Solving Polynomial Equations and Applications

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Abstract

These notes accompany an introductory lecture given by the author at the workshop on solving polynomial equations & applications at CWI Amsterdam in the context of the 2022 fall semester programme on polynomial optimization & applications. We introduce systems of polynomial equations and the main approaches for solving them. We also discuss applications and solution counts. The theory is illustrated by many examples.

1 Systems of polynomial equations

Let K be a field with algebraic closure \overline{K} . The polynomial ring with n variables and coefficients in K is $R = K[x_1, \ldots, x_n]$. We abbreviate $x = (x_1, \ldots, x_n)$ and use variable names x, y, z rather than x_1, x_2, x_3 when n is small. Elements of R are polynomials, which we think of as functions $f: \overline{K}^n \to \overline{K}$ of the form

$$f(x) = \sum_{\alpha \in \mathbb{N}^n} c_{(\alpha_1, \dots, \alpha_n)} x_1^{\alpha_1} \cdots x_n^{\alpha_n} = \sum_{\alpha \in \mathbb{N}^n} c_{\alpha} x^{\alpha},$$

with finitely many nonzero coefficients $c_{\alpha} \in K$. A system of polynomial equations is

$$f_1(x) = \dots = f_s(x) = 0, \tag{1}$$

where $f_1, \ldots, f_s \in R$. By a solution of (1) we mean a point $x \in \overline{K}^n$ satisfying all of these s equations. Solving usually means finding coordinates for all solutions. This makes sense only when the set of solutions is finite, which typically happens when $s \geq n$. However, systems with infinitely many solutions can be 'solved' too, in an appropriate sense [40]. We point out that one is often mostly interested in solutions $x \in K^n$ over the ground field K. The reason for allowing solutions over the algebraic closure \overline{K} is that many solution methods, like those discussed in Section 4, intrinsically compute all such solutions. Here are some examples.

Example 1.1 (Univariate polynomials: n = 1). When n = s = 1, solving the polynomial system defined by $f = c_0 + c_1 x + \cdots + c_d x^d \in K[x]$, with $c_d \neq 0$, amounts to finding the roots of f(x) = 0 in \overline{K} . These are the eigenvalues of the $d \times d$ companion matrix

$$C_f = \begin{pmatrix} & & -c_0/c_d \\ 1 & & -c_1/c_d \\ & \ddots & & \vdots \\ & & 1 & -c_{d-1}/c_d \end{pmatrix}$$
 (2)

of f, whose characteristic polynomial is $\det(x \cdot \mathrm{id} - C_f) = c_d^{-1} \cdot f$.

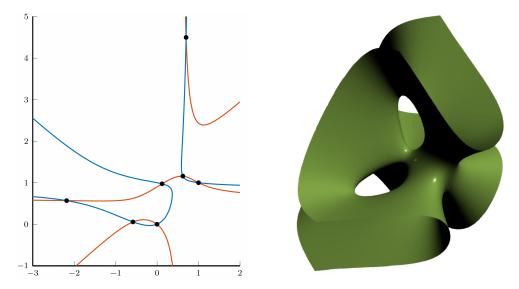


Figure 1: Algebraic curves in the plane (n = 2) and an algebraic surface (n = 3).

Example 1.2 (Linear equations). When $f_i = \sum_{j=1}^n a_{ij} x_j - b_i$ are given by affine-linear functions, (1) is a linear system of the form Ax = b, with $A \in K^{s \times n}$, $b \in K^s$.

Examples 1.1 and 1.2 show that, after a trivial rewriting step, the univariate and affinelinear cases are reduced to a *linear algebra* problem. Here, we are mainly interested in the case where n > 1, and some equations are of higher degree in the variables x. Such systems require tools from *nonlinear algebra* [34]. We proceed with an example in two dimensions.

Example 1.3 (Intersecting two curves in the plane). Let $K = \mathbb{Q}$ and n = s = 2. We work in the ring $R = \mathbb{Q}[x, y]$ and consider the system of equations f(x, y) = g(x, y) = 0 where

$$f = -7x - 9y - 10x^{2} + 17xy + 10y^{2} + 16x^{2}y - 17xy^{2},$$

$$g = 2x - 5y + 5x^{2} + 5xy + 5y^{2} - 6x^{2}y - 6xy^{2}.$$

Geometrically, we can think of f(x,y)=0 as defining a curve in the plane. This is the red curve shown in Figure 1. The curve defined by g(x,y)=0 is shown in blue. The set of solutions of f=g=0 consists of points $(x,y)\in\overline{\mathbb{Q}}^2$ satisfying f(x,y)=g(x,y)=0. These are the intersection points of the two curves. There are seven such points in $\overline{\mathbb{Q}}^2$, of which two lie in \mathbb{Q}^2 . These are the points (0,0) and (1,1). Note that all seven solutions are real. replacing \mathbb{Q} by $K=\mathbb{R}$, we count as many solutions over K as over $\overline{K}=\mathbb{C}$.

The set of solutions of the polynomial system (1) is called an affine variety. We denote this by $V_{\overline{K}}(f_1,\ldots,f_s)=\{x\in\overline{K}^n\mid f_1(x)=\cdots=f_s(x)=0\}$, and replace \overline{K} by K in this notation to mean only the solutions over the ground field. Examples of affine varieties are the red curve $V_{\overline{K}}(f)$ and the set of black dots $V_{\overline{K}}(f,g)$ in Example 1.3. In the case of $V_{\overline{K}}(f)$, Figure (1) only shows the real part $V_{\overline{K}}(f)\cap\mathbb{R}^2$ of the affine variety.

Example 1.4 (Surfaces in \mathbb{R}^3). Let $K = \mathbb{R}$ and consider the affine variety $V = V_{\mathbb{C}}(f)$ where

$$f = 81(x^3 + y^3 + z^3) - 189(x^2y + x^2z + y^2x + y^2z + xz^2 + yz^2) + 54xyz + 126(xy + xz + yz) - 9(x^2 + y^2 + z^2) - 9(x + y + z) + 1.$$
(3)

Its real part $V_{\mathbb{R}}(f)$ is the surface shown in the right part of Figure 1. The variety V is called the Clebsch surface. It is a cubic surface because it is defined by an equation of degree three. Note that f is invariant under permutations of the variables, i.e. f(x,y,z) = f(y,x,z) = f(z,y,x) = f(x,z,y) = f(y,z,x). This reflects in symmetries of the surface $V_{\mathbb{R}}(f)$. Many polynomials from applications have similar symmetry properties. Exploiting this in computations is an active area of research, see for instance [28].

More pictures of real affine varieties can be found, for instance, in [16, Chapter 1, §2], or in the algebraic surfaces gallery hosted at https://homepage.univie.ac.at/herwig.hauser/bildergalerie/gallery.html.

We now briefly discuss commonly used fields K. In many engineering applications, the coefficients of f_1, \ldots, f_s live in \mathbb{R} or \mathbb{C} . Computations in such fields use floating point arithmetic, yielding approximate results. The required quality of the approximation depends on the application. Other fields show up too: polynomial systems in cryptography often use $K = \mathbb{F}_q$, see for instance [39]. Equations of many prominent algebraic varieties have integer coefficients, i.e. $K = \mathbb{Q}$. Examples are determinantal varieties (e.g. the variety of all 2×2 matrices of rank ≤ 1), Grassmannians in their Plücker embedding [34, Chapter 5], discriminants and resultants [44, Sections 3.4, 5.2] and toric varieties obtained from monomial maps [45, Section 2.3]. In number theory, one is interested in studying rational points $V_{\mathbb{Q}}(f_1, \ldots, f_s) \subset V_{\mathbb{Q}}(f_1, \ldots, f_s)$ on varieties defined over \mathbb{Q} . Recent work in this direction for del Pezzo surfaces can be found in [35, 17]. Finally, in tropical geometry, coefficients come from valued fields such as the p-adic numbers \mathbb{Q}_p or Puiseux series $\mathbb{C}\{\{t\}\}$ [32]. Solving over the field of Puiseux series is also relevant for homotopy continuation methods, see [44, Section 6.2.1]. We end the section with two examples highlighting the difference between $V_K(f_1, \ldots, f_s)$ and $V_{\overline{K}}(f_1, \ldots, f_s)$.

Example 1.5 (Fermat's last theorem). Let $k \in \mathbb{N} \setminus \{0\}$ be a positive integer and consider the equation $f = x^k + y^k - 1 = 0$. For any k, the variety $V_{\overline{\mathbb{Q}}}(f)$ has infinitely many solutions in $\overline{\mathbb{Q}}^2$. For k = 1, 2, there are infinitely many solutions in \mathbb{Q}^2 . For $k \geq 3$, the only solutions in \mathbb{Q}^2 are (1,0),(0,1) and, when k is even, (-1,0),(0,-1).

Example 1.6 (Computing real solutions). The variety $V_{\mathbb{C}}(x^2+y^2)$ consists of the two lines $x+\sqrt{-1}\cdot y=0$ and $x-\sqrt{-1}\cdot y=0$ in \mathbb{C}^2 . However, the real part $V_{\mathbb{R}}(x^2+y^2)=\{(0,0)\}$ has only one point. If we are interested only in this real solution, we may replace x^2+y^2 by the two polynomials x,y, which have the property that $V_{\mathbb{R}}(x^2+y^2)=V_{\mathbb{R}}(x,y)=V_{\mathbb{C}}(x,y)$. An algorithm that computes all *complex* solutions will still recover only the interesting solutions, after this replacing step. It turns out that such a 'better' set of equations can always be computed. The new polynomials generate the *real radical ideal* associated to the original equations [34, Section 6.3]. For recent computational progress, see for instance [1].

2 Applications

Polynomial equations appear in many fields of science and engineering. Some examples are molecular biology [22], computer vision [29], economics and game theory [42, Chapter 6], topological data analysis [9] and partial differential equations [42, Chapter 10]. For an overview and more references, see [14, 10]. In this section, we present a selection of other applications in some detail. Much of the material is taken from [44, Section 1.2].

2.1 Polynomial optimization

In the spirit of the semester programme on polynomial optimization & applications, let us consider the problem of minimizing a polynomial objective function $g(x_1, \ldots, x_k) \in \mathbb{R}[x_1, \ldots, x_k]$ over the real affine variety $V_{\mathbb{R}}(h_1, \ldots, h_\ell) \subset \mathbb{R}^k$, with $h_1, \ldots, h_\ell \in \mathbb{R}[x_1, \ldots, x_k]$. That is, we consider the polynomial optimization problem [31]

$$\min_{x \in \mathbb{R}^k} g(x_1, \dots, x_k),
\text{subject to} h_1(x_1, \dots, x_k) = \dots = h_{\ell}(x_1, \dots, x_k) = 0.$$
(4)

Introducing new variables $\lambda_1, \ldots, \lambda_\ell$ we obtain the Lagrangian $L = g - \lambda_1 h_1 - \cdots - \lambda_\ell h_\ell$, whose partial derivatives give the optimality conditions

$$\frac{\partial L}{\partial x_1} = \dots = \frac{\partial L}{\partial x_k} = h_1 = \dots = h_\ell = 0.$$
 (5)

This is a polynomial system with $n = s = k + \ell$, and our field is $K = \mathbb{R}$. Only real solutions are candidate minimizers. The methods in Section 4 compute all complex solutions and select the real ones among them. The number of solutions over \mathbb{C} is typically finite.

Example 2.1 (Euclidean distance degree). Given a general point $y = (y_1, ..., y_k) \in \mathbb{R}^k$, we consider the (squared) Euclidean distance function $g(x_1, ..., x_k) = ||x - y||_2^2 = (x_1 - y_1)^2 + ... + (x_k - y_k)^2$. Let Y be the real affine variety defined by $h_1, ..., h_\ell \in \mathbb{R}[x_1, ..., x_k]$:

$$Y = \{x \in \mathbb{R}^k \mid h_1 = \dots = h_\ell = 0\}.$$

Consider the optimization problem (4) given by these data. The solution x^* is the point on Y that's closest to y. The number of *complex* solutions of (5) is called the *Euclidean distance degree* of Y [19]. For a summary and examples, see [43, Section 2].

Example 2.2 (Parameter estimation for system identification). System identification is an engineering discipline that aims at constructing models for dynamical systems from measured data. A model explains the relation between input, output, and noise. It depends on a set of model parameters, which are selected to best fit the measured data. A discrete time, single-input single-output linear time invariant system with input sequence $u: \mathbb{Z} \to \mathbb{R}$, output sequence $y: \mathbb{Z} \to \mathbb{R}$ and white noise sequence $e: \mathbb{Z} \to \mathbb{R}$ is often modeled by

$$A(q) y(t) = \frac{B_1(q)}{B_2(q)} u(t) + \frac{C_1(q)}{C_2(q)} e(t).$$

Here $A, B_1, B_2, C_1, C_2 \in \mathbb{C}[q]$ are unknown polynomials of a fixed degree in the *backward* shift operator q, acting on $s : \mathbb{Z} \to \mathbb{R}$ by qs(t) = s(t-1). The model parameters are the coefficients of these polynomials, which are to be estimated. Clearing denominators gives

$$A(q)B_2(q)C_2(q)y(t) = B_1(q)C_2(q)u(t) + B_2(q)C_1(q)e(t).$$
(6)

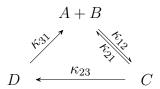
Suppose we have measured $u(0), \ldots, u(N), y(0), \ldots, y(N)$. Then we can find algebraic relations among the coefficients of A, B_1, B_2, C_1, C_2 by writing (6) down for $t = d, d + 1, \ldots, N$ where $d = \max(d_A + d_{B_2} + d_{C_2}, d_{B_1} + d_{C_2}, d_{B_2} + d_{C_1})$ and $d_A, d_{B_1}, d_{B_2}, d_{C_1}, d_{C_2}$ are the degrees of our polynomials. The model parameters are estimated by solving

$$\min_{\Theta \in \mathbb{R}^k} \quad e(0)^2 + \ldots + e(N)^2 \quad \text{subject to} \quad (6) \text{ is satisfied for } t = d, \ldots, N$$

where Θ consists of $e(0), \ldots, e(N)$ and the coefficients of A, B_1, B_2, C_1, C_2 . We refer to [3, Section 1.1.1] for a worked out example and more references.

2.2 Chemical reaction networks

The equilibrium concentrations of the chemical species occuring in a *chemical reaction network* satisfy algebraic relations. Taking advantage of the algebraic structure of these networks has led to advances in the understanding of their dynamical behaviour. We refer the interested reader to [18, 13] and references therein. The network below involves 4 species A, B, C, D and models T cell signal transduction (see [18]).



The parameters κ_{12} , κ_{21} , κ_{31} , $\kappa_{23} \in \mathbb{R}_{>0}$ are the reaction rate constants. Let x_A, x_B, x_C, x_D denote the time dependent concentrations of the species A, B, C, D respectively. The law of mass action gives the relations

$$f_A = \frac{dx_A}{dt} = -\kappa_{12}x_Ax_B + \kappa_{21}x_C + \kappa_{31}x_D, \quad f_B = \frac{dx_B}{dt} = -\kappa_{12}x_Ax_B + \kappa_{21}x_C + \kappa_{31}x_D,$$

$$f_C = \frac{dx_C}{dt} = \kappa_{12}x_Ax_B - \kappa_{21}x_C - \kappa_{23}x_C, \qquad f_D = \frac{dx_D}{dt} = \kappa_{23}x_C - \kappa_{31}x_D.$$

The set $\{(x_A, x_B, x_C, x_D) \in (\mathbb{R}_{>0})^4 \mid f_A = f_B = f_C = f_D = 0\}$ is called the *steady state* variety of the chemical reaction network. By the structure of the equations, for given initial concentrations, the solution (x_A, x_B, x_C, x_D) cannot leave its *stoichiometric compatibility* class, which is an affine subspace of $(\mathbb{R}_{>0})^4$. Adding the affine equations of the stoichiometric compatibility class to the system, we get the set of all candidate steady states. In this application, we are interested in *positive* solutions, rather than *real* solutions. It is important to characterize parameter values for which there is a unique positive solution, see [37].

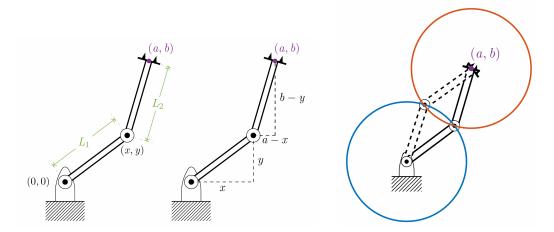


Figure 2: The two configurations of a robot arm are the intersection points of two circles.

2.3 Robotics

We present a simple example of how polynomial equations arise in robotics. Consider a planar robot arm whose shoulder is fixed at the origin (0,0) in the plane, and whose two arm segments have fixed lengths L_1 and L_2 . Our aim is to determine the possible positions of the elbow (x,y), given that the hand of the robot touches a given point (a,b). The situation is illustrated in Figure 2. The Pythagorean theorem gives the identities

$$x^{2} + y^{2} - L_{1}^{2} = (a - x)^{2} + (b - y)^{2} - L_{2}^{2} = 0,$$
(7)

which is a system of s = 2 equations in n = 2 variables x, y. The plane curves corresponding to these equations (see Example 1.3) are shown in blue and orange in Figure 2. Their intersection points, i.e. the solutions to the equations, correspond to the possible configurations. One easily imagines more complicated robots leading to more involved equations, see [51].

3 Number of solutions

It is well known that a univariate polynomial $f \in \mathbb{C}[x]$ of degree d has at most d roots in \mathbb{C} . Moreover, d is the *expected* number of roots, in the following sense. Consider the *family*

$$\mathcal{F}(d) = \{ a_0 + a_1 x + \dots + a_d x^d \mid (a_0, \dots, a_d) \in \mathbb{C}^{d+1} \} \simeq \mathbb{C}^{d+1}$$
 (8)

of polynomials of degree at most d. There is an affine variety $\nabla_d \subset \mathbb{C}^{d+1}$, such that all $f \in \mathcal{F}(d) \setminus \nabla_d$ have precisely d roots in \mathbb{C} . Here $\nabla_d = V_{\mathbb{C}}(\Delta_d)$, where Δ_d is a polynomial in the coefficients a_i of $f \in \mathcal{F}(d)$. Equations for small d are

$$\begin{split} &\Delta_1 = a_1, \\ &\Delta_2 = a_2 \cdot (a_1^2 - 4 \, a_0 a_2), \\ &\Delta_3 = a_3 \cdot (a_1^2 a_2^2 - 4 \, a_0 a_2^3 - 4 \, a_1^3 a_3 + 18 \, a_0 a_1 a_2 a_3 - 27 \, a_0^2 a_3^2), \\ &\Delta_4 = a_4 \cdot (a_1^2 a_2^2 a_3^2 - 4 \, a_0 a_2^3 a_3^2 - 4 \, a_1^3 a_3^3 + 18 \, a_0 a_1 a_2 a_3^3 + \dots + 256 \, a_0^3 a_4^3). \end{split}$$

Notice that $\Delta_d = a_d \cdot \widetilde{\Delta}_d$, where $\widetilde{\Delta}_d$ is the <u>discriminant</u> for degree d polynomials. There exist similar results for families of polynomial systems with n > 1, which bound the number of isolated solutions from above by the <u>expected</u> number. This section states some of these results. The field $K = \overline{K}$ is algebraically closed throughout the section.

3.1 Bézout's theorem

Let $R = K[x] = K[x_1, \ldots, x_n]$. A monomial in R is a finite product of variables: $x^{\alpha} = x^{\alpha_1} \cdots x^{\alpha_n}$, $\alpha \in \mathbb{N}$. The degree of the monomial x^{α} is $\deg(x^{\alpha}) = \sum_{i=1}^{n} \alpha_i$, and the degree of a polynomial $f = \sum_{\alpha} c_{\alpha} x^{\alpha}$ is $\deg(f) = \max_{\{\alpha: c_{\alpha} \neq 0\}} \deg(x^{\alpha})$. We define the vector subspaces

$$R_d = \{ f \in R : \deg(f) \le d \}, \quad d \ge 0.$$

For an *n*-tuple of degrees (d_1, \ldots, d_n) , we define the family of polynomial systems

$$\mathcal{F}(d_1,\ldots,d_n) = R_{d_1} \times \cdots \times R_{d_n}.$$

That is, $F = (f_1, \ldots, f_n) \in \mathcal{F}(d_1, \ldots, d_n)$ satisfies $\deg(f_i) \leq d_i, i = 1, \ldots, n$, and represents the polynomial system F = 0 with s = n. We leave the fact that $\mathcal{F}(d_1, \ldots, d_n) \simeq K^D$, with $D = \sum_{i=1}^n \binom{n+d_i}{n}$, as an exercise to the reader. Note that this is a natural generalization of (8). The set of solutions of F = 0 is denoted by $V_K(F) = V_K(f_1, \ldots, f_n)$, and a point in $V_K(F)$ is isolated if it does not lie on a component of $V_K(F)$ with dimension ≥ 1 .

Theorem 3.1 (Bézout). For any $F = (f_1, \ldots, f_n) \in \mathcal{F}(d_1, \ldots, d_n)$, the number of isolated solutions of $f_1 = \cdots = f_n = 0$, i.e., the number of isolated points in $V_K(F)$, is at most $d_1 \cdots d_n$. Moreover, there exists a proper subvariety $\nabla_{d_1, \ldots, d_n} \subsetneq K^D$ such that, when $F \in \mathcal{F}(d_1, \ldots, d_n) \setminus \nabla_{d_1, \ldots, d_n}$, the variety $V_K(F)$ consists of precisely $d_1 \cdots d_n$ isolated points.

The proof of this theorem can be found in [21, Theorem III-71]. As in our univariate example, the variety ∇_{d_1,\dots,d_n} can be described using discriminants and resultants. See for instance the discussion at the end of [44, Section 3.4.1]. Theorem 3.1 is an important result and gives an easy way to bound the number of isolated solutions of a system of n equations in n variables. The bound is almost always tight, in the sense that the only systems with fewer solutions lie in ∇_{d_1,\dots,d_n} . Unfortunately, many systems coming from applications lie inside ∇_{d_1,\dots,d_n} . For instance, the system (7) with $K = \mathbb{C}$ lies in $\nabla_{2,2}$, by the fact that the two real solutions seen in 2 are the only complex solutions, and $2 < d_1 \cdot d_2 = 4$. One also checks that for f, g as in Example 1.3, we have $(f, g) \in \nabla_{3,3} \subset \mathcal{F}(3,3)$ (use $K = \overline{\mathbb{Q}}$).

Remark 3.2. Bézout's theorem more naturally counts solutions in projective space \mathbb{P}_K^n , and it accounts for solutions with multiplicity > 1. More precisely, if $f_i \in H^0(\mathbb{P}_K^n, \mathcal{O}_{\mathbb{P}_K^n}(d_i))$ is a homogeneous polynomial in n+1 variables of degree d_i , and $f_1 = \cdots = f_n = 0$ has finitely many solutions in \mathbb{P}_K^n , the number of solutions (counted with multiplicity) is always $d_1 \cdots d_n$. We encourage the reader who is familiar with projective geometry to check that (7) defines two solutions at infinity, when each of the equations is viewed as a global section of $\mathcal{O}_{\mathbb{P}_2^n}(2)$.

3.2 Kushnirenko's theorem

An intuitive consequence of Theorem 3.1 is that random polynomial systems given by polynomials of fixed degree always have the same number of solutions. Looking at f and g from Example 1.3, we see that they do not look so random, in the sense that some monomials of degree ≤ 3 are missing. For instance, x^3 and y^3 do not appear. Having zero coefficients standing with some monomials in $\mathcal{F}(d_1,\ldots,d_n)$ is sometimes enough to conclude that the system lies in ∇_{d_1,\ldots,d_n} . That is, the system is not random in the sense of Bézout's theorem.

The *(monomial)* support supp(f) of $f = \sum_{\alpha} c_{\alpha} x^{\alpha}$ is the set of exponents appearing in f:

$$\operatorname{supp}\left(\sum_{\alpha} c_{\alpha} x^{\alpha}\right) = \left\{\alpha : c_{\alpha} \neq 0\right\} \subset \mathbb{N}^{n}.$$

This subsection considers families of polynomial systems whose equations have a fixed support. Let $\mathcal{A} \subset \mathbb{N}^n$ be a finite subset of exponents in \mathbb{N}^n of cardinality $|\mathcal{A}|$. We define

$$\mathcal{F}(\mathcal{A}) = \{ (f_1, \dots, f_n) \in \mathbb{R}^n : \operatorname{supp}(f_i) \subset \mathcal{A}, i = 1, \dots, n \} \simeq K^{n \cdot |\mathcal{A}|}.$$

Kushnirenko's theorem counts the number of solutions for systems in the family $\mathcal{F}(\mathcal{A})$. It expresses this in terms of the volume $\operatorname{Vol}(\mathcal{A}) = \int_{\operatorname{Conv}(\mathcal{A})} d\alpha_1 \cdots d\alpha_n$ of the convex polytope

$$\operatorname{Conv}(\mathcal{A}) = \left\{ \sum_{\alpha \in \mathcal{A}} \lambda_{\alpha} \cdot \alpha : \lambda_{\alpha} \ge 0, \sum_{\alpha \in \mathcal{A}} \lambda_{\alpha} = 1 \right\} \subset \mathbb{R}^{n}.$$
 (9)

The normalized volume vol(A) is defined as $n! \cdot Vol(A)$.

Theorem 3.3 (Kushnirenko). For any $F = (f_1, \ldots, f_n) \in \mathcal{F}(\mathcal{A})$, the number of isolated solutions of $f_1 = \cdots = f_n = 0$ in $(K \setminus \{0\})^n$, i.e., the number of isolated points in $V_K(F) \cap (K \setminus \{0\})^n$, is at most $\operatorname{vol}(\mathcal{A})$. Moreover, there exists a proper subvariety $\nabla_{\mathcal{A}} \subsetneq K^{n \cdot |\mathcal{A}|}$ such that, when $F \in \mathcal{F}(\mathcal{A}) \setminus \nabla_{\mathcal{A}}$, $V_K(F) \cap (K \setminus \{0\})^n$ consists of precisely $\operatorname{vol}(\mathcal{A})$ isolated points.

For a proof, see [30]. The theorem necessarily counts solutions in $(K \setminus \{0\})^n \subset K^n$, as multiplying all equations with a monomial x^{α} may change the number of solutions in the coordinate hyperplanes (i.e., there may be new solutions with zero-coordinates), but it does not change the normalized volume $\operatorname{vol}(\mathcal{A})$. The statement can be adapted to count solutions in K^n , but becomes more involved [27]. We point out that, with the extra assumption that $0 \in \mathcal{A}$, one may replace $(K \setminus \{0\})^n$ by K^n in Theorem 3.3.

Example 3.4 (Kushnirenko VS Bézout). If $\mathcal{A} = \{\alpha \in \mathbb{N}^n : \deg(x^{\alpha}) \leq d\}$, we have $\mathcal{F}(\mathcal{A}) = \mathcal{F}(d, \ldots, d)$ and $d^n = \operatorname{vol}(\mathcal{A})$. Theorem 3.3 recovers Theorem 3.1 for $d_1 = \cdots = d_n$.

Example 3.5 (Example 1.3 continued). The polynomial system f = g = 0 from Example 1.3 belongs to the family $\mathcal{F}(\mathcal{A})$ with $\mathcal{A} = \{(1,0),(0,1),(2,0),(1,1),(0,2),(2,1),(1,2)\}$. The convex hull $\operatorname{Conv}(\mathcal{A})$ is a hexagon in \mathbb{R}^2 , see Figure 3. Its normalized volume is $\operatorname{vol}(\mathcal{A}) = n! \cdot \operatorname{Vol}(\mathcal{A}) = 2! \cdot 3 = 6$. Theorem 3.3 predicts 6 solutions in $(\overline{\mathbb{Q}} \setminus \{0\})^2$. These are six out of the seven black dots seen in the left part of Figure 1: the solution (0,0) is not counted. We have a chain of inclusions $\nabla_{\mathcal{A}} \subset \mathcal{F}(\mathcal{A}) \subset \nabla_{3,3} \subset \mathcal{F}(3,3)$ and $(f,g) \in \mathcal{F}(\mathcal{A}) \setminus \nabla_{\mathcal{A}}$.

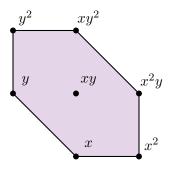


Figure 3: The polytope Conv(A) from Example 3.5 is a hexagon.

Remark 3.6. The analogue of Remark 3.2 for Theorem 3.3 is that vol(A) counts solutions on the *projective toric variety* X_A associated to A. It equals the degree of X_A in its embedding in $\mathbb{P}_K^{|A|-1}$ (after multiplying with a lattice index). When A is as in Example 3.4, we have $X_A = \mathbb{P}^n$. A toric proof of Kushnirenko's theorem and examples are given in [45, Section 3.4].

Remark 3.7. The convex polytope Conv(supp(f)) is called the *Newton polytope* of f. Its importance goes beyond counting solutions: it is dual to the *tropical hypersurface* defined by f, which is a *combinatorial shadow* of $V_K(f) \cap (K \setminus \{0\})^n$ [32, Proposition 3.1.6].

3.3 Bernstein's theorem

There is a generalization of Kushnirenko's theorem which allows different supports for the polynomials f_1, \ldots, f_n . We fix n finite subsets of exponents $\mathcal{A}_1, \ldots, \mathcal{A}_n$ with respective cardinalities $|\mathcal{A}_i|$. These define the family of polynomial systems

$$\mathcal{F}(\mathcal{A}_1,\ldots,\mathcal{A}_n) = \{ (f_1,\ldots,f_n) \in \mathbb{R}^n : \operatorname{supp}(f_i) \subset \mathcal{A}_i, i = 1,\ldots,n \} \simeq \mathbb{K}^D,$$

where $D = |\mathcal{A}_1| + \cdots + |\mathcal{A}_n|$. The number of solutions is characterized by the mixed volume of $\mathcal{A}_1, \ldots, \mathcal{A}_n$, which we now define. The Minkowski sum S + T of two sets $S, T \subset \mathbb{R}^n$ is $\{s + t : s \in S, t \in T\}$, where s + t is the usual addition of vectors in \mathbb{R}^n . For a nonnegative real number λ , the λ -dilation of $S \subset \mathbb{R}^n$ is $\lambda \cdot S = \{\lambda \cdot s : s \in S\}$, where $\lambda \cdot s$ is the usual scalar multiplication in \mathbb{R}^n . Each of the supports \mathcal{A}_i gives a convex polytope $\operatorname{Conv}(\mathcal{A}_i)$ as in (9). The function $\mathbb{R}^n_{\geq 0} \to \mathbb{R}_{\geq 0}$ given by

$$(\lambda_1, \ldots, \lambda_n) \longmapsto \operatorname{Vol}(\lambda_1 \cdot \operatorname{Conv}(\mathcal{A}_1) + \cdots + \lambda_n \cdot \operatorname{Conv}(\mathcal{A}_n))$$

is a homogeneous polynomial of degree n, meaning that all its monomials have degree n [15, Chapter 7, §4, Proposition 4.9]. The *mixed volume* $MV(A_1, \ldots, A_n)$ is the coefficient of that polynomial standing with $\lambda_1 \cdots \lambda_n$. Note that $MV(A, \ldots, A) = vol(A)$.

Theorem 3.8 (Bernstein-Kushnirenko). For any $F = (f_1, \ldots, f_n) \in \mathcal{F}(\mathcal{A}_1, \ldots, \mathcal{A}_n)$, the number of isolated solutions of $f_1 = \cdots = f_n = 0$ in $(K \setminus \{0\})^n$, i.e., the number of isolated points in $V_K(F) \cap (K \setminus \{0\})^n$, is at most $MV(\mathcal{A}_1, \ldots, \mathcal{A}_n)$. Moreover, there exists a proper

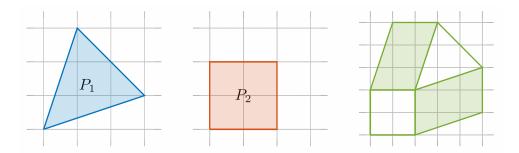


Figure 4: The green area counts the solutions to equations with support in P_1, P_2 .

subvariety $\nabla_{\mathcal{A}_1,...,\mathcal{A}_n} \subset K^D$ such that, when $F \in \mathcal{F}(\mathcal{A}_1,...,\mathcal{A}_n) \setminus \nabla_{\mathcal{A}_1,...,\mathcal{A}_n}$, $V_K(F) \cap (K \setminus \{0\})^n$ consists of precisely $MV(\mathcal{A}_1,...,\mathcal{A}_n)$ isolated points.

This theorem was originally proved by Bernstein for $K = \mathbb{C}$ in [6]. The proof by Kushnirenko in [30] works for algebraically closed fields. Several alternative proofs were found by Khovanskii. Theorem 3.8 is sometimes called the *BKK theorem*, after the aforementioned mathematicians. Like Kushnirenko's theorem, Theorem 3.8 can be adapted to count roots in K^n rather than $(K \setminus \{0\})^n$ [27], and if $0 \in \mathcal{A}_i$ for all i, one may replace $(K \setminus \{0\})^n$ by K^n .

Example 3.9. When $\mathcal{A}_1 = \cdots = \mathcal{A}_n = \mathcal{A}$, we have $\mathcal{F}(\mathcal{A}_1, \ldots, \mathcal{A}_n) = \mathcal{F}(\mathcal{A})$, and when $\mathcal{A}_i = \{\alpha \in \mathbb{N}^n : \deg(x^{\alpha}) \leq d_i\}$, we have $\mathcal{F}(\mathcal{A}_1, \ldots, \mathcal{A}_n) = \mathcal{F}(d_1, \ldots, d_n)$. Hence, all families we have seen before are of this form, and Theorem 3.8 generalizes Theorems 3.1 and 3.3. \diamond

Example 3.10. A useful formula for n = 2 is $MV(A_1, A_2) = Vol(A_1 + A_2) - Vol(A_1) - Vol(A_2)$. For instance, the following two polynomials appear in [44, Example 5.3.1]:

$$f = a_0 + a_1 x^3 y + a_2 x y^3$$
, $g = b_0 + b_1 x^2 + b_2 y^2 + b_3 x^2 y^2$.

The system f = g = 0 is a general member of the family $\mathcal{F}(\mathcal{A}_1, \mathcal{A}_2) \simeq K^7$, where $\mathcal{A}_1 = \{(0,0),(3,1),(1,3)\}$ and $\mathcal{A}_2 = \{(0,0),(2,0),(0,2),(2,2)\}$. The Newton polygons, together with their Minkowski sum, are shown in Figure 4. By applying the formula for $MV(\mathcal{A}_1, \mathcal{A}_2)$ seen above, we find that the mixed volume for the system f = g = 0 is the area of the green regions in the right part of Figure 4, which is 12. Note that the Bézout bound (Theorem 3.1) is 16. and Theorem 3.3 also predicts 12 solutions, with $\mathcal{A} = \mathcal{A}_1 \cup \mathcal{A}_2$. Hence $\mathcal{F}(\mathcal{A}_1, \mathcal{A}_2) \subset \mathcal{F}(\mathcal{A}) \subset \nabla_{4,4} \subset \mathcal{F}(4,4)$ and $\mathcal{F}(\mathcal{A}_1, \mathcal{A}_2) \not\subset \nabla_{\mathcal{A}}$.

Theorem 3.8 provides an upper bound on the number of isolated solutions to any system of polynomial equations with n = s. Although it improves significantly on Bézout's bound for many systems, it still often happens that the bound is not tight for systems in applications. That is, one often encounters systems $F \in \nabla_{\mathcal{A}_1,\dots,\mathcal{A}_n}$. Even more refined root counts exist, such as those based on *Khovanskii bases*. In practice, with today's computational methods (see Section 4), we often count solutions reliably by simply solving the system. *Certification methods* provide a proof for a *lower* bound on the number of solutions [11]. The actual number of solutions is implied if one can match this with a theoretical upper bound.

4 Computational methods

We give a brief introduction to two of the most important computational methods for solving polynomial equations. The first method uses normal forms, the second is based on homotopy continuation. We keep writing $F = (f_1, \ldots, f_s) = 0$ for the system we want to solve. We require $s \geq n$, and assume finitely many solutions over \overline{K} . All methods discussed here compute all solutions over \overline{K} , so we will assume that $K = \overline{K}$ is algebraically closed. An important distinction between normal forms and homotopy continuation is that the former works over any field K, while the latter needs $K = \mathbb{C}$. If the coefficients are contained in a subfield (e.g. $\mathbb{R} \subset \mathbb{C}$), a significant part of the computation in normal form algorithms can be done over this subfield. Also, homotopy continuation is most natural when n = s, whereas s > n is not so much a problem for normal forms. However, if $K = \mathbb{C}$ and n = s, continuation methods are extremely efficient and can compute millions of solutions.

4.1 Normal form methods

Let $I = \langle f_1, \ldots, f_s \rangle \subset R = K[x_1, \ldots, x_n]$ be the ideal generated by our polynomials. For ease of exposition, we assume that I is radical, which is equivalent to all points in $V_K(I) = V_K(f_1, \ldots, f_s)$ having multiplicity one. In other words, the Jacobian matrix $(\partial f_i/\partial x_j)$, evaluated at any of the points in $V_K(I)$, has rank n. Let us write $V_K(I) = \{z_1, \ldots, z_\delta\} \subset K^n$ for the set of solutions, and R/I for the quotient ring obtained from R by the equivalence relation $f \sim g \Leftrightarrow f - g \in I$. The main observation behind normal form methods is that the coordinates of z_i are encoded in the eigenstructure of the K-linear endomorphisms $M_g: R/I \to R/I$ given by $[f] \mapsto [g \cdot f]$, where [f] is the residue class of f in R/I.

We will now make this precise. First, we show that $\dim_K R/I = \delta$. We define $\operatorname{ev}_i : R/I \to K$ as $\operatorname{ev}_i([f]) = f(z_i)$, and combine these to get

$$\operatorname{ev} = (\operatorname{ev}_1, \dots, \operatorname{ev}_{\delta}) : R/I \longrightarrow K^{\delta}, \text{ given by } \operatorname{ev}([f]) = (f(z_1), \dots, f(z_{\delta})).$$

By Hilbert's Nullstellensatz [16, Chapter 4], a polynomial $f \in R$ belongs to I if and only if $f(z_i) = 0, i = 1, ..., \delta$. In other words, the map ev is injective. It is also surjective: there exist Lagrange polynomials $\ell_i \in R$ satisfying $\ell_i(z_j) = 1$ if i = j and $\ell_i(z_j) = 0$ for $i \neq j$ [44, Lemma 3.1.2]. We conclude that ev establishes the K-vector space isomorphism $R/I \simeq K^{\delta}$.

The following statement makes our claim that the zeros z_1, \ldots, z_δ are encoded in the eigenstructure of M_q concrete.

Theorem 4.1. The left eigenvectors of the K-linear map M_g are the evaluation functionals $\operatorname{ev}_i, i = 1, \ldots, \delta$. The eigenvalue corresponding to ev_i is $g(z_i)$.

Proof. We have $(ev_i \circ M_g)([f]) = ev_i([g \cdot f]) = g(z_i)f(z_i) = g(z_i) \cdot ev_i([f])$, which shows that ev_i is a left eigenvector with eigenvalue $g(z_i)$. Moreover, ev_1, \ldots, ev_δ form a complete set of eigenvectors, since $ev : R/I \to K^\delta$ is a K-vector space isomorphism.

We encourage the reader to check that the residue classes of the Lagrange polynomials $[\ell_i] \in R/I$ form a complete set of right eigenvectors. We point out that, after choosing a basis of R/I, the functional ev_i is represented by a row vector w_i^{T} of length δ , and M_q is a

 $\delta \times \delta$ matrix. Such matrices are called *multiplication matrices*. The eigenvalue relation in the proof of Theorem 4.1 reads more familiarly as $w_i^{\top} M_g = g(z_i) \cdot w_i^{\top}$. Theorem 4.1 suggests to break up the task of computing $V_K(I) = \{z_1, \dots, z_{\delta}\}$ into two parts:

- (A) Compute multiplication matrices M_g and
- (B) extract the coordinates of z_i from their eigenvectors or eigenvalues.

For step (B), let $\{[b_1], \ldots, [b_{\delta}]\}$ be a K-basis for R/I, with $b_j \in R$. The vector w_i is explicitly given by $w_i = (b_1(z_i), \ldots, b_{\delta}(z_i))$. If the coordinate functions x_1, \ldots, x_n are among the b_j , one reads the coordinates of z_i directly from the entries of w_i . If not, some more processing might be needed. Alternatively, one can choose $g = x_j$ and read the j-th coordinates of the z_i from the eigenvalues of M_{x_j} . There are many things to say about these procedures, in particular about their efficiency and numerical stability. We refer the reader to [44, Remark 4.3.4] for references and more details, and do not elaborate on this here.

We turn to step (A), which is where *normal forms* come into play. Suppose a basis $\{[b_1], \ldots, [b_{\delta}]\}$ of R/I is fixed. We identify R/I with the K-vector space $B = \operatorname{span}_K(b_1, \ldots, b_{\delta}) \subset R$. For any $f \in R$, there are unique constants $c_j(f) \in K$ such that

$$f - \sum_{j=1}^{\delta} c_j(f) \cdot b_j \in I. \tag{10}$$

These are the coefficients in the unique expansion of $[f] = \sum_{j=1}^{\delta} c_j(f) \cdot [b_j]$ in our basis. The K-linear map $\mathcal{N}: R \to B$ which sends f to $\sum_{j=1}^{\delta} c_j(f) \cdot b_j$ is called a normal form. Its key property is that \mathcal{N} projects R onto B along I, meaning that $\mathcal{N} \circ \mathcal{N} = \mathcal{N}$ ($\mathcal{N}_{|B}$ is the identity), and $\ker \mathcal{N} = I$. The multiplication map $M_g: B \to B$ is simply given by $M_g(b) = \mathcal{N}(g \cdot b)$. More concretely, the i-th column of the matrix representation of M_g contains the coefficients $c_j(g \cdot b_i), j = 1, \ldots, \delta$ of $\mathcal{N}(g \cdot b_i)$. Here is a familiar example.

Example 4.2. Let $I = \langle f \rangle = \langle c_0 + c_1 x + \dots + c_d x^d \rangle$ be the ideal generated by the univariate polynomial $f \in K[x]$ in Example 1.1. For general c_i , there are $\delta = d$ roots with multiplicity one, hence I is radical. The dimension $\dim_K K[x]/I$ equals d, and a canonical choice of basis is $\{[1], [x], \dots, [x^{d-1}]\}$. Let us construct the matrix M_x in this basis. That is, we set g = x. We compute the normal forms $\mathcal{N}(x \cdot x^{i-1}), i = 1, \dots, d$:

$$\mathcal{N}(x^i) = x^i, i = 1, \dots, d-1 \text{ and } \mathcal{N}(x^d) = -c_d^{-1}(c_0 + c_1x + \dots + c_{d-1}x^{d-1}).$$

One checks this by verifying that $x^i - \mathcal{N}(x^i) \in \langle f \rangle$. The coefficient vectors $c_j(x^i), j = 1, \ldots, d$ of the $\mathcal{N}(x^i)$ are precisely the columns of the companion matrix C_f . We conclude that $M_x = C_f$ and Theorem 4.1 confirms that the eigenvalues of C_f are the roots of f.

Computing normal forms can be done using linear algebra on certain structured matrices, called *Macaulay matrices*. We illustrate this with an example from [47].

Example 4.3. Consider the ideal $I = \langle f, g \rangle \subset \mathbb{Q}[x, y]$ given by $f = x^2 + y^2 - 2$, $g = 3x^2 - y^2 - 2$. The variety $V_{\overline{\mathbb{Q}}}(I) = V_{\mathbb{Q}}(I)$ consists of four points $\{(-1, -1), (-1, 1), (1, -1), (1, 1)\}$,

as predicted by Theorem 3.1. We construct a <u>Macaulay matrix</u> whose rows are indexed by f, xf, yf, g, xg, yg, and whose columns are indexed by all monomials of degree at most 3:

$$\mathcal{M} = \begin{bmatrix} x^3 & x^2y & xy^2 & y^3 & x^2 & y^2 & 1 & x & y & xy \\ f & & & & 1 & 1 & & -2 & & \\ xf & 1 & & 1 & & & & -2 & & \\ yf & & 1 & & 1 & & & & -2 & & \\ g & & & 3 & -1 & & & -2 & & \\ yg & & 3 & & -1 & & & & -2 & & \end{bmatrix}.$$

The first row reads $f = 1 \cdot x^2 + 1 \cdot y^2 - 2 \cdot 1$. A basis for $\mathbb{Q}[x,y]/I$ is $\{[1],[x],[y],[xy]\}$. These monomials index the last four columns. We now invert the leftmost 6×6 block and apply this inverse from the left to M:

$$\tilde{\mathcal{M}} = \begin{bmatrix} x^3 & x^2y & xy^2 & y^3 & x^2 & y^2 & 1 & x & y & xy \\ x^3 - x & 1 & & & & & & & & & -1 \\ x^2y - y & 1 & & & & & & -1 \\ & 1 & & & & & & -1 \\ & & 1 & & & & & -1 \\ & & & 1 & & & & -1 \\ & & & & 1 & & & -1 \\ & & & & 1 & & -1 \\ & & & & 1 & & -1 \end{bmatrix}.$$

The rows of $\tilde{\mathcal{M}}$ are linear combinations of the rows of \mathcal{M} , so they represent polynomials in I. The first row reads $x^3 - 1 \cdot x \in I$. Comparing this with (10), we see that we have found that the normal form of x^3 is x. Using the information in $\tilde{\mathcal{M}}$ we can construct M_x and M_y :

$$M_{x} = \begin{bmatrix} [x] & [x^{2}] & [xy] & [x^{2}y] \\ 0 & 1 & 0 & 0 \\ 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 \\ [xy] & [xy] & 0 & 0 & 1 \end{bmatrix}, \quad M_{y} = \begin{bmatrix} [y] & [xy] & [y^{2}] & [xy^{2}] \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \end{bmatrix}.$$

Remark 4.4. The entries of a Macaulay matrix \mathcal{M} are the coefficients of the polynomials f_1, \ldots, f_s . An immediate consequence of the fact that normal forms are computed using linear algebra on Macaulay matrices, is that when the coefficients of f_i are contained in a subfield $\tilde{K} \subset K$, all computations in step (A) can be done over \tilde{K} . That is assuming the polynomials g for which we want to compute M_g have coefficients in \tilde{K} .

As illustrated in Example 4.3, to compute the matrices M_g it is sufficient to determine the restriction of the normal form $\mathcal{N}: R \to B$ to a finite dimensional K-vector space $V \subset R$, containing $g \cdot B$. The restriction $\mathcal{N}_{|V}: V \to B$ is called a truncated normal form, see [46] and [44, Chapter 4]. The dimension of the space V counts the number of columns of the Macaulay matrix.

Usually, one chooses the basis elements b_j of B to be monomials, and g to be a coordinate function x_i . The basis elements may arise, for instance, as standard monomials from a $Gr\ddot{o}bner\ basis$ computation. We briefly discuss this important concept.

4.1.1 Gröbner bases.

Gröbner bases are a powerful tool for symbolic computations in algebraic geometry. A nice way to motivate their definition is by thinking of *Euclidean division* as a candidate for a normal form map. In the case n = 1 from Example 4.2, this rewrites $g \in K[x]$ as

$$g = q \cdot f + r$$
, where $\deg(r) < d$. (11)

Clearly [g] = [r] in $K[x]/\langle f \rangle$, and $r \in B = \operatorname{span}_K(1, x, \dots, x^{d-1})$ implies that $\mathcal{N}(g) = r$.

To generalize this to n > 1 variables, we fix a monomial order \leq on $R = K[x_1, \ldots, x_n]$ and write LT(f) for the leading term of f with respect to \leq . The reader who is not familiar with monomial orders can consult [16, Chapter 2, §2]. As above, let $I = \langle f_1, \ldots, f_s \rangle \subset R$ be a radical ideal such that $|V_K(I)| = \delta < \infty$. As basis elements b_1, \ldots, b_δ of $B \simeq R/I$, we use the $\delta \leq$ -smallest monomials which are linearly independent modulo our ideal I. They are also called standard monomials. By [16, Chapter 9, §3, Theorem 3], there exists an algorithm which, for any input $g \in R$, computes $q_1, \ldots, q_s, r \in R$ such that

$$g = q_1 \cdot f_1 + \dots + q_s \cdot f_s + r$$
, where LT (f_i) does not divide any term of r , for all i . (12)

This algorithm is called multivariate Euclidean division. Note how the condition "LT(f_i) does not divide any term of r, for all i" generalizes $\deg(r) < d$ in (11). From (12), it is clear that [g] = [r]. However, we do not have $r \in B$ in general. Hence, unfortunately, sending g to its remainder r is usually not a normal form... but it is when f_1, \ldots, f_s is a Gröbner basis!

A set of polynomials $g_1, \ldots, g_k \in I$ forms a $Gr\"{o}bner\ basis$ of the ideal I if the leading terms $LT(g_1), \ldots, LT(g_k)$ generate the leading term ideal $\langle LT(g) : g \in I \rangle$. We point out that no finiteness of $V_K(I)$ or radicality of I is required for this definition. The remainder r in $g = q_1 \cdot g_1 + \cdots + q_k \cdot g_k + r$ where $LT(f_i)$ does not divide any term of r, for all i, now satisfies [g] = [r] and $r \in B$. This justifies the following claim:

taking remainder upon Euclidean division by a Gröbner basis is a normal form.

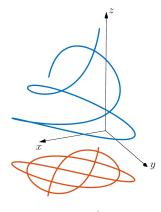
Computing a Gröbner basis g_1, \ldots, g_k from a set of input polynomials f_1, \ldots, f_s can be interpreted as Gaussian elimination on a Macaulay matrix [23]. Once this has been done, multiplication matrices are computed via taking remainder upon Euclidean division by $\{g_1, \ldots, g_k\}$.

On a sidenote, we point out that Gröbner bases are often used for elimination of variables. For instance, if g_1, \ldots, g_k form a Gröbner basis of an ideal I with respect to a lex monomial order for which $x_1 \prec x_2 \prec \cdots \prec x_n$, we have for $j = 1, \ldots n$ that the j-th elimination ideal

$$I_j = I \cap K[x_1, \dots, x_j] = \langle g_i : g_i \in K[x_1, \dots, x_j] \rangle$$

is generated by those elements of our Gröbner basis which involve only the first j variables, see [16, Chapter 3, §1, Theorem 2]. In our case, a consequence is that one of the g_i is univariate in x_1 , and its roots are the x_1 -coordinates of z_1, \ldots, z_{δ} . The geometric counterpart of computing the j-th elimination ideal is projection onto a j-dimensional coordinate space: the variety $V_K(I_j) \subset K^j$ is obtained from $V_K(I) \subset K^n$ by forgetting the final n-j coordinates $(x_1, \ldots, x_n) \mapsto (x_1, \ldots, x_j)$ and taking the closure of the image. Here are two examples.

Example 4.5. To the right we show a blue curve in \mathbb{R}^3 defined by an ideal $I \subset \mathbb{R}[x, y, z]$. Its Gröbner basis with respect to the lex ordering $x \prec y \prec z$ contains $g_1 \in \mathbb{R}[x, y]$, which generates I_2 . The variety $V_{\mathbb{R}}(I_2) = V_{\mathbb{R}}(g_1) \subset \mathbb{R}^2$ is the orange curve seen in the picture. \diamond



Example 4.6. In [35], the authors study del Pezzo surfaces of degree 4 in \mathbb{P}^4 with defining equations $x_0x_1 - x_2x_3 = a_0x_0^2 + a_1x_1^2 + a_2x_2^2 + a_3x_3^2 + a_4x_4^2 = 0$. We will substitute $x_4 = 1 - x_0 - x_1 - x_2 - x_3$ to reduce to the affine case. It is claimed that the smooth del Pezzo surfaces of this form are those for which the parameters a_0, \ldots, a_4 lie outside the hypersurface $H = \{a_0a_1a_2a_3a_4(a_0a_1 - a_2a_3) = 0\}$. This hypersurface is the projection of the variety

$$\left\{ (a,x) \in \mathbb{Q}^5 \times \mathbb{Q}^4 : x_0 x_1 - x_2 x_3 = \sum_{i=0}^3 a_i x_i^2 + a_4 \left(1 - \sum_{i=0}^3 x_i\right)^2 = 0 \text{ and } \operatorname{rank}(J) < 2 \right\}$$

onto \mathbb{Q}^5 . Here J is the 2×4 Jacobian matrix of our two equations with respect to the 4 variables x_0, x_1, x_2, x_3 . The defining equation of H is computed in Macaulay2 [24] as follows:

```
 \begin{array}{l} R = QQ[x_{-0}..x_{-3},a_{-0}..a_{-4}] \\ x_{-4} = 1-x_{-0}-x_{-1}-x_{-2}-x_{-3} \\ I = ideal(\ x_{-0}+x_{-1}-x_{-2}+x_{-3}), \ a_{-0}+x_{-0}+a_{-1}+x_{-1}+a_{-2}+x_{-2}+a_{-3}+x_{-3}+a_{-4}+x_{-4}+a_{-2}) \\ M = submatrix(\ transpose\ jacobian\ I\ ,\ 0...3\ ) \\ radical\ eliminate(\ I+minors(2,M)\ ,\ \{x_{-0},x_{-1},x_{-2},x_{-3}\}\ ) \\ \end{array}
```

The work behind the final command is a Gröbner basis computation.

Remark 4.7. In a numerical setting, it is better to use border bases or more general bases to avoid the amplification of rounding errors. Border bases use basis elements b_i for B whose elements satisfy a connectedness property, see for instance [36] for details. They do not depend on a monomial order. For a summary and comparison between Gröbner bases and border bases, see [44, Sections 3.3.1, 3.3.2]. Recently, in truncated normal form algorithms, bases are selected adaptively by numerical linear algebra routines, such as QR decomposition with optimal column pivoting or singular value decomposition. This often yields a significant improvement in terms of accuracy. See for instance Section 7.2 in [47].

4.2 Homotopy Continuation

The goal of this subsection is to briefly introduce the method of homotopy continuation for solving polynomial systems. For more details, we refer to the standard textbook [41], and to the lecture notes and exercises from the 2021 Workshop on Software and

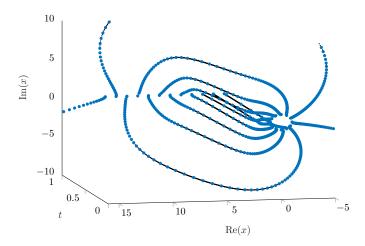


Figure 5: The twelfth roots of unity travel to $1, \ldots, 12$ along continuous paths.

Applications of Numerical Nonlinear Algebra, available at https://github.com/PBrdng/Workshop-on-Software-and-Applications-of-Numerical-Nonlinear-Algebra.

We set $K = \mathbb{C}$ and n = s. We think of $F = (f_1, \ldots, f_n) \in R$ as an element of a family \mathcal{F} of polynomial systems. The reader can replace \mathcal{F} by any of the families seen in Section 3. A homotopy in \mathcal{F} with target system $F \in \mathcal{F}$ and start system $G \in \mathcal{F}$ is a continuous deformation of the map $G = (g_1, \ldots, g_n) : \mathbb{C}^n \to \mathbb{C}^n$ into F, in such a way that all systems obtained throughout the deformation are contained in \mathcal{F} . For instance, When $F \in \mathcal{F}(d_1, \ldots, d_n)$ as in Section 3.1 and G is any other system in $\mathcal{F}(d_1, \ldots, d_n)$, a homotopy is $H(x;t) = t \cdot F(x) + (1-t) \cdot G(x)$, where t runs from 0 to 1. Indeed, for any fixed $t^* \in [0,1]$ the degrees of the equations remain bounded by (d_1, \ldots, d_n) , hence $H(x;t^*) \in \mathcal{F}(d_1, \ldots, d_n)$.

The method of homotopy continuation for solving the target system $f_1 = \cdots = f_n = 0$ assumes that a start system $g_1 = \cdots = g_n = 0$ can easily be solved. The idea is that transforming G continuously into F via a homotopy H(x;t) in \mathcal{F} transforms the solutions of G continuously into those of F. The following example appears in [48].

Example 4.8 $(n = s = 1, \mathcal{F} = \mathcal{F}(12))$. Let $f = (x - 1)(x - 2) \cdots (x - 12)$ be the Wilkinson polynomial of degree 12. We view f = 0 as a member of $\mathcal{F}(12)$ and choose the start system $g = x^{12} - 1 = 0$. The start solutions, i.e. the solutions of g = 0, are the twelfth roots of unity. The solutions of $H(x;t) = t \cdot f(x) + (1-t) \cdot g(x)$ travel from these roots of unity to the integers $1, \ldots, 12$ as t moves from 0 to 1. This is illustrated in Figure 5.

More formally, if $H(x;t) = (h_1(x;t), \ldots, h_n(x;t))$ is a homotopy with $t \in [0,1]$, the solutions describe continuous paths x(t) satisfying H(x(t);t) = 0. Taking the derivative with respect to t gives the *Davidenko differential equation*

$$\frac{\mathrm{d}H(x(t),t)}{\mathrm{d}t} = J_x \cdot \dot{x}(t) + \frac{\partial H}{\partial t}(x(t),t) = 0, \quad \text{with } J_x = \left(\frac{\partial h_j}{\partial x_i}\right)_{i,i}.$$
 (13)

Each start solution x^* of $g_1 = \cdots = g_n = 0$ gives an initial value problem with $x(0) = x^*$, and the corresponding solution path x(t) can be approximated using any numerical ODE

method. This typically leads to a discretization of the solution path, see the blue dots in Figure 5. The solutions of the target system $f_1 = \cdots = f_n = 0$ are obtained by evaluating the solution paths at t = 1. The following are important practical remarks.

Avoid the discriminant. Each of the families seen in Section 3 has a subvariety $\nabla_{\mathcal{F}} \subset \mathcal{F}$ consisting of systems with non-generic behavior in terms of their number of solutions. This subvariety is sometimes referred to as the discriminant of \mathcal{F} , for reasons alluded to in the beginning of Section 3. When the homotopy H(x;t) crosses $\nabla_{\mathcal{F}}$, i.e. $H(x;t^*) \in \nabla_{\mathcal{F}}$ for some $t^* \in [0,1)$, two or more solution paths collide at t^* , or some solution paths diverge. This is not allowed for the numerical solution of (13). Fortunately, crossing $\nabla_{\mathcal{F}}$ can be avoided. The discriminant $\nabla_{\mathcal{F}}$ has complex codimension at least one, hence real codimension at least two. Since the homotopy $t \mapsto H(x;t)$ describes a one-real-dimensional path in \mathcal{F} , it is always possible to go around the discriminant. See for instance [41, Section 7]. When the target system \mathcal{F} belongs to the discriminant, end games are used to deal with colliding/diverging paths at t = 1 [41, Section 10]. Note that this discussion implies that the number of paths tracked in a homotopy algorithm equals the generic number of solutions of the family \mathcal{F} . In that sense, results like Theorems 3.1, 3.3 and 3.8 characterize the complexity of homotopy continuation in the respective families.

Predict and correct. Naively applying ODE methods for solving the Davidenko equation (13) is not the best we can do. Indeed, we have the extra information that the solution paths x(t) satisfy the implicit equation H(x(t),t)=0. This is typically used to improve the accuracy of the ODE solver in each step. Given an approximation of $x(t^*)$ at any fixed $t^* \in [0,1)$ and a step size $0 < \Delta t \ll 1$, one approximates $x(t^* + \Delta t)$ by \tilde{x} using, for instance, a step of Euler's method. This is called the *predictor* step. Then, one refines \tilde{x} to a satisfactory approximation of $x(t^* + \Delta t)$ by using \tilde{x} as a starting point for Newton iteration on $H(x, t^* + \Delta t) = 0$. This is the *corrector step*. The two-step process is illustrated in Figure 6 (predict in orange, correct in green, solution paths x(t) in blue). In the right part of the figure, the predict-correct procedure goes wrong: because of a too large step size Δt in the predictor step, the Newton correction converges to a different path. This phenomenon is called path jumping, and to avoid it one must choose the stepsize Δt carefully and adaptively. Recent work in this direction uses Padé approximants, see [48, 49].

Start systems in practice. There are recipes for start systems in the families \mathcal{F} from Section 3. For instance, we use $G = (x_1^{d_1} - 1, \dots, x_n^{d_n} - 1)$ for $\mathcal{F}(d_1, \dots, d_n)$. The $d_1 \cdots d_n$ solutions can easily be written down. Note that $G \notin \nabla_{d_1,\dots,d_n}$ by Theorem 3.1. For the families $\mathcal{F}(\mathcal{A})$ and $\mathcal{F}(\mathcal{A}_1,\dots,\mathcal{A}_n)$, an algorithm to solve start systems was developed in [26]. For solving start systems of other families, one may use *monodromy loops* [20].

5 Case study: 27 lines on the Clebsch surface

A classical result from intersection theory states that every smooth cubic surface in complex three-space contains exactly 27 lines. In this final section, we use Gröbner bases and

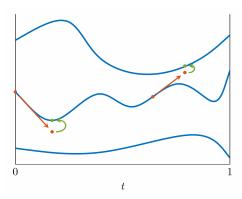


Figure 6: Predictor and corrector steps in numerical continuation. Beware of path jumping.

homotopy continuation to compute lines on the Clebsch surface from Example 1.4. This particular surface is famous for the fact that all its 27 lines are real. Let f(x, y, z) be as in (3). A line in \mathbb{R}^3 parameterized by $(a_1 + t \cdot b_1, a_2 + t \cdot b_2, a_3 + t \cdot b_3)$ is contained in our Clebsch surface if and only if $f(a_1 + t \cdot b_1, a_2 + t \cdot b_2, a_3 + t \cdot b_3) \equiv 0$. The left hand side evaluates to a cubic polynomial in t with coefficients in the ring $\mathbb{Z}[a_1, a_2, a_3, b_1, b_2, b_3] = \mathbb{Z}[a, b]$:

$$f(a_1 + t \cdot b_1, a_2 + t \cdot b_2, a_3 + t \cdot b_3) = f_1(a, b) \cdot t^3 + f_2(a, b) \cdot t^2 + f_3(a, b) \cdot t + f_4(a, b).$$

The lines contained in the Clebsch surface satisfy

$$f_1(a,b) = f_2(a,b) = f_3(a,b) = f_4(a,b) = 0.$$
 (14)

We further reduce this to a system of s = 4 equations in n = 4 unknowns by removing the redundancy in our parameterization of the line: we may impose a random affine-linear relation among the a_i and among the b_i . We choose to substitute

$$a_3 = -(7 + a_1 + 3a_2)/5$$
, $b_3 = -(11 + 3b_1 + 5b_2)/7$.

Implementations of Gröbner bases are available, for instance, in Maple [33] and in msolve [7], which can be used in julia via the package msolve.jl. This is also available in the package Oscar.jl [38]. The following snippet of Maple code constructs our system (14) and computes a Gröbner basis with respect to the graded lexicographic monomial ordering with $a_1 \succ a_2 \succ b_1 \succ b_2$. This basis consists of 23 polynomials g_1, \ldots, g_{23} .

The set of standard monomials is the first output of the command NormalSet. It consists of 27 elements, and the multiplication matrix with respect to a_1 in this basis is constructed using MultiplicationMatrix:

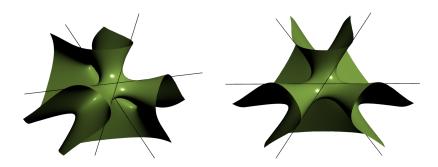


Figure 7: Two views on the Clebsch surface with three of its 27 lines.

```
R,(a1,a2,b1,b2) = PolynomialRing(QQ,["a1","a2","b1","b2"])
I = ideal(R, [-189*b2*b1^2 - 189*b2^2*b1 + 27*(11 + 3*b1 + 5*b2)*b1^2 + ...
A, B = msolve(I)
```

The output B contains 4 rational coordinates (a_1, a_2, b_1, b_2) of 27 lines which approximate the solutions. To see them in floating point format, use for instance

```
[convert.(Float64,convert.(Rational{BigInt},b)) for b in B]
```

We have drawn three of these lines on the Clebsch surface in Figure 7 as an illustration. Other software systems supporting Gröbner bases are Macaulay2 [24], Magma [8] and Singular [25].

Homotopy continuation methods provide an alternative way to compute our 27 lines. Here we use the julia package HomotopyContinuation.jl [12].

The output is shown in Figure 8. There are 27 solutions, as expected. The last line indicates that a :polyhedral start system was used. In our terminology, this means that the system was solved using a homotopy in the family $\mathcal{F}(\mathcal{A}_1,\ldots,\mathcal{A}_4)$ from Section 3.3. The number of tracked paths is 45, which is the mixed volume $MV(\mathcal{A}_1,\ldots,\mathcal{A}_4)$ of this family. The discrepancy 27 < 45 means that our system F lies in the discriminant $\nabla_{\mathcal{A}_1,\ldots,\mathcal{A}_4}$. The 18 'missing' solutions are explained in [4, Section 3.3]. The output also tells us that all solutions have multiplicity one (this is the meaning of non-singular) and all of them are real. Other software implementing homotopy continuation are Bertini [2] and PHCpack [50]. Numerical normal form methods are used in EigenvalueSolver.jl [5].

Result with 27 solutions

- 45 paths tracked
- 27 non-singular solutions (27 real)
- random seed: 0xb8b6077c
- start system: :polyhedral

Figure 8: The julia output when computing 27 lines on the Clebsch surface.

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