] := <[],[]>;
        end if;
        rescpl := res[m];
```

```
        res[m] := <Append(rescpl[1],s),rescpl[2] cat cpl[2]>;
    end for;
    return res;
end function;
function ListReductionsDownwards (Calc)
    // Returns a list of <s,l> where s is a sig and l the list of
    // lines it reduced
        res := AssociativeArray();
        for i in [1..Calc`m] do
            SfM := Calc`SigsFromMons[i];
            SfL := Calc`SigsFromLines;
            for m in Keys(SfM) do
                    s := SfM[m];
                    L := s`log_redlist;
                    for line in L do
                    s2 := SfL[s`d][line];
                    i2 := s2`i;
                    mu2 := s2`mu;
                    if not IsDefined(res,<i2,mu2>) then
                                    res[<i2,mu2>] := <s2,[]>;
                    end if;
                    cpl := res[<i2,mu2>];
                    l2 := Append(cpl[2],s);
                    res[<i2,mu2>] := <cpl[1],l2>;
                    end for;
        end for;
    end for;
    return res;
end function;
```


## A.1.2. Weighted homogeneous systems

File Whomo.m

```
function TotalWeightedDegree (f)
    res := 0;
    for m in Monomials(f) do
            d := WeightedDegree(m);
            res := Max(res,d);
    end for;
    return res;
end function;
function Homogenize (F)
    if F eq [] then return [];
    else
        P := Parent(F[1]);
        n := Rank(P);
```

```
        ord := MonomialOrder(P);
        W := Grading(P);
        Wh := Append(W,1);
        if ord[1] eq "grevlexw" then
        nouvord := <"grevlexw",Wh>;
        else
        nouvord := ord;
    end if;
    R := CoefficientRing(P);
    Q := PolynomialRing(R,Wh,nouvord);
    Fh := [];
    for f in F do
        D := TotalWeightedDegree(f);
        fh := 0;
        M := Monomials(f);
        C := Coefficients(f);
        r := #M;
        for i in [1.. r] do
            e := Exponents(M[i]);
            d := &+[W[i] * e[i] : i in [1..n]];
            eh := Append(e,D - d);
            fh := fh + C[i] * Monomial(Q,eh);
        end for;
        Append(~Fh,fh);
    end for;
        return Fh,Q;
    end if;
end function;
function DeHomogenize (Fh,Ph,P)
    n := Rank(P);
    F := [];
    for f in Fh do
        C,M := CoefficientsAndMonomials(f);
        M2 := [];
        for i in [1..#M] do
            e := Exponents (M[i]);
            e2 := [];
            for k := 1 to n do // Remove the last variable
                e2[k] := e[k];
            end for;
            M2[i] := Monomial(P,e2);
        end for;
        f2 := Polynomial(C,M2);
        Append(~F,f2);
    end for;
    return F;
end function;
```

```
function WHomoToHomo (F)
    if F eq [] then return []; end if;
    P := Parent(F[1]);
    W := Grading(P);
    n := #W;
    R := CoefficientRing(P);
    Q := PolynomialRing(R,n,"grevlex");
    h := hom<P -> Q | [(Q.i)^W[i] : i in [1..n]]>;
    return h(F);
end function;
function FindWeights(F)
    // Find weights wl..wn such that all polynomials of F can be
    // written as polynomials in X1^w1..Xn^wn
    if F eq [] then error("Empty sequence"); end if;
    P := Parent(F[1]);
    n := Rank(P);
    W := [0 : i in [1..n]];
    for f in F do
        for m in Monomials(f) do
            e := Exponents(m);
            for i in [1..n] do
                W[i] := Gcd(W[i],e[i]);
            end for;
        end for;
    end for;
    return W;
end function;
function HomoToWHomo(F : W := [])
    if F eq [] then return []; end if;
    P := Parent(F[1]);
    n := Rank(P);
    if IsEmpty(W) then
        W := FindWeights(F);
    end if;
    assert #W eq n;
    Q := PolynomialRing(CoefficientRing(P),W);
    F2 := [];
    for i in [1..#F] do
        f := F[i];
        f2 := 0;
        C,M := CoefficientsAndMonomials(f);
        for j in [1..#C] do
            c := C[j];
            e := Exponents(M[j]);
            e2 := [e[k] div W[k] : k in [1..n]];
            m2 := Monomial(Q,e2);
            f2 +:= c*m2;
```

```
        end for;
        Append(~F2,f2);
    end for;
    return F2;
end function;
```

File tests_GB.m
load "F5Mtx.m";
load "Whomo.m";

SetVerbose("Groebner",1);
SetVerbose("User2",1);

```
W := [4,2,1,1];
P<X,Y,Z,T> := PolynomialRing(GF(65521),W);
F := [
    11541*X^2 + 53990*X*Y^2 + 37592*Y^4 + 27834*X*Y*Z^2 +
    30004*Y^3*Z^2 + 8314*X*Z^4 + 4825*Y^2*Z^4 + 28158*Y*Z^6 +
    19050*Z^8 + 26015*X*Y*Z*T + 557*Y^3*Z*T + 58180*X*Z^3*T +
    42158*Y^2* 2*^3*T + 34612*Y*Z^5*T + 45718*Z^7*T + 34908*X*Y*T^2
    + 23734*Y^3*T^2 + 57434*X*Z^2*T^2 + 28743*Y^2*Z^2*T^2 +
    18580*Y*Z^4*T^2 + 10813*Z^^ 6*T^2 + 23568*X*Z*T^3 +
    65399* *^2*Z*T^3 + 3761*Y*Z^3*T^3 + 51403*Z^5*T` 3 +
    20744*X*T^4 + 46116*Y`2* 2*^4 + 44410* Y*Z^2*T^4 + 37156*Z^4*T^4
    + 13625*Y*Z*T^5 + 24856*Z^3*T^5 + 39692*Y*T^6 + 34019*Z^2*T^6
    + 53446*Z*T^7 + 1929*T^8,
    33098*\mp@subsup{X}{}{\wedge}2 + 54375*X*Y^2 + 14826*Y^4 + 3377*X*Y*Z^2 +
    23825*\mp@subsup{Y}{}{\wedge}3*\mp@subsup{Z}{}{\wedge}2 + 7023*X*Z^4 + 52920*Y^2* *^4 + 23589*Y*Z^^6 +
    50207*Z^8 + 49154*X*Y*Z*T + 3423*Y^3*Z*T + 225*X*Z^3*T +
    43216*Y^2*ZZ^3*T + 56563*Y*Z^5*T + 8367*Z^^7*T + 64946*X*Y*T^2
    + 60800*Y^3*T^2 + 65240*X*Z^2 2*T^2 + 34223*Y^2*Z`2*T^2 +
    39536*Y*Z^4*T^2 + 15290*Z^6*T^2 + 33901*X*Z*T^3 +
    21724*Y^2*Z*T^3 + 1521*Y*Z^3*T^3 + 32997*Z^5*T^^3 +
    11568*X*T^4 + 35339*Y^2*T^4 + 39002*Y*Z^}2**\mp@subsup{T}{}{\wedge}4+6645*\mp@subsup{Z}{}{\wedge}4*\mp@subsup{T}{}{\wedge}
    + 38351*Y*Z*T^5 + 52984*Z^3*T^5 + 50226*Y*T^6 + 12349*Z^2*T^6
    + 50105*Z*T^7 + 8711*T^8,
    35779* X^2 + 7522* X*Y^2 + 49129* Y^4 + 6315*X*Y*Z^2 +
    33104*Y^3*Z^2 + 59779*X*Z^4 + 47533*Y^2*Z^4 + 14703*Y*Z^6 +
    51257*Z^8 + 11944*X*Y*Z*T + 27755*Y^3*Z*T + 55382*X*Z^3*T +
    1529*Y^2*Z^3*T + 24962*Y*Z^5*T + 6285*Z^7*T + 62689*X*Y*T^2 +
    61121*Y^3*T^2 + 2430*X*ZZ^2*T^2 + 52833*Y`2*Z^2*T^2 +
    48100*Y*Z^4*T^2 + 44590*Z^6*T^2 + 59954*X*Z*T^3 +
    43520*Y`2 2*Z*T^3 + 16159*Y*Z^3*T^3 + 112*Z^5*T^3 + 48989*X*T^4
    + 60810*Y^2*T^4 + 51752*Y*Z^^2*T^4 + 54345*Z^4*T^4 +
    12978*Y*Z*T^5 + 14453*Z^3*T^5 + 20513*Y*T^6 + 46414*Z^2*T^6 +
    40677*Z*T^7 + 35597*T^8,
```

```
    59914*X^2 + 5160*X*Y^2 + 51256*Y^4 + 688*X*Y*Z^2 +
    37687*Y^3*Z^2 + 7159*X*Z^4 + 60215*Y^2*Z^4 + 39922*Y*Z^6 +
    11761*Z^8 + 22386*X*Y*Z*T + 45230*Y^3*Z*T + 28100*X*Z^3*T +
    27809*Y^2*Z^3*T + 54179*Y*Z^5*T + 18563*Z^7*T + 38052*X*Y*T^2
    +64357*Y^3**^2 + 26800*X*Z^2*T^2 + 37661*Y^2*Z^2*T^2 +
    57333*Y*Z^4*T^2 + 37124*Z^6*T^2 + 18538*X*Z*T^3 +
    29552*Y^2*Z*T^3 + 44263*Y*Z^3*T^3 + 17949*Z^5*T^3 +
    55615*X*T^4 + 32246*Y^2*T^4 + 61559*Y*Z^2*T^4 + 46600*Z^4*T^4
    + 59196*Y*Z**`5 + 19035*Z^3*T^5 + 43459*Y*T^6 + 60433*Z^2*T^6
    + 10492*Z*T^7 + 57095*T^8
];
FH := WHomoToHomo(F);
//print FH;
assert IsHomogeneous(FH);
FF1 := HomoToWHomo(FH);
//print FF1;
assert [P!f : f in FF1] eq F;
FF2 := HomoToWHomo(FH : W := W);
//print FF2;
assert [P!f : f in FF2] eq F;
G1 := HomoToWHomo(GroebnerBasis(FH) : W := W);
G1 := [P!f : f in G1];
GG2 := F5_Mtx(F);
G2 := GG2[4];
SetVerbose("Groebner",0);
assert IsGroebner(G1);
assert IsGroebner(G2);
assert Ideal(G1) eq Ideal(F);
assert Ideal(G2) eq Ideal(F);
```


## A.2. Maple source code for determinantal varieties (Chapter 4)

## A.2.1. Algorithms

File functions.mpl
\# read "functions.mpl";
with(LinearAlgebra):
with (FGb) :
with(VectorCalculus):
with(combinat):
randomize():

```
randval := rand(100..10000):
Minors := proc(s,K)
description "Input: s integer, K matrix"
    "Output: the list of (s X s) minors of K";
local d1,d2,MIN,j,l,listofminors,subm:
    (d1,d2):=Dimension(K):
    if s > min(d1,d2) then
        printf("ERROR: INPUT s TOO LARGE\n");
        return;
    fi;
    MIN:=Matrix(binomial(d1,s),binomial(d2,s)):
    for j from l to binomial(d1,s) do
        for l from 1 to binomial(d2,s) do
            subm := SubMatrix(K,(choose([seq(i,i=1..d1)],s))[j],
                (choose([seq(i,i=1..d2)],s))[l]);
            MIN[j,l]:=Determinant(subm):
        od:
    od:
    listofminors:=seq(seq(MIN[j,l],l=1..binomial(d2,s)),
                    j=1..binomial(d1,s)):
    return(listofminors):
end:
dimension := proc(sys, vars := indets(sys))
description "Find the dimension of the variety defined by the system "
    "by cutting it with random hyperplanes";
    local dim, hyp, randgen,hh, v, gg,sysCut,j, firstdone;
    dim := nops(vars):
    hyp := []:
    randgen := rand(-500..500):
    for j in seq(i,i=1..dim) do
        hh := randgen():
        for v in vars do
        hh := hh + randgen()*v:
        od;
        hyp := [op(hyp),hh]:
    od:
    sysCut := [op(sys),op(hyp)]:
    gg := [1]:
    firstdone := false;
    while gg = [1] and dim >= 0 do
        if firstdone then
```

```
            dim := dim -1:
            sysCut := sysCut[1..-2]:
        else
            firstdone := true:
        fi:
        printf("Trying dimension %d\n",dim):
        gg := fgb_gbasis(sysCut,65521,vars,[],{"index"=2000000}):
    od:
    return dim:
end:
sqfr := proc(f)
    description "Square-free reduction of f";
    return mul(ff,ff in map(x -> x[1],factors(f)[2])):
end:
GB_interpolate := proc(sys,charac,vars1,vars2,varinter,deg)
description "Use interpolation on varinter up to degree deg to compute "
    "a Groebner basis of sys for the order eliminating vars1";
local points, bases, v, i, gb, res, t0, safety, tgtdeg, m,
    bases_i, pol, frct;
    v := -1:
    points := []:
    bases := []:
    safety := 1:
    tgtdeg := 2*deg + 2 + 2*safety:
    for i from 1 to tgtdeg do
        printf("GB %a/%a: ", i,tgtdeg):
        t0 := time():
        while v = -1 or v in points do
            v := randval():
        od:
        gb := fgb_gbasis_elim(eval(sys,varinter=v),charac,vars1,vars2,
                            {"index"=20000000}):
        points := [op(points),v]:
        bases := [op(bases),gb]:
        printf("done (%a) [v=%a]\n", time()-t0,v):
    od:
    res := []:
    m := mul(varinter-v,v in points):
    for i from 1 to nops(bases[1]) do
        bases_i := map(x -> x[i]/lcoeff(x[i])
            , bases):
        pol := CurveFitting[PolynomialInterpolation](points,
                                    bases_i,
                                    varinter):
        frct := ratrecon(pol,m,varinter,deg+safety,deg+safety):
```

```
        res := [op(res),numer(frct)]:
    od:
    return res:
end:
GB_find_deg := proc(sys,vars1,vars2,varinter,ntrials := 10)
description "Find the maximal degree in variable varinter of "
    "the polynomials in a Groebner basis of sys eliminating vars1";
local deg, gb, ss, i, v, t0;
    deg := -1:
    printf("Finding degree (%a tests): ", ntrials):
    i := 0:
    facts := []:
    while i < ntrials do
        ss := {}:
        for v in vars2 do
            if v <> varinter then
                ss := {op(ss), v=randval()}:
            fi:
        od:
        gb := fgb_gbasis_elim(eval(sys,ss),0,vars1,vars2):
        for j from 1 to nops(gb) do
            if i = 0 then
                facts := [op(facts),0]:
            fi:
            facts[j] := gcd(facts[j],gb[j]);
        od:
        deg := max(map(degree,gb,varinter),deg):
        printf("|"):
        i := i+1:
    od:
    printf("\n"):
    facts_prod := sqfr(mul(p,p in facts));
    return deg, facts_prod:
end:
Elimination_codim1 := proc(F,vars,params,{algo:=default})
description "Compute a polynomial whose zeroes cover the "
    "projection of V(F) on the parameter space"
    ""
    "Compute a system of generators for the elimination ideal "
    "of F obtained by eliminated the variables vars"
    "Uses the algorithm 'algo' to compute the basis. Admissible "
    "values for 'algo' are:"
    "- gb_direct : compute an elimination Groebner basis directly "
    "- gb_interp : compute an elimination Groebner basis using "
    "evaluation/interpolation"
    "Any other value defaults to gb_interp."
```

```
    ;
local var,deg_interp,G:
    if algo = gb_direct then
                ### Direct
                G := fgb_gbasis_elim(F,0,vars,params ,{"verb"=3}):
                return G[1];
    else
        ### Interpolation with first param
        var := params[l]:
        deg_interp, facts := GB_find_deg(F,vars,params,var):
        G := GB_interpolate(F,0,vars,params,var,deg_interp):
        return facts*G[1];
    fi:
end:
```

DiscriminatingPolynomial := proc(M,r0, H, vars, params,
eqs $:=[]$, neqs $:=[]$, \{algo:=default\})
description "Compute a polynomial P as in section 3."
" "
"The argument ‘eqs‘ is a set of extra equations and inequations "
"restricting the solutions: the solutions returned are "
"projections in the parameter space of points at which all "
"polynomials of eqs vanish."
""
"The polynomials in 'neqs' are saturated in the computations."
""
"See Elimination for the parameter 'algo‘":
local res,res1, res2,res3,t0,t1,t2,t3:
t0 := time():
res1 := RankExactly(M,r0,vars,params,eqs,neqs,
':-algo'=algo):
t1 := time():
res2 := DeterminantCritVals(M,r0,vars,params,eqs,neqs,true,
':-algo'=algo) :
t2 := time():
res3 := DeterminantBoundary(M,r0,H,vars,params,eqs,neqs,true,
':-algo'=algo):
t3 := time():
printf("RankExactly \t\%as $\backslash n ", t 1-t 0):$
printf("DeterminantCritVals $\backslash t \% a s \backslash n ", t 2-t 1):$
printf("DeterminantBoundary $\backslash t \% a s \backslash n ", t 3-t 2):$
return res1*res2*res3, res1, res2, res3;
end:
IncidenceVariety := proc(M,r)
description "Compute a system of generators for the incidence variety "
"of rank $r$ of M"
;

```
local MatrixU, MatrixY,k,i,j,Prod1,Sys1,Prod2,Sys2,Sys:
    k := RowDimension(M):
    MatrixU := Matrix(k-r,k):
    for i from 1 to k-r do
        for j from 1 to k do
            MatrixU[i,j] := randval():
        od:
    od:
    print(MatrixU);
    MatrixY := Matrix(k,k-r):
    for i from 1 to k do
        for j from 1 to k-r do
            MatrixY[i,j] := Y[i,j]:
        od:
    od:
    Prod1 := M . MatrixY:
    Sys1 := [seq(seq(Prod1[i][j],j=1..k-r),i=1..k)]:
    Prod2 := MatrixU . MatrixY - Matrix(k-r,k-r,shape=identity):
    Sys2 := [seq(seq(Prod2[i][j],j=1..k-r),i=1..k-r)]:
    Sys := [op(Sys1),op(Sys2)]:
    return Sys, indets(MatrixY):
end:
RankExactly := proc(M,r0,vars,params,eqs := [], neqs := [],
                    {algo := default})
description "Algorithm RankExactly (section 3.4)"
    ""
    "See DiscriminatingPolynomial for a description of parameters "
    "'eqs' and 'neqs`."
    "See Elimination for a description of the parameter 'algo"";
local res,k,n,FVr,dim,JVr,FV,F0,varsY,Mins,i,Sysi,G, codim:
    k := RowDimension(M) :
    n := nops(vars):
    res := 1:
    FVr := [Minors(r0+1,M)]:
    codim := (k-r0)^2:
    JVr := Jacobian(FVr,[op(vars)]):
    FV := [op(FVr), Minors(codim,JVr),op(eqs),u*mul(f, f in neqs)-1]:
    F0, varsY := IncidenceVariety(M, r0):
    Mins := [Minors(r0,M)]:
    for i from 1 to nops(Mins) do
        Sysi := [op(FV),op(F0),op(Mins[1..i-1]),uu*Mins[i]-1]:
        g := Elimination_codim1(Sysi,[uu,u,op(varsY),op(vars)],params,
                        ':-algo'=algo):
        if g <> FAIL then
```

```
            res := res*g:
        fi:
    od:
    printf("RankExactly: %a\n", factor(res));
    return res:
end:
DeterminantCritVals := proc(M,r0,vars,params,eqs := [], neqs := [],
                            skipRankExactly := false, {algo:=default})
    description "Algorithm DeterminantCritVals (section 3.5)"
    ""
    "See DiscriminatingPolynomial for a description of parameters "
    "‘eqs` and 'neqs`."
    "See Elimination for a description of the parameter 'algo"":
local res,F0,varsY,N,J,F1,G:
    if not skipRankExactly then
        res := RankExactly(M,r0,vars,params):
    else
        res := 1:
    fi:
    F0, varsY := IncidenceVariety(M,r0-1):
    N := nops(F0):
    J := Jacobian(F0,[op(varsY),op(vars)]):
    F1 := [op(F0),Minors(N,J),op(eqs),u*mul(f, f in neqs)-1]:
    g := Elimination_codim1(F1,[u,op(varsY),op(vars)],params,
        ':-algo'=algo):
    res := res*g:
    printf("CritVals: %a\n", factor(res));
    return res:
end:
DeterminantBoundary := proc(M,r0,H,vars,params,
                    eqs := [], neqs := [],
                    skipRankExactly := false,
                            {algo := default})
    description "Algorithm DeterminantBoundary (section 3.6)"
    ""
    "See DiscriminatingPolynomial for a description of parameters "
    "'eqs` and 'neqs`."
    "See Elimination for a description of the parameter 'algo"";
local res, F0, varsY, h, F1, G:
    if not skipRankExactly then
        res := RankExactly(M,r0,vars,params):
    else
        res := 1:
    fi:
    F0, varsY := IncidenceVariety(M,r0-1):
    for h in H do
        F1 := [op(F0),h,op(eqs),u*mul(f, f in neqs)-1]:
```

```
        g := Elimination_codim1(F1,[u,op(varsY),op(vars)],params,
                ':-algo'=algo):
        res := res * g:
    od:
    printf("Boundary: %a\n", factor(res));
    return res:
end:
```


## A.2.2. Case of water

## Definitions

File water_def.mpl

```
# read "water_def.mpl":
d1:=g1-G1:d2:=g2-G2:
MatrixD:=Matrix([[-G1*y1, -z1-1,-G1+d1*z1,2*d1*y1],
    [-g1*z1,y1,d1*y1,(G1-d1)-2*d1*z1],
    [-G2*y2, -z2-1,-G2+d2*z2,2*d2*y2],
    [-g2*z2,y2,d2*y2,(G2-d2)-2*d2*z2]]):
g1 := 1: G1 := 1:
vars := [y1,y2,z1,z2]:
params := [g2,G2]:
H := [1-y1^2-(z1+1)^2, 1-y2^2-(z2+1)^2]:
k := 4: n := 4: t := 2;
DetD:=Determinant(MatrixD);
```


## Computations using the functions

File water_functions.mpl

```
# read "water_functions.mpl":
# Function definitions
read "functions.mpl";
# Definitions for water
read "water_def.mpl";
# Computations
```

```
t0 := time():
res_general,res1,res2,res3 :=
    DiscriminatingPolynomial(MatrixD,3,H,vars,params,[],[],
                                    algo = gb_direct):
t1 := time():
# Results
# RankExactly
ff1 := map(x -> x[1],factors(res1)[2]):
nops(ff1);
# 6
map(degree,ff1);
# [1,1,1,2,1,1]
# Critical values
ff2 := map(x -> x[1],factors(res2)[2]):
nops(ff2);
# 5
map(degree,ff2);
# [1, 1, 1, 1, 1]
# Boundary
ff3 := map(x -> x[1],factors(res3)[2]):
nops(ff3);
# 7
map(degree,ff3);
# [1, 1, 1, 1, 2, 2, 3]
allfacts := [op({op(ff1),op(ff2),op(ff3)})]:
nops(allfacts);
# 13
map(degree,allfacts);
# [1, 1, 1, 1, 1, 1, 1, 1, 1, 2, 2, 2, 3]
filename := "water_results.mpl";
fid := fopen(filename,WRITE):
fprintf(fid,"## RankExactly:\n##---------\n"):
fprintf(fid,"res1 := [\n"):
for i from 1 to nops(ff1)-1 do
    f := ff1[i]:
    fprintf(fid," %a,\n", f):
od:
fprintf(fid," %a\n]:\n\n",ff1[-1]):
fprintf(fid,"## CritVals:\n##---------\n"):
fprintf(fid,"res2 := [\n"):
for i from 1 to nops(ff2)-1 do
```

```
    f := ff2[i]:
    fprintf(fid," %a,\n", f):
od:
fprintf(fid," %a\n]:\n\n",ff2[-1]):
fprintf(fid,"## Boundary\n##---------\n"):
fprintf(fid,"res3 := [\n"):
for i from 1 to nops(ff3)-1 do
    f := ff3[i]:
    fprintf(fid," %a,\n", f):
od:
fprintf(fid," %a\n]:\n\n",ff3[-1]):
fprintf(fid,"## All (without duplicates)\n##---------\n"):
fprintf(fid,"allfacts := [\n"):
for i from 1 to nops(allfacts)-1 do
    f := allfacts[i]:
    fprintf(fid," %a,\n", f):
od:
fprintf(fid," %a\n]:\n\n",allfacts[-1]):
fclose(fid):
```


## Step-by-step computations

```
File water_computations.mpl
```

```
# read "water_computations.mpl";
```


# read "water_computations.mpl";

read "functions.mpl";

# System

read "water_def.mpl":
DetD:=Determinant(MatrixD):
gradD_yz := [diff(DetD, y1), diff(DetD, z1),
diff(DetD, y2), diff(DetD, z2)]:
r0 := 3:
Sys := [DetD,op(gradD_yz)]:

# Incidence varieties

F2,varsY2 := IncidenceVariety(MatrixD,2):
F3,varsY3 := IncidenceVariety(MatrixD,3):

# Singular points and critical values

# Rank(M) <= 2

```
```

Jac := Jacobian(F2,[z1,z2,y1,y2,op(varsY2)]):
Sys_rk2 := [op(F2),Minors(2*k+(k-r0+1)^2,Jac),u*g2-1]:
GB_rk2 := fgb_gbasis_elim(Sys_rk2,0,
[u,op(varsY2),z1,z2,y1,y2],[G2,g2]):
lprint(factor(GB_rk2));
(* Results:
[(G2-1)^2*(-2*g2+1+G2)^2*(-g2-1+2*G2)^2*(2*G2-g2)^2]
*)

# Intersection with the boundary

# Side 1

Sys_bnd1 := [op(F2),H[1]]:
GB_bnd1 := fgb_gbasis_elim(Sys_bnd1,0,
[op(varsY2),z1,z2,y1,y2],[G2,g2]):
lprint(factor(GB_bnd1));
(* Results:
[g2*(G2-1)^2*(-2*g2-1+3*G2)*(3*G2^2-5*G2*g2+g2^2+2*G2-2*g2+1)]
*)

# Side 2

Sys_bnd2 := [op(F2),H[2]]:
GB_bnd2 := fgb_gbasis_elim(Sys_bnd2,0,
[op(varsY2),z1,z2,y1,y2],[G2,g2]):
lprint(factor(GB_bnd2));
(* Results:
[(G2-1)^2*(2*G2^2-5*G2*g2+2*g2^2-2*G2+3*g2)\
*(2*G2^3-3*G2^}2*\textrm{g}2-3*\textrm{G}2*\textrm{g}2^^2+2*g2^ 3+2*G2^^2+9*G2*g2-11*g2^2-4*G2+6*g2)\
*(-g2+2*G2)^2]
*)

# Rank(M) = 3

allGB_rk3 := []:
MM := [Minors(k-1,MatrixD)]:
Sys_rk3 := [DetD,op(gradD_yz),op(F3),u1*g2-1]:
for i from 1 to k^2 do
Sys_rk3_i := [op(Sys_rk3),op(MM[1..i-1]),u2*MM[i]-1]:
GB := fgb_gbasis_elim(Sys_rk3_i,0,
[u1,u2,op(varsY3),z1,z2,y1,y2],
[G2,g2],

```
```

            {"verb"=3}):
    allGB_rk3 := [op(allGB_rk3),[i,GB]]:
    od:
for g in allGB_rk3 do
if g[2] <> [1] then
printf("i%a Basis:%a\n",g[1],factor(g[2])):
fi:
od:
(* Results:
i=1 Basis:[(-g2+2*G2)*(g2-2+G2)*(2*G2^2-5*G2*g2+2*g2^2+1)]
*)

```

\section*{Cylindrical algebraic decomposition}

File water_cad.mpl
```


# restart; read "water_cad.mpl";

read "water_def.mpl";
SDetD := [DetD,diff(DetD,y1),diff(DetD,y2),
diff(DetD,z1),diff(DetD,z2)]:
Pols := [g2,
G2,
g2-2*G2,
G2-1,
3*G2-1-2*g2,
3*G2^2-5*G2*g2+g2^2+2*G2-2*g2+1,
2*G2^2-5*G2*g2+2*g2^2-2*G2+3*g2,
(2*G2^3-3*G2^2*g2-3*G2*g2^2+2*g2^3
+2*G2^2+9*G2*g2-11*g2^2-4*G2+6*g2),
-2*g2+G2+1,
2*G2-1-g2,
G2-2+g2,
2*G2^2-5*G2*g2+2*g2^2+1];
xi := product(Pols[i], i = 1 .. nops(Pols));
with(RegularChains); with(ChainTools); with(SemiAlgebraicSetTools);
R := PolynomialRing([G2,g2]);
cadfull := CylindricalAlgebraicDecompose(Pols,R,output=cadcell):
sols := [];
for j from 1 to nops(cadfull) do
sols := [op(sols),
subs(op(2,op(1,subs(op(2,op(1,cadfull[j])),SamplePoint))),
box_bwe)[1]
]:
od:

```
```

print(nops(sols));
print(sols[1]);
avg2 := proc(l):
return (l[1]+l[2])/2;
end:
insideBB := x -> (evala(eval(y1^2+(z1+1)^2-1, x)) <= 0 and
evala(eval(y2^2+(z2+1)^^2-1, x)) <= 0):
PolsRed := remove(x -> x = g2 or x = G2 or x = 2*G2-g2, Pols):
sAvg := map(x -> [g2 = avg2(eval(g2,x)),
G2 = avg2(eval(G2,x))],
sols):
sAvgValid := select(x ->(eval(xi,x) <> 0
and 0 < eval(g2,x)
and 0 < eval(G2,x)
and eval(g2,x) < 2*eval(G2,x)),
sAvg):
print(nops(%)):

# 570

Sols := [seq(select(x -> x <> {y1 = 0, y2 = 0,
z1 = -1, z2 = -1},
[solve(Groebner[Basis](eval(SDetD,
sAvgValid[i]),
plex(y1, y2, z1, z2)))]),
i = 1..nops(sAvgValid))]:
print(map(nops, Sols));

# [2..2]

SolsInBB := [seq(select(x -> (evalb(0 <= evala(eval(y1^2, x)))
and evalb(0 <= evala(eval(y2^2, x)))
and insideBB(x)), Sols[i]),
i = 1 .. nops(Sols))]:
print(nops(select(x -> x <> [], SolsInBB)));

# 187

SolsBBSym := [seq(select(x -> (evalb(evala(eval(y1, x)) = 0)
and evalb(evala(eval(y2, x)) = 0)

```
```

        and insideBB(x)), Sols[i]),
    i = 1 .. nops(Sols))]:
    print(nops(select(x -> ( x <> []), SolsBBSym)));

# 156

IndBBSym := select(x -> ( SolsBBSym[x] <> []),
[seq(i, i = 1 .. nops(SolsBBSym))]):
SolsBBNoSym := [seq(select(x -> (evalb(0 < evala(eval(y1^2, x)))
and evalb(0 < evala(eval(y2^2, x)))
and insideBB(x)), Sols[i]),
i = 1 .. nops(Sols))]:
print(nops(select(x -> ( x <> []), SolsBBNoSym)));

# 31

IndBBNoSym := select(x -> ( SolsBBNoSym[x] <> []),
[seq(i, i = 1 .. nops(SolsBBNoSym))]):
g2G2_1 := sAvgValid[IndBBSym]:
g2G2_2 := sAvgValid[IndBBNoSym]:
filename1 := "pts_cad_1.txt";
fid1 := fopen(filename1,WRITE);
for s in g2G2_1 do
fprintf(fid1,"%a %a\n",evalf(eval(g2,s)), evalf(eval(G2,s)));
od:
fclose(fid1);
filename2 := "pts_cad_2.txt";
fid2 := fopen(filename2,WRITE);
for s in g2G2_2 do
fprintf(fid2,"%a %a\n",evalf(eval(g2,s)), evalf(eval(G2,s)));
od:
fclose(fid2);
f0,f1,f2,f3,f4,f5,f6,f7,f8,f9 := op(Pols[3..-1]):
test_all := x -> map(f -> evalb(eval(f,x) > 0),
[f1,f2,f3,f4,f5,f6,f7,f8,f9]);
crit_1 := x -> ((eval(f1,x) > 0
and eval(f4,x) > 0
and eval(f2,x) < 0)
or (eval(f1,x) < 0

```
```

    and eval(f2,x) > 0));
    print({op(map(crit_1,g2G2_1))});
print(nops(g2G2_1));

# 156

print(nops(select(crit_1,sAvgValid)));

# 156

crit_2 := x -> ((eval(f1,x) < 0
and eval(f6,x) > 0
and eval(f3,x) < 0)
or (eval(f1,x) > 0
and eval(f6,x) < 0
and eval(f5,x) > 0)):
print({op(map(crit_2,g2G2_2))});
print(nops(g2G2_2));

# 31

print(nops(select(crit_2,sAvgValid)));

# 31

```

\section*{A.2.3. General case}

\section*{Definitions}

File general_def.mpl
```


# read "general_def.mpl";

d1:=g1-G1:d2:=g2-G2:
MatrixD:=Matrix([[-G1*y1, -z1-1,-G1+d1*z1,2*dl*y1],
[-gl*z1,y1,dl*y1,(G1-d1)-2*d1*zl],
[-G2*y2, -z2-1,-G2+d2*z2,2*d2*y2],
[-g2*z2,y2,d2*y2,(G2-d2)-2*d2*z2]]):

```
g1 := 1:
vars := \([y 1, y 2, z 1, z 2]:\)
params := [G1,g2,G2]:
\(H:=\left[1-y 1^{\wedge} 2-(z 1+1)^{\wedge} 2,1-y 2^{\wedge} 2-(z 2+1)^{\wedge} 2\right]:\)
\(\mathrm{k}:=4: \mathrm{n}:=4: \mathrm{t}:=3\) :

DetD:=Determinant (MatrixD);

\section*{Using the functions}

File general_functions.mpl
```


# read "general_functions.mpl":

# Function definitions

read "functions.mpl";

# Definitions

read "general_def.mpl";

# Computations

t0 := time():
res_general,res1,res2,res3 :=
DiscriminatingPolynomial(MatrixD,3,H,vars,params,[],[]):
t1 := time():

# Results

# RankExactly

ff1 := map(x -> x[1],factors(res1)[2]):
nops(ff1);

# 6

map(degree,ff1);

# [1,1,1,1,2,1]

# Critical values

ff2 := map(x -> x[1],factors(res2)[2]):
nops(ff2);

# 7

map(degree,ff2);

# [1, 1, 1, 14, 1, 1, 1]

# Boundary

ff3 := map(x -> x[1],factors(res3)[2]):
nops(ff3)

# 7

map(degree,ff3);

# [2, 9, 9, 2, 1, 1, 1]

allfacts := [op({op(ff1),op(ff2),op(ff3)})]:
nops(allfacts);

# 15

map(degree,allfacts);

# [1, 1, 1, 1, 1, 1, 1, 1, 1, 2, 2, 2, 9, 9, 14]

```
```

filename := "general_results.mpl";
fid := fopen(filename,WRITE):
fprintf(fid,"\#\# RankExactly:\n\#\#---------\n"):
fprintf(fid,"res1 := [\n"):
for i from 1 to nops(ff1)-1 do
f := ffl[i]:
fprintf(fid," %a,\n", f):
od:
fprintf(fid," %a\n]:\n\n",ff1[-1]):
fprintf(fid,"\#\# CritVals:\n\#\#---------\n"):
fprintf(fid,"res2 := [\n"):
for i from 1 to nops(ff2)-1 do
f := ff2[i]:
fprintf(fid," %a,\n", f):
od:
fprintf(fid," %a\n]:\n\n",ff2[-1]):
fprintf(fid,"\#\# Boundary\n\#\#---------\n"):
fprintf(fid,"res3 := [\n"):
for i from 1 to nops(ff3)-1 do
f := ff3[i]:
fprintf(fid," %a,\n", f):
od:
fprintf(fid," %a\n]:\n\n",ff3[-1]):
fprintf(fid,"\#\# All (without duplicates)\n\#\#---------\n"):
fprintf(fid,"allfacts := [\n"):
for i from 1 to nops(allfacts)-1 do
f := allfacts[i]:
fprintf(fid," %a,\n", f):
od:
fprintf(fid," %a\n]:\n\n",allfacts[-1]):
fclose(fid):

```

\section*{Appendix B}

\section*{Bibliography}
[Agnoz] G. Agnarsson. "On the Sylvester denumerants for general restricted partitions". In: Proceedings of the Thirty-third Southeastern International Conference on Combinatorics, Graph Theory and Computing (Boca Raton, FL, 2002). Vol. 154. 2002, pp. 49-60 (cit. on p. 31).
[Alfo5] J. L. R. Alfonsín. The Diophantine Frobenius Problem. Oxford Lecture Series in Mathematics and Its Applications. Oxford: Oxford University Press, 2005 (cit. on pp. 31, 32, 95).
[Alf98] J. L. R. Alfonsin. "On variations of the subset sum problem". In: Discrete Applied Mathematics 81.1-3 (1998), pp. 1-7. ISSN: 0166-218X. DoI: 10.1016/So166-218X(96) ooro5-9 (cit. on p. 95).
[ALM99] P. Aubry, D. Lazard, and M. Moreno Maza. "On the theories of triangular sets". In: Эournal of Symbolic Computation, Special Issue on Polynomial Elimination 28 (1999), pp. 105-124. DOI: 10.1006/jsco.1999.0269 (cit. on p. 10).
[Bar+05] M. Bardet, J.-C. Faugère, B. Salvy, and B.-Y. Yang. "Asymptotic Behaviour of the Degree of Regularity of Semi-Regular Polynomial Systems". In: MEGA'05, 2005. Eighth International Symposium on Effective Methods in Algebraic Geometry. 2005 (cit. on pp. 88, 91, 114).
[Baro4] M. Bardet. "Étude des systèmes algébriques surdéterminés. Applications aux codes correcteurs et à la cryptographie". Français. PhD thesis. Université Pierre et Marie Curie - Paris VI, Dec. 2004. URL: http://tel.archives-ouvertes.fr/tel-oo449609 (cit. on pp. 34, 35, 43, 49, 74, 91, 114-116).
[BCR98] J. Bochnak, M. Coste, and M.-F. Roy. Real algebraic geometry. Vol. 36. Ergebnisse der Mathematik und ihrer Grenzgebiete (3) [Results in Mathematics and Related Areas (3)]. Translated from the 1987 French original, Revised by the authors. Springer-Verlag, Berlin, 1998, pp. x+430. ISBN: 3-540-64663-9. DOI: \(10.1007 / 978-3-\) 662-03718-8 (cit. on p. 139).
[BDo7] C. W. Brown and J. H. Davenport. "The complexity of quantifier elimination and cylindrical algebraic decomposition". In: ISSAC 2007. ACM, New York, 2007, pp. 54-60. DOI: \(10.1145 / 1277548.1277557\) (cit. on p. 135).
[BFSo4] M. Bardet, J.-C. Faugère, and B. Salvy. "On the complexity of Gröbner basis computation of semi-regular overdetermined algebraic equations". In: International Conference on Polynomial System Solving - ICPSS. Paris, France, Nov. 2004, pp. 71-75. URL: http://www-salsa.lip6.fr/ jcff/Papers/43BF.pdf (cit. on p. 17).
[BFS14] M. Bardet, J.-C. Faugère, and B. Salvy. "On the Complexity of the F5 Gröbner basis Algorithm". In: Journal of Symbolic Computation (Sept. 2014), pp. 1-24. DOI: 10.1016/j.jsc.2014.09.025 (cit. on pp. 15, 17, 37, 62, 70-72, 78-81).
[Bon+12] B. Bonnard et al. "Geometric optimal control of the contrast imaging problem in nuclear magnetic resonance". In: IEEE Trans. Automat. Control 57.8 (2012), pp. 1957-1969. ISSN: 0018-9286. DOI: 10.1109/TAC.2012.2195859 (cit. on pp. 134, 137).
[Bon +13 ] B. Bonnard, M. Chyba, A. Jacquemard, and J. Marriott. "Algebraic geometric classification of the singular flow in the contrast imaging problem in nuclear magnetic resonance". In: Math. Control Relat. Fields 3.4 (2013), pp. 397-432. ISsN: 2156-8472. DOI: \(10.3934 /\) mcrf.2013.3.397 (cit. on pp. 21, 134, 135, 137, 146).
[Bon+16] B. Bonnard et al. "Determinantal sets, singularities and application to optimal control in medical imagery". In: Proceedings of the 2016 International Symposium on Symbolic and Algebraic Computation. ISSAC '16. Waterloo, Canada, 2016 (cit. on p. 133).
[BP99] M. de Boer and R. Pellikaan. "Gröbner bases for codes". In: Some tapas in computer algebra. Ed. by A. M. Cohen, H. Cuypers, and H. Sterk. Algorithms and Computation in Mathematics 4. Springer, 1999, pp. 237-259. URL: http://www.win. tue.nl/\%5C \(\sim \%{ }_{7}\) B\% \({ }_{7}\) Druudp/paper/34.pdf (cit. on p. 91).
[BS88] D. Bayer and M. Stillman. "On the complexity of computing syzygies". In: fournal of Symbolic Computation 6.2-3 (Oct. 1988), pp. 135-147. ISSN: 07477171. DOI: 10. 1016/S0747-7171(88)80039-7 (cit. on p. 78).
[Buc76] B. Buchberger. "A theoretical basis for the reduction of polynomials to canonical forms". In: ACM SIGSAM Bulletin 10.3 (1976), pp. 19-29 (cit. on pp. 10, 67, 69).
[Buc79] B. Buchberger. "A Criterion for Detecting Unnecessary Reductions in the Construction of Gröbner Bases". In: 72 (Jan. 1979), pp. 3-21. DOI: 10.1007/3-540-09519-5 (cit. on p. 69).
[BV88] W. Bruns and U. Vetter. Determinantal rings. Vol. 1327. Lecture Notes in Mathematics. Springer-Verlag, Berlin, 1988, pp. viii+236. ISBN: 3-540-19468-1 (cit. on pp. 12, 16, 53).
[BW93] T. Becker and V. Weispfenning. Gröbner bases. Vol. 141. Graduate Texts in Mathematics. A computational approach to commutative algebra, In cooperation with Heinz Kredel. New York: Springer-Verlag, 1993, pp. xxii+574. ISBN: 0-387-97971-9. DOI: 10.1007/978-1-4612-0913-3 (cit. on p. 64).
[CDo5] E. Cattani and A. Dickenstein. "Introduction to residues and resultants". In: Solving polynomial equations. Springer Berlin Heidelberg, 2005, pp. 1-61 (cit. on p. 10).
[CDR96] M. Caboara, G. de Dominicis, and L. Robbiano. "Multigraded Hilbert Functions and Buchberger Algorithm". In: Proceedings of the 1996 International Symposium on Symbolic and Algebraic Computation, ISSAC '96, Zurich, Switzerland, Fuly 2426, 1996. Ed. by E. Engeler, B. F. Caviness, and Y. N. Lakshman. ACM, 1996, pp. 7278. ISBN: o-89791-796-O. DOI: 10.1145/236869.236901 (cit. on p. 88).
[Chy+15] M. Chyba et al. "Optimal Geometric Control Applied to the Protein Misfolding Cyclic Amplification Process". English. In: Acta Applicandae Mathematicae 135.1 (2015), pp. 145-173. ISSN: 0167-8019. DOI: 10.1007/S10440-014-9950-8 (cit. on p. 135).
[CKM97] S. Collart, M. Kalkbrenner, and D. Mall. "Converting bases with the Gröbner walk". English. In: Journal of Symbolic Computation. Special issue on computational algebra and number theory: proceedings of the first MAGMA conference 24.3-4 (1997), pp. 465-469. DOI: \(10.1006 /\) jsco.1996.0145 (cit. on p. 14).
[CLOo7] D. Cox, J. Little, and D. O'Shea. Ideals, varieties, and algorithms. Third. Undergraduate Texts in Mathematics. An introduction to computational algebraic geometry and commutative algebra. New York: Springer, 2007, pp. xvi+551. ISBN: 978-0-387-35650-1; 0-387-35650-9. DOI: 10.1007/978-o-387-35651-8 (cit. on pp. 29, 40, 65, 67, 69).
[Col75] G.E. Collins. "Quantifier elimination for real closed fields by cylindrical algebraic decomposition". In: Automata theory and formal languages (Second GI Conf., Kaiserslautern, 1975). Vol. 33. Lecture Notes in Computer Science. Springer, Berlin, 1975, pp. 134-183 (cit. on pp. 10, 135).
[Com74] L. Comtet. Advanced combinatorics. enlarged. The art of finite and infinite expansions. D. Reidel Publishing Co., Dordrecht, 1974, pp. xi+343. ISBN: 90-277-0441-4 (cit. on pp. 31, 32).
[CS95] M. Coste and M. Shiota. "Thom's first isotopy lemma: a semialgebraic version, with uniform bound". In: Real analytic and algebraic geometry (Trento, 1992). de Gruyter, Berlin, 1995, pp. 83-101 (cit. on p. 47).
[Demoo] M. Demazure. Bifurcations and catastrophes. Universitext. Geometry of solutions to nonlinear problems, Translated from the 1989 French original by D. Chillingworth. Springer-Verlag, Berlin, 2000, pp. viii+303. ISBN: 3-540-52118-6. DOI: 10.1007/978-3-642-57134-3 (cit. on p. 46).
[DH88] J. H. Davenport and J. Heintz. "Real quantifier elimination is doubly exponential". In: Journal of Symbolic Computation 5.1-2 (1988), pp. 29-35. ISSN: 0747-7171. DOI: 10.1016/So747-7171(88)80004-X (cit. on p. 135).
[DSo6] G. Dalzotto and E. Sbarra. "Computations in weighted polynomial rings". In: Analele Stiintifice ale Universitatii Ovidius Constanta 14(2) (2006), pp. 31-44. URL: http://www.anstuocmath.ro/mathematics/pdfiz/31_44_GDalzotto_ESbarra.pdf (cit. on p. 89).
[Eis95] D.Eisenbud. Commutative algebra. Vol. 150. Graduate Texts in Mathematics. With a view toward algebraic geometry. New York: Springer-Verlag, 1995, pp. xvi+785. ISBN: 0-387-94268-8; 0-387-94269-6. DOI: 10.1007/978-1-4612-5350-1 (cit. on pp. 27, 29, 32, 33, 45, 50, 51, 53, 89, 104).
[Fau+13] J.-C. Faugère, P. Gaudry, L. Huot, and G. Renault. "Using Symmetries in the Index Calculus for Elliptic Curves Discrete Logarithm". English. In: Journal of Cryptology (2013), pp. 1-41. ISSN: 0933-2790. DOI: 10.1007/soo145-013-9158-5 (cit. on pp. 11, 20, 91, 126).
\([\mathrm{Fau}+14] \quad\) J.-C. Faugère, P. Gaudry, L. Huot, and G. Renault. "Sub-cubic change of ordering for Gröbner basis". In: Proceedings of the 39th International Symposium on Symbolic and Algebraic Computation - ISSAC '14. Kobe, Japan: ACM Press, July 2014, pp. 170-177. ISBN: 9781450325011. DOI: \(10.1145 / 2608628.2608669\) (cit. on pp. 83, 123).
[Fau+93] J.-C. Faugère, P. Gianni, D. Lazard, and T. Mora. "Efficient computation of zerodimensional Gröbner bases by change of ordering". In: Journal of Symbolic Computation 16.4 (1993), pp. 329-344. ISSN: 0747-7171. DOI: 10.1006/jsco.1993.1051 (cit. on pp. 14, 75, 77, 82).
[Fauo2] J.-C. Faugère. "A new efficient algorithm for computing Gröbner bases without reduction to zero ( \(F_{5}\) )". In: Proceedings of the 2002 International Symposium on Symbolic and Algebraic Computation. ACM, 2002, 75-83 (electronic). DOI: 10. 1145/780506.780516 (cit. on pp. 10, 69-71, 136).
[Fau1o] J.-C. Faugère. "FGb: A Library for Computing Gröbner Bases". In: Mathematical Software - ICMS 2010. Ed. by K. Fukuda, J. Hoeven, M. Joswig, and N. Takayama. Vol. 6327. Lecture Notes in Computer Science. Kobe, Japan: Springer Berlin / Heidelberg, Sept. 2010, pp. 84-87. DOI: 10.1007/ 978-3-642-15582-6_17 (cit. on pp. 69, 91, 136).
[Fau99] J.-C. Faugère. "A new efficient algorithm for computing Gröbner bases \(\left(F_{4}\right)\) ". In: Journal of Pure and Applied Algebra 139.1-3 (1999). Effective methods in algebraic geometry (Saint-Malo, 1998), pp. 61-88. ISSN: 0022-4049. DOI: 10.1016/Soo22-4049(99)00005-5 (cit. on pp. 10, 69).
[FLPo8] J.-C. Faugère, F. Levy-dit-Vehel, and L. Perret. "Cryptanalysis of Minrank". In: Advances in Cryptology CRYPTO 2008. Ed. by D. Wagner. Vol. 5157. Lecture Notes in Computer Science. Santa Barbara, CA, USA: Springer-Verlag, Aug. 2008, pp. 280-296. ISBN: 978-3-540-85173-8. DOI: 10.1007/978-3-540-85174-5_16 (cit. on p. 12).
[FM13] J.-C. Faugère and C. Mou. "Sparse FGLM algorithms". Anglais. Preprint available at http://hal.inria.fr/hal-oo807540. 2013 (cit. on p. 82).
[Frö85] R. Fröberg. "An inequality for Hilbert series of graded algebras". eng. In: Mathematica Scandinavica 56 (1985), pp. 117-144. URL: http://eudml.org/doc/166929 (cit. on p. 50).
[FSo9] P. Flajolet and R. Sedgewick. Analytic combinatorics. Cambridge University Press, Cambridge, 2009, pp. xiv+810. ISBN: 978-0-521-89806-5. DOI: 10.1017 / CBO9780511801655. URL: http://dx.doi.org/10.1017/CBO9780511801655 (cit. on pp. 19, 31, 32).
[FS12] J.-C. Faugère and J. Svartz. "Solving Polynomial Systems Globally Invariant Under an Action of the Symmetric Group and Application to the Equilibria of N vortices in the Plane". In: ISSAC '12: Proceedings of the 2012 international symposium on Symbolic and algebraic computation. 2012, pp. 170-178. URL: http: \(/ / w w w-p o l s y s . l i p 6 . f r / \sim j c f /\) Papers/FS12.pdf (cit. on p. 16).
[FSS11] J.-C. Faugère, M. Safey El Din, and P.-J. Spaenlehauer. "Gröbner bases of bihomogeneous ideals generated by polynomials of bidegree ( 1,1 ): algorithms and complexity". In: fournal of Symbolic Computation 46.4 (2011), pp. 406-437. IssN: 0747-7171. DOI: 10.1016/j.jsc.2010.10.014 (cit. on p.16).
[FSS12] J.-C. Faugère, M. Safey El Din, and P.-J. Spaenlehauer. "Critical Points and Gröbner Bases: the Unmixed Case". In: Proceedings of the 2012 International Symposium on Symbolic and Algebraic Computation (ISSAC 2012). 2012, pp. 162-169 (cit. on p. 136).
[FSS13] J.-C. Faugère, M. Safey El Din, and P.-J. Spaenlehauer. "On the Complexity of the Generalized MinRank Problem". In: fournal of Symbolic Computation 55 (2013), pp. 30-58 (cit. on pp. 16, 136).
[FSV13] J.-C. Faugère, M. Safey El Din, and T. Verron. "On the complexity of computing Gröbner bases for quasi-homogeneous systems". In: Proceedings of the 2013 International Symposium on Symbolic and Algebraic Computation. ISSAC '13. Boston, USA: ACM, 2013 (cit. on pp. 87, 89).
[FSV15] J.-C. Faugère, M. Safey El Din, and T. Verron. "On the complexity of computing Gröbner bases for weighted homogeneous systems". In: (2015). To appear in Journal of Symbolic Computations. DOI: 10.1016/j.jsc.2015.12.001 (cit. on pp. 87, 89).
[Gauo9] P. Gaudry. "Index calculus for abelian varieties of small dimension and the elliptic curve discrete logarithm problem". In: Journal of Symbolic Computation 44.12 (2009). Gröbner Bases in Cryptography, Coding Theory, and Algebraic Combinatorics, pp. 1690-1702. ISSN: 0747-7171. DOI: 10.1016/j.jsc.2008.08.005 (cit. on p. 126).
[Gio+91] A. Giovini et al. ""One sugar cube, please" or selection strategies in the Buchberger algorithm". In: Proceedings of the 1991 international symposium on Symbolic and algebraic computation - ISSAC '91. Vol. 91. 1991, pp. 49-54. ISBN: o897914376. DOI: 10.1145/120694.120701 (cit. on p. 69).
[GKZ94] I. M. Gelfand, M. M. Kapranov, and A. V. Zelevinsky. Discriminants, Resultants, and Multidimensional Determinants. Boston, MA: Birkhäuser Boston, 1994. ISBN: 978-o-8176-4770-4, 978-o-8176-4771-1. DOI: 10.1007/978-0-8176-4771-1 (cit. on p. 34).
[GLSo1] M. Giusti, G. Lecerf, and B. Salvy. "A Gröbner Free Alternative for Polynomial System Solving". In: fournal of Complexity 17.1 (2001), pp. 154-211. IssN: 0885064X. DOI: \(10.1006 /\) jcom.2000.0571 (cit. on p. 10).
[GM89] P. Gianni and T. Mora. "Algebraic solution of systems of polynomial equations using Groebner bases". In: Applied algebra, algebraic algorithms and errorcorrecting codes (Menorca, 1987). Vol. 356. Lecture Notes in Comput. Sci. Springer, Berlin, 1989, pp. 247-257. DOI: 10.1007/3-540-51082-6_83. URL: http://dx.doi.org/10. 1007/3-540-51082-6_83 (cit. on p. 66).
[GRo9] E. Guerrini and A. Rimoldi. "FGLM-Like Decoding: from Fitzpatrick's Approach to Recent Developments". English. In: Gröbner Bases, Coding, and Cryptography. Ed. by M. Sala et al. Springer Berlin Heidelberg, 2009, pp. 197-218. ISBN: 978-3-540-93805-7. DOI: 10.1007/978-3-540-93806-4_12 (cit. on p. 91).
[Har77] R. Hartshorne. Algebraic geometry. Graduate Texts in Mathematics, No. 52. New York: Springer-Verlag, 1977, pp. xvi+496. ISBN: 0-387-90244-9 (cit. on pp. 40-42).
[HNS15a] D. Henrion, S. Naldi, and M. Safey El Din. "Exact algorithms for linear matrix inequalities". Submitted. Aug. 2015. URL: https:// hal.archives-ouvertes.fr/hal-o1184320 (cit. on pp. 135, 136).
[HNS15b] D. Henrion, S. Naldi, and M. Safey El Din. "Real root finding for low rank linear matrices". Submitted. June 2015. URL: https://hal.archives-ouvertes.fr/hal-o1159210 (cit. on pp. 53, 135, 136, 143).
[ \(\left.\mathrm{HNS}_{15} \mathrm{c}\right]\) D. Henrion, S. Naldi, and M. Safey El Din. "Real root finding for rank defects in linear Hankel matrices". In: Proceedings of the 2015 ACM on International Symposium on Symbolic and Algebraic Computation. Bath, United Kingdom, July 2015. URL: https:// hal.archives-ouvertes.fr/hal-o1114378 (cit. on pp. 135, 136).
[HNS16] D. Henrion, S. Naldi, and M. Safey El Din. "Real root finding for determinants of linear matrices". In: Journal of Symbolic Computation 74 (2016), pp. 205-238. ISSN: 0747-7171. DOI: http://dx.doi.org/10.1016/j.jsc.2015.06.010 (cit. on pp. 135, 136).
[HS12] H. Hong and M. Safey El Din. "Variant quantifier elimination". In: fournal of Symbolic Computation 47.7 (2012). International Symposium on Symbolic and Algebraic Computation (ISSAC 2009), pp. 883-901. ISSN: 0747-7171. DOI: http: \(/ / d x\).doi.org/10.1016/j.jsc.2011.05.014 (cit. on p. 135).
[Lap+12] M. Lapert et al. "Exploring the Physical Limits of Saturation Contrast in Magnetic Resonance Imaging". In: Scientific Reports 2.589 (2012). DOI: 10.1038/srepoo589. URL: http://dx.doi.org/10.1038/srepoo589 (cit. on p. 20).
[Laz83] D. Lazard. "Gröbner-Bases, Gaussian elimination and resolution of systems of algebraic equations". In: EUROCAL. Ed. by J. A. van Hulzen. Vol. 162. Lecture Notes in Computer Science. Springer, 1983, pp. 146-156. ISBN: 3-540-12868-9. DOI: 10.1007/3-540-12868-9_99 (cit. on pp. 15, 70).
[LeGalli4] F. Le Gall. "Powers of Tensors and Fast Matrix Multiplication". In: CoRR (2014). arXiv: 1401/7714 (cit. on p. 78).
[Leoo9] D. A. Leonard. "A weighted module view of integral closures of affine domains of type I". In: Advances in Mathematics of Communications 3.1 (2009), pp. 1-11. ISSN: 1930-5346. DOI: 10.3934 /amc.2009.3.1 (cit. on p. 91).
[LMXo5] F. Lemaire, M. Moreno Maza, and Y. Xie. "The RegularChains library". In: Maple conference. Vol. 5. 2005, pp. 355-368 (cit. on pp. 10, 136).
[LRo7] D. Lazard and F. Rouillier. "Solving parametric polynomial systems". In: fournal of Symbolic Computation 42.6 (2007), pp. 636-667 (cit. on pp. 21, 135, 139).
[Luc91] É. Lucas. Théorie des nombres. Vol. 1. Théorie des nombres. Gauthier-Villars et fils, 1891 (cit. on p. 97).
[Macaulay2] D. R. Grayson and M. E. Stillman. Macaulay2, a software system for research in algebraic geometry. Available at http://www.math.uiuc.edu/Macaulayz/. 2014 (cit. on pp. 69, 93).
[Magma] W. Bosma, J. Cannon, and C. Playoust. "The Magma algebra system. I. The user language". In: fournal of Symbolic Computation 24•3-4 (1997). Computational algebra and number theory (London, 1993), pp. 235-265. ISSN: 0747-7171. DOI: 10.1006/ jsco.1996.0125 (cit. on pp. 69, 88, 91).
[Maple] Waterloo Maple (Maplesoft). Computer Algebra System. URL: www.maplesoft.com/ products/maple/ (cit. on p. 69).
[McC99] S. McCallum. "On projection in CAD-based quantifier elimination with equational constraint". In: Proceedings of the 1999 International Symposium on Symbolic and Algebraic Computation (Vancouver, BC). ACM, New York, 1999, 145-149 (electronic). DOI: \(10.1145 / 309831.309892\) (cit. on p. 135).
[MM82] E. W. Mayr and A. R. Meyer. "The complexity of the word problems for commutative semigroups and polynomial ideals". In: Advances in Mathematics 46.3 (Dec. 1982), pp. 305-329. ISSN: 00018708. DOI: 10.1016/0001-8708(82)90048-2 (cit. on p. 78).
[Moro3] G. Moreno-Socías. "Degrevlex Gröbner bases of generic complete intersections". In: Journal of Pure and Applied Algebra 180 (2003), pp. 263-283. URL: http://www. sciencedirect.com/science/article/pii/So022404902002979 (cit. on p. 50).
[Mor96] G. Moreno-Socías. Revlex standard bases of generic complete intersections. Technical report. 1996 (cit. on pp. 17, 35, 90, 100-102).
[Nal15] S. Naldi. "Exact algorithms for determinantal varieties and semidefinite programming". PhD thesis. Sept. 2015 (cit. on pp. 16, 46, 47).
[NZM91] I. M. Niven, H. S. Zuckerman, and H. L. Montgomery. An introduction to the theory of numbers. Wiley, 1991. ISBN: 9780471625469 (cit. on p. 97).
[Ott+15] J. C. Ottem, K. Ranestad, B. Sturmfels, and C. Vinzant. "Quartic spectrahedra". en. In: Mathematical Programming 151.2 (July 2015), pp. 585-612. ISSN: 0025-5610, 1436-4646. DOI: \(10.1007 /\) s10107-014-0844-3 (cit. on p. 12).
\begin{tabular}{|c|c|}
\hline [Parıo] & K. Pardue. "Generic sequences of polynomials". In: fournal of Algebra 324.4 (2010), pp. 579-590. ISSN: 0021-8693. DOI: 10.1016/j.jalgebra.2010.04.018 (cit. on pp. 39, 49). \\
\hline [Pon+62] & L. S. Pontryagin, V. G. Boltyanskii, R. V. Gamkrelidze, and E. F. Mishchenko. The mathematical theory of optimal processes. Translated from the Russian by K. N. Trirogoff; edited by L. W. Neustadt. Interscience Publishers John Wiley \& Sons, Inc. New York-London, 1962, pp. viii+36o (cit. on p. 137). \\
\hline [RRR91] & L. Reid, L. G. Roberts, and M. Roitman. "On complete intersections and their Hilbert functions". In: Canadian Mathematical Bulletin \(34 \cdot 4\) (1991), pp. 525-535. ISSN: 0008-4395. DOI: 10.4153/CMB-1991-083-9 (cit. on pp. 116, 117). \\
\hline [ \(\mathrm{RT}_{15}\) ] & J. I. Rodriguez and X. Tang. "Data-Discriminants of Likelihood Equations". In: Proceedings of the 2015 ACM on International Symposium on Symbolic and Algebraic Computation. ISSAC '15. Bath, United Kingdom: ACM, 2015, pp. 307-314 ISBN: 978-1-4503-3435-8. DOI: 10.1145/2755996.2756649 (cit. on p. 135). \\
\hline [SF82] & J. J. Sylvester and F. Franklin. "A Constructive Theory of Partitions, Arranged in Three Acts, an Interact and an Exodion". In: Amer. 7. Math. 5.1-4 (1882), pp. 251330. ISSN: OOO2-9327. DOI: 10.2307/2369545. URL: http://dx.doi.org/10.2307/2369545 (cit. on p. 31). \\
\hline [Singular] & W. Decker, G.-M. Greuel, G. Pfister, and H. Schönemann. Singular 3-1-6 - A computer algebra system for polynomial computations. http://www.singular.unikl.de. 2012 (cit. on pp. 69, 124). \\
\hline [Spa12] & P.-J. Spaenlehauer. "Résolution de systèmes multi-homogènes et déterminantiels" PhD thesis. Oct. 2012. URL: http://www.pjspaenlehauer. net/data/these\% \({ }_{7 B} \%_{5} C_{-}\) \%7Dspaenlehauer.pdf (cit. on pp. 18, 43, 44). \\
\hline [Spa14] & P.-J. Spaenlehauer. "On the complexity of computing critical points with Gröbner bases". In: SIAM fournal on Optimization 24.3 (2014), pp. 1382-1401 (cit. on p. 136) \\
\hline [Str69] & V. Strassen. "Gaussian elimination is not optimal". In: Numerische Mathematik 13.4 (Aug. 1969), pp. 354-356. ISSN: 0029-599X. DOI: 10.1007/BFo2165411 (cit. on p. 78). \\
\hline [Stuo8] & B. Sturmfels. Algorithms in Invariant Theory (Texts and Monographs in Symbolic Computation). 2nd ed.; VII, 197 pp.; 5 figs. Springer Publishing Company, Incorporated, 2008. ISBN: 3211774165,9783211774168 (cit. on pp. 91, 129). \\
\hline [SWo5] & A. J. Sommese and C. W. Wampler. The Numerical Solution of Systems of Polynomials Arising in Engineering and Science. World Scientific Publishing Company, Incorporated, 2005. ISBN: 9789812567727 (cit. on p. 10). \\
\hline [Tra96] & C. Traverso. "Hilbert Functions and the Buchberger Algorithm". In: fournal of Symbolic Computation 22.4 (1996), pp. 355-376. IssN: 0747-7171. DOI: http: //dx.doi.org/10.1006/jsco.1996.0056 (cit. on pp. 20, 69, 88, 129). \\
\hline
\end{tabular}
[YHXor] L. Yang, X. Hou, and B. Xia. "A complete algorithm for automated discovering of a class of inequality-type theorems". In: Sci. China Ser. F 44.1 (2001), pp. 33-49. IsSN: 1009-2757. Doi: \(10.1007 /\) BFoz713938 (cit. on pp. 21, 135, 139).

\section*{Index}

Algorithm \(\mathrm{F}_{5}, 71\)
(weighted case), 118
Complexity, 78, 81
Complexity (weighted case), 89, 118, 120
Algorithm FGLM, 76
Complexity, 82
Complexity (weighted case), 90
Bézout's bound, 44, 90
birational equivalence, 41
Buchberger's algorithm, 69
codimension
of a variety, 42
of an ideal, 33
complete intersection, 34
critical
point, 45
value, 45
degree, 28
of a polynomial, 29
of an ideal, 33
degree fall, 74
degree of regularity, 78
weighted, 93
dehomogenization, 29
denumerant, 31
determinantal
ideal, 52
variety, 52
dimension
of a manifold, 47
of a variety, 42
of an ideal, 33
divisor of zero, 27
elimination order, 63
equidimensional component, 42
equidimensional variety, 42
\(\mathrm{F}_{5}\) criterion, 70
Fröberg's conjecture, 50
Frobenius number, 31
generating series, 32
genericity, 48
graded morphism, 28
graded reverse-lexicographical ordering, 61
graded ring, 28
Gröbner basis, 63
reduced, 65
Hilbert series, 32
of a homogeneous semi-regular sequence, 39
of a polynomial algebra, 32
of a regular sequence, 34
of a weighted-homogeneous semi-regular sequence, 109
Hilbert's Nullstellensatz, 40
homogeneous, 28
component, 28
ideal, 28
polynomial, 29
homogenization, 29
incidence variety, 56
index of regularity, 33
initial ideal, 60
integral, 36
irreducibility, 40
irreducible component, 41

Jacobian criterion, 46
Jacobian matrix, 45
Krull dimension, 42
leading
coefficient, 60
ideal, 60
monomial, 60
term, 60
lexicographical ordering, 60
local dimension
(algebraic geometry), 42
(manifold), 47
locally closed, 41
Macaulay matrix, 33
Macaulay's bound, 79
weighted, 89,105
weighted (simultaneous Noether position), 106
monomial ordering, 59
Elim, 63
GRevLex, 61
Lex, 60
W-GRevLex, 61
morphism of varieties, 41
Noether position, 37
Noether's normalization theorem, 50
weighted, 96
Noetherian ring, 27
normal Form, 64
Nullstellensatz, 40
proper map, 47
rational map, 41
real semi-algebraic set, 46
reduction algorithm, 68
regular point, 45
regular sequence, 34
in the affine sense, 74
regular value, 45
reverse chain-divisible, 95

Sard's theorem, 46
saturation, 67
Schur complement, 55
semi-regular sequence, 39
weighted homogeneous, 109
shape position, 66
signature, 70
simultaneous Noether position, 37
singular point, 45
smooth variety, 45
\(S\)-polynomial, 67
strongly \(W\)-compatible, 98
structure lemma, 8 o
(weighted case), 120
Tarski-Seidenberg theorem, 47
Thom's isotopy lemma, 47
Thom's weak transversality theorem, 46 total degree, 29
variety, 40
W-compatible, 98
strongly, 98
weighted degree, 29, 87
weighted degree of regularity, 93
weighted graded reverse-lexicographical ordering, 61
weighted homogeneous, 29, 87
Zariski topology, 40

\title{
Régularisation du calcul de bases de Gröbner pour des systèmes avec poids et déterminantiels, et application en imagerie médicale
}

\author{
Thibaut Verron
}

\begin{abstract}
Résumé
La résolution de systèmes polynomiaux est un problème aux multiples applications, et les bases de Gröbner sont un outil important dans ce cadre. Il est connu que de nombreux systèmes issus d'applications présentent une structure supplémentaire par rapport à des systèmes arbitraires, et que ces structures peuvent souvent être exploitées pour faciliter le calcul de bases de Gröbner.
Dans cette thèse, on s'intéresse à deux exemples de telles structures, pour différentes applications. Tout d'abord, on étudie les systèmes homogènes avec poids, qui sont homogènes si on calcule le degré en affectant un poids spécifique à chaque variable. Cette structure apparaitt naturellement dans de nombreuses applications, dont un problème de cryptographie (logarithme discret). On montre comment les algorithmes existants, efficaces pour les polynômes homogènes, peuvent être adaptés au cas avec poids. Les bornes de complexité pour des systèmes homogènes avec poids génériques sont alors divisées par un facteur polynomial en le produit des poids, par rapport aux bornes pour des systèmes homogènes génériques de même degré.
Par ailleurs, on étudie un problème de classification de racines réelles pour des variétés définies par des déterminants. Ce problème a une application directe en théorie du contrôle, pour l'optimisation de contraste de l'imagerie à résonance magnétique (IRM). Ce problème particulier s'avère hors de portée des algorithmes préexistants pour la classification. On montre comment, en adaptant ces stratégies pour tirer profit de la structure déterminantielle du système, les calculs deviennent praticables. De plus, on illustre ce procédé en apportant des réponses aux questions posées par le problème d'optimisation de contraste.
\end{abstract}

\begin{abstract}
Polynomial system solving is a problem with numerous applications, and Gröbner bases are an important tool in this context. Previous studies have shown that systèmes arising in applications usually exhibit more structure than arbitrary systems, and that these structures can be used to make Gröbner bases computations easier.
In this thesis, we consider two examples of such structures. First, we study weighted homogeneous systems, which are homogeneous if we give to each variable a specific weight when computing the degree. This structure appears naturally in many applications, including a cryptographical problem (discrete logarithm). We show how existing algorithms, which are efficient for homogeneous systems, can be adapted to a weighted setting. The complexity bounds for generic weighted homogeneous systems are divided by a factor polynomial in the product of the weights, when compared to the bounds for generic homogeneous systems with the same degree.

Then we consider a real roots classification problem for varieties defined by determinants. This problem has a direct application in control theory, for contrast optimization in magnetic resonance imagery (MRI). This specific problem appears to be out of reach of previous algorithms. We show how, by adapting these algorithms to the determinantal structure of the system, the computations become practicable. As an illustration, we answer the questions from the application to contrast optimization.
\end{abstract}```


[^0]:    ${ }^{1}$ In [Baro4, Chap. 4], the remainder $O\left(n^{1 / 4}\right)$ was written as $o(\sqrt{n})$. However, it appears that in the proof, this $o(\sqrt{n})$ is a rewriting of $\sqrt{n} \cdot O(\sqrt{\Delta z})$, where $\Delta z=O(1 / \sqrt{n})$.

[^1]:    ${ }^{2}$ The Gröbner basis algorithms from Magma seem to behave strangely with elimination orders, as seen in the detailed logs, and it coincides with significant slowdowns. This behavior was not observed on other implementations of the same algorithms: $\mathrm{F}_{4}$ from FGb and Buchberger from Singular [Singular]. For example, for the system in the first line of table 3.2c, without the weights, with Magma's $F_{4}$ algorithm, the first degree fall comes at step 4, and the algorithm needs more than 66 steps to compute the basis. With FGb's implementation of $F_{4}$ in Maple, the first degree fall appears at step 13, and the algorithm finishes at step 32.

[^2]:    ${ }^{1}$ The full results, together with the source code which produced them are available at: mercurey.gforge.inria.fr

