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# Formal methods for systems of partial differential equations 

Daniel Robertz

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## 1. Introduction

Given a system of differential equations, we would like to be able to solve the following tasks:
(a) determine all analytic solutions;
(b) obtain an overview of all consequences of the system; in particular, given another differential equation, decide whether it is a consequence of the system or not;
(c) among the consequences find the ones which involve only certain specified unknowns.

Throughout these notes we shall consider partial differential equations (PDEs) for unknown functions $u_{1}\left(z_{1}, \ldots, z_{n}\right), \ldots, u_{m}\left(z_{1}, \ldots, z_{n}\right)$. Since we are going to employ formal methods, we restrict our attention to formal power series solutions in (a). Convergence of these power series on certain regions of $\mathbb{R}^{n}$ or $\mathbb{C}^{n}$ is to be investigated after the formal treatment. In fact, the formal treatment may reveal conditions on how the region in $\mathbb{R}^{n}$ or $\mathbb{C}^{n}$ should be chosen. Singular points will be excluded from consideration.

One of the first existence theorems for a large class of PDEs is the Cauchy-Kovalevskaya Theorem (cf., e.g., [Kov75], [RR04], [Eva10]).

[^0]Theorem 1.1 (Cauchy-Kovalevskaya, 1875). The Cauchy problem
where $a_{i, j, k}$ and $b_{i}$ are real analytic functions around the origin of $\mathbb{R}^{m+n-1}$, has a unique real analytic solution $\left(u_{1}, \ldots, u_{m}\right)$ in a neighborhood of $\left(z_{1}, \ldots, z_{n}\right)=(0, \ldots, 0)$.

Note that any system of differential equations can be rewritten as a system of first order differential equations by introducing new unknown functions, if necessary. The differential equations in Theorem 1.1 are quasilinear in the sense that each equation is linear in the highest derivatives of the unknown functions. Analytic coordinate changes may be used to transform boundary data on an analytic hypersurface which is non-characteristic for the first order PDE system to the hypersurface $z_{1}=0$. Theorem 1.1 is also valid for complex analytic functions. However, the assumption of analyticity is necessary (cf. [Lew57]).

In work of C. Méray [Mér80] and C. Riquier [Riq10] in the second half of the 19th century a generalization of the Cauchy-Kovalevskaya Theorem was obtained. Riquier's Existence Theorem asserts the existence of analytic solutions to systems of PDEs of a certain class (cf. also [Tho28, Tho34], [Rit34, Chap. IX], [Rit50, Chap. VIII]). The equations are assumed to be solved for certain distinct partial derivatives and their right hand sides are analytic functions of $z_{1}, \ldots, z_{n}$ and of partial derivatives of $u_{1}, \ldots, u_{m}$ which are ranked lower than the ones on the respective left hand side with respect to a certain kind of total ordering. Moreover, the system is supposed to incorporate all integrability conditions in some sense discussed below.

These notes consist of two sections following the Introduction. Section 2 treats the problems outlined above for systems of linear PDEs, whereas Section 3 is dedicated to the more general case of systems of nonlinear PDEs. The discussion of the linear case leads to the notion of Janet basis. A basic variant of an algorithm computing Janet bases is outlined in Subsection 2.2, which builds on a method for partitioning certain sets of monomials into disjoint cones, as introduced in Subsection 2.1. The concept of Thomas decomposition is central for the nonlinear case. It is introduced for algebraic systems in Subsection 3.1 and is then adapted to differential systems in Subsection 3.2. The final Subsection 3.3 explains how to apply the Thomas decomposition technique for eliminating unknown functions from a system of nonlinear PDEs.

It is a non-trivial task to include here all relevant references. Among the most important historical ones we select: C. Méray [Mér80], C. Riquier [Riq10], M. Janet (1888-1983) [Jan29], J. M. Thomas (1898-1979) [Tho37, Tho62], J. F. Ritt [Rit34], [Rit50], E. R. Kolchin [Kol73] and A. Seidenberg [Sei56]. Related references are [Olv93], [Pom78], [Pom94], [Sch08a] and many more.

Closely related to the method of Thomas decomposition is the Rosenfeld-Gröbner algorithm and its implementation in the Maple package diffalg resp. DifferentialAlgebra (cf., e.g., [BLOP95], [BLOP09], [Hub97], [Hub00], [Bou]), but also the method of regular chains [LMMX05] and the rifsimp algorithm [RWB96]. Moreover, the notion of a characteristic set, introduced by J. F. Ritt and Wen-tsün Wu, again belongs to the same circle of ideas, cf., e.g., [Wu00], [Wu89] [Wan98], [Wan01], [Wan04], [Dio92]. Janet bases are related to Gröbner bases [Buc06, Buc87] and are particular instances of involutive bases [GB98a, GB98b, ZB96].

It is essential to note that the presented methods are also fundamental for further effective module-theoretic constructions for rings of linear functional operators and their implementations, on which applications, e.g., to systems theory are built (cf., e.g., [CQR05], [CQR07], [QR07], [CQ08], [CQ09], [Qua10a], [Qua10b], [QR14], [Rob15]). Efficient versions of the algorithms discussed in these notes have been implemented in Maple packages (Involutive, Janet, JanetOre, LDA, AlgebraicThomas, DifferentialThomas).

This exposition is based, in particular, on [Rob07], [Rob14], [Rob16], [LHR], [GLR19].

## 2. Systems of linear differential equations

In this section we assume that the given system of differential equations is linear (and homogeneous). In other words, for some $l, m, n \in \mathbb{N}$, some ring $D$ of differential operators, some matrix of operators $R \in D^{l \times m}$ and some left $D$-module $\mathcal{F}$ we can write the system as

$$
R u=0, \quad \text { where } \quad u=\left(\begin{array}{c}
u_{1}  \tag{2.1}\\
u_{2} \\
\vdots \\
u_{m}
\end{array}\right)
$$

for the unknown functions $u_{i}=u_{i}\left(z_{1}, \ldots, z_{n}\right) \in \mathcal{F}, i=1, \ldots, m$. The consequences of (2.1) are the left $D$-linear combinations of the rows of $R$, i.e., the elements of $D^{1 \times l} R$. (The functions in $\mathcal{F}$ need to be infinitely often differentiable at least.)

Example 2.1 ([Rob14], Ex. 3.2.49). An example of a system of linear PDEs with constant coefficients for one unknown function $u=u(x, y)$ of $x=z_{1}$ and $y=z_{2}$ is

$$
\left\{\begin{align*}
\frac{\partial^{2} u}{\partial x^{2}} & =0  \tag{2.2}\\
\frac{\partial^{2} u}{\partial y^{2}}+\frac{\partial u}{\partial x}+\frac{\partial u}{\partial y} & =0
\end{align*}\right.
$$

We may choose $D$ to be the commutative polynomial algebra $K\left[\partial_{x}, \partial_{y}\right]$, where $K$ is a field of characteristic zero (e.g., $\mathbb{Q}, \mathbb{R}, \mathbb{C}, \ldots$ ) and where $\partial_{x}$ and $\partial_{y}$ are the partial differential operators with respect to $x$ and $y$, respectively. The multiplication in $D$ is composition of operators.

Example 2.2 ([Rob14], Ex. 3.2.38). A system of linear PDEs with non-constant coefficients for $u=u(x, y)$ is given by

$$
\left\{\begin{array}{l}
\frac{\partial^{3} u}{\partial x \partial y^{2}}-\frac{\partial^{3} u}{\partial y^{3}}-(2 y+1) \frac{\partial^{2} u}{\partial y^{2}}-4 \frac{\partial u}{\partial y}=0 \\
\frac{\partial^{3} u}{\partial x^{2} \partial y}-\frac{\partial^{3} u}{\partial y^{3}}-2(2 y+1) \frac{\partial^{2} u}{\partial x \partial y}+\left(4 y^{2}+4 y-5\right) \frac{\partial u}{\partial y}=0
\end{array}\right.
$$

We may choose $K$ to be $\mathbb{Q}(x, y)$ or the field or meromorphic functions on some open and connected subset $\Omega$ of $\mathbb{C}^{2}$. Moreover, we let $D=K\left\langle\partial_{x}, \partial_{y}\right\rangle$ be the ring of differential operators

$$
\sum_{i=0}^{r} \sum_{j=0}^{s} a_{i, j} \partial_{x}^{i} \partial_{y}^{j}, \quad a_{i, j} \in K, \quad r, s \in \mathbb{Z}_{\geq 0}
$$

which are (skew) polynomials in $\partial_{x}$ and $\partial_{y}$, where composition is non-commutative in general.
Example 2.3 ([Rob14], Ex. 2.1.46). Linearizing the system of nonlinear PDEs

$$
\left\{\begin{align*}
\frac{\partial u}{\partial x}-u^{2} & =0  \tag{2.3}\\
\frac{\partial^{2} u}{\partial y^{2}}-u^{3} & =0
\end{align*}\right.
$$

for one unknown function $u$ of $x$ and $y$, we obtain the system of linear PDEs

$$
\left\{\begin{align*}
\frac{\partial U}{\partial x}-2 u U & =0  \tag{2.4}\\
\frac{\partial^{2} U}{\partial y^{2}}-3 u^{2} U & =0
\end{align*}\right.
$$

for one unknown function $U$ of $x$ and $y$, where $u$ is a solution of (2.3). In this case a preparatory treatment of the nonlinear system (2.3) is necessary to deal with the linearized system (2.4). The methods to be discussed in Section 3 allow to split system (2.3) into two systems

$$
\begin{aligned}
& \underline{u_{x}}-u^{2}=0 \\
&\left.2 \underline{u}_{y}^{2}-u^{4}, \partial_{y}\right\} \\
& u=0
\end{aligned}
$$

$$
\begin{array}{|cc|}
\hline \\
u=0 \quad\left\{\partial_{x}, \partial_{y}\right\} \\
\hline
\end{array}
$$

(where subscripts of $u$ indicate differentiation and where the meaning of the sets on the right will become clear later). The set of analytic solutions of the original system (2.3) is the disjoint union of the sets of analytic solutions of the above two systems. Choosing the first system in that splitting, we define the differential polynomial ring $R=\mathbb{Q}(\sqrt{2})\left[u, u_{x}, u_{y}, u_{x, x}, u_{x, y}, u_{y, y}, \ldots\right]$ and the ideal $I$ of $R$ which consists of all $R$-linear combinations of

$$
\begin{array}{llll}
u_{x}-u^{2}, & \partial_{x}\left(u_{x}-u^{2}\right), & \partial_{y}\left(u_{x}-u^{2}\right), & \partial_{x}^{2}\left(u_{x}-u^{2}\right), \\
u_{y}-\frac{\sqrt{2}}{2} u^{2}, & \partial_{x}\left(u_{y}-\frac{\sqrt{2}}{2} u^{2}\right), & \partial_{y}\left(u_{y}-\frac{\sqrt{2}}{2} u^{2}\right), & \partial_{x}^{2}\left(u_{y}-\frac{\sqrt{2}}{2} u^{2}\right),
\end{array} \ldots
$$

i.e., which consists of all $R$-linear combinations of the partial derivatives (of any order) of $u_{x}-u^{2}$ and $u_{y}-\frac{\sqrt{2}}{2} u^{2}$. Then $I$ is closed under partial differentiation, $R / I$ is an integral domain, and we may choose $K$ as the field of fractions of $R / I$, which is a differential field. In order to deal with system (2.4), we then define the skew polynomial ring $D=K\left\langle\partial_{x}, \partial_{y}\right\rangle$. (Instead of $u_{y}-\frac{\sqrt{2}}{2} u^{2}$ one may also choose $u_{y}+\frac{\sqrt{2}}{2} u^{2}$.)
Remark 2.4. An essential remark for what follows is that the given linear PDEs translate into linear equations for the Taylor coefficients $c_{i, j}$ of power series solutions

$$
u(x, y)=\sum_{i, j \geq 0} c_{i, j} \frac{\left(x-x_{0}\right)^{i}}{i!} \frac{\left(y-y_{0}\right)^{j}}{j!}, \quad\left(x_{0}, y_{0}\right) \in \mathbb{C}^{2}
$$

by substituting this ansatz into the PDEs and comparing coefficients (and similarly for a different number of independent variables and unknown functions). However, in order for the resulting system of linear equations in $c_{i, j}$ to characterize the power series solutions of the PDE system correctly (around a sufficiently generic point $\left(x_{0}, y_{0}\right)$ ), an overview of all consequences of the PDE system needs to be obtained first. Interesting new consequences are usually found by differentiating two known consequences so that in a suitable linear combination of these derivatives the highest derivatives of the unknown function cancel.

Example 2.5. Considering again Example 2.3, differentiation of the two PDEs in (2.4) yields

$$
\begin{aligned}
\frac{\partial^{2}}{\partial y^{2}}\left(\frac{\partial U}{\partial x}-2 u U\right) & =\frac{\partial^{3} U}{\partial x \partial y^{2}}-2\left(\frac{\partial^{2} u}{\partial y^{2}} U+2 \frac{\partial u}{\partial y} \frac{\partial U}{\partial y}+u \frac{\partial^{2} U}{\partial y^{2}}\right) \\
& =\frac{\partial^{3} U}{\partial x \partial y^{2}}-2 u^{3} U-2 \sqrt{2} u^{2} \frac{\partial U}{\partial y}-6 u^{3} U
\end{aligned}
$$

and

$$
\begin{aligned}
\frac{\partial}{\partial x}\left(\frac{\partial^{2} U}{\partial y^{2}}-3 u^{2} U\right) & =\frac{\partial^{3} U}{\partial x \partial y^{2}}-3\left(2 u \frac{\partial u}{\partial x} U+u^{2} \frac{\partial U}{\partial x}\right) \\
& =\frac{\partial^{3} U}{\partial x \partial y^{2}}-6 u^{3} U-6 u^{3} U
\end{aligned}
$$

Hence, we obtain

$$
\frac{\partial^{2}}{\partial y^{2}}\left(\frac{\partial U}{\partial x}-2 u U\right)-\frac{\partial}{\partial x}\left(\frac{\partial^{2} U}{\partial y^{2}}-3 u^{2} U\right)=4 u^{3} U-2 \sqrt{2} u^{2} \frac{\partial U}{\partial y}
$$

which yields the consequence

$$
\frac{\partial U}{\partial y}-\sqrt{2} u U=0
$$

This procedure for finding new consequences shall be studied now systematically. We start by considering the special case of ideals which are generated by monomials.
2.1. Monomial ideals. We denote by $\partial_{1}, \ldots, \partial_{n}$ the partial differential operators with respect to $z_{1}, \ldots, z_{n}$ and define the commutative polynomial algebra $D=K\left[\partial_{1}, \ldots, \partial_{n}\right]$ for some field $K$. For the sake of simplicity, we assume in this subsection that $\partial_{1}, \ldots, \partial_{n}$ act trivially on $K$, i.e., $K$ consists of constants, and that the differential equations involve one unknown function only.

The simplest operators in $D$ are given by monomials

$$
\partial^{J}:=\partial_{1}^{j_{1}} \ldots \partial_{n}^{j_{n}}, \quad \text { where } \quad J=\left(j_{1}, \ldots, j_{n}\right) \in\left(\mathbb{Z}_{\geq 0}\right)^{n}
$$

For $\mu \subseteq\left\{\partial_{1}, \ldots, \partial_{n}\right\}$ we consider the monoid

$$
\operatorname{Mon}(\mu):=\left\{\partial^{J} \mid J=\left(j_{1}, \ldots, j_{n}\right) \in\left(\mathbb{Z}_{\geq 0}\right)^{n}, j_{i}=0 \text { for all } i \text { such that } \partial_{i} \notin \mu\right\}
$$

with the usual divisibility relation $\mid$, and we let $\operatorname{Mon}(D):=\operatorname{Mon}\left(\left\{\partial_{1}, \ldots, \partial_{n}\right\}\right)$. An ideal of $D$ which is generated by monomials is called a monomial ideal.
Example 2.6 ([Rob14], Ex. 2.1.69). The system of linear PDEs

$$
\begin{equation*}
\frac{\partial^{2} u}{\partial x \partial y}=0, \quad \frac{\partial^{4} u}{\partial x^{3} \partial z}=0, \quad \frac{\partial^{4} u}{\partial x \partial y^{2} \partial z}=0, \quad \frac{\partial^{5} u}{\partial x^{2} \partial y \partial z^{2}}=0 \tag{2.5}
\end{equation*}
$$

for the unknown function $u=u(x, y, z)$ defines the monomial ideal $I$ of $K\left[\partial_{x}, \partial_{y}, \partial_{z}\right]$ which is generated by $\partial_{x} \partial_{y}, \partial_{x}^{3} \partial_{z}, \partial_{x} \partial_{y}^{2} \partial_{z}, \partial_{x}^{2} \partial_{y} \partial_{z}^{2}$. The ideal $I$ encodes all consequences of (2.5).

Remark 2.7. Let the ideal $I$ of $D$ be generated by monomials $m_{1}, \ldots, m_{r}$. Then every monomial in $I$ is a multiple of some $m_{i}$. The set of all monomials in $I$ is a multiple-closed subset of $\operatorname{Mon}(D)$ in the sense of the following definition.
Definition 2.8. A set $S \subseteq \operatorname{Mon}(D)$ is said to be $\operatorname{Mon}(\mu)$-multiple-closed, $\mu \subseteq\left\{\partial_{1}, \ldots, \partial_{n}\right\}$, if

$$
m s \in S \quad \text { for all } \quad m \in \operatorname{Mon}(\mu), \quad s \in S
$$

Every set $G \subseteq \operatorname{Mon}(D)$ satisfying

$$
\operatorname{Mon}(\mu) G=\{m g \mid m \in \operatorname{Mon}(\mu), g \in G\}=S
$$

is called a generating set for the $\operatorname{Mon}(\mu)$-multiple-closed set $S$.
Example 2.9. Let $D=K\left[\partial_{1}, \partial_{2}\right]$ and $G:=\left\{\partial_{1} \partial_{2}^{2}, \partial_{1}^{3} \partial_{2}, \partial_{1}^{4}\right\}$. We consider the $\operatorname{Mon}(D)$-multipleclosed set $S$ generated by $G$. If we visualize the monomial $\partial_{1}^{i} \partial_{2}^{j}$ as the point $(i, j)$ in the positive quadrant of a two-dimensional coordinate system, then the set $S$ of monomials can be viewed as the discrete set of points in the upper-right region in Figure 2.1.

The following combinatorial fact is also referred to as Dickson's Lemma.
Lemma 2.10. Every $\operatorname{Mon}(D)$-multiple-closed subset of $\operatorname{Mon}(D)$ has a finite generating set. Equivalently, every ascending chain of $\operatorname{Mon}(D)$-multiple-closed subsets of $\operatorname{Mon}(D)$ terminates.

In other words, every sequence of monomials in which no monomial has a divisor among the previous ones is finite.
Exercise. Prove Lemma 2.10 by induction on $n$.
Remark 2.11. Every multiple-closed set has a unique minimal generating set. It is obtained from any generating set $G$ by removing all elements which have a proper divisor in $G$.
Example 2.12. The multiple-closed set generated by $\partial_{x} \partial_{y}, \partial_{x}^{3} \partial_{z}, \partial_{x} \partial_{y}^{2} \partial_{z}, \partial_{x}^{2} \partial_{y} \partial_{z}^{2}$ in Example 2.6 has minimal generating set $\left\{\partial_{x} \partial_{y}, \partial_{x}^{3} \partial_{z}\right\}$.


Figure 2.1: $\operatorname{Mon}(D)$-multiple-closed set $S$ in Example 2.9


Figure 2.2: Cone decomposition of $S$ in Example 2.14

We are going to partition multiple-closed sets (and, more importantly, their complements in $\operatorname{Mon}(D)$ ) into cones of monomials, where a cone is a $\operatorname{Mon}(\mu)$-multiple-closed set generated by one monomial, for some $\mu \subseteq\left\{\partial_{1}, \ldots, \partial_{n}\right\}$. For a set $S$ let $\mathcal{P}(S)$ be its power set.

Definition 2.13. (a) A pair $(C, \mu) \in \mathcal{P}(\operatorname{Mon}(D)) \times \mathcal{P}\left(\left\{\partial_{1}, \ldots, \partial_{n}\right\}\right)$ is called a cone if there exists $v \in C$ such that

$$
\operatorname{Mon}(\mu) v=\{m v \mid m \in \operatorname{Mon}(\mu)\}=C
$$

The elements of $\mu$ are called the multiplicative variables, those of $\bar{\mu}:=\left\{\partial_{1}, \ldots, \partial_{n}\right\} \backslash \mu$ the non-multiplicative variables for $(C, \mu)$ (or simply for $C$, or for $v$ ). We often also refer to the cone $C$ by the pair $(v, \mu)$, where $v$ is the generator of $C$.
(b) Let $S \subseteq \operatorname{Mon}(D)$. A cone decomposition of $S$ is a finite set $\left\{\left(m_{1}, \mu_{1}\right), \ldots,\left(m_{r}, \mu_{r}\right)\right\}$ of cones such that the sets $C_{1}, \ldots, C_{r}$ defined by $C_{i}:=\operatorname{Mon}\left(\mu_{i}\right) m_{i}$ satisfy

$$
\bigcup_{i=1}^{r} C_{i}=S \quad \text { and } \quad C_{i} \cap C_{j}=\emptyset \quad \text { for all } \quad i \neq j
$$

Example 2.14. A cone decomposition of the multiple-closed set $S$ defined in Example 2.9 is

$$
\left\{\left(\partial_{1}^{4},\left\{\partial_{1}, \partial_{2}\right\}\right), \quad\left(\partial_{1}^{3} \partial_{2},\left\{\partial_{2}\right\}\right), \quad\left(\partial_{1}^{2} \partial_{2}^{2},\left\{\partial_{2}\right\}\right), \quad\left(\partial_{1} \partial_{2}^{2},\left\{\partial_{2}\right\}\right)\right\},
$$

which is visualized in Figure 2.2.
Remark 2.15. A cone decomposition of $S \subseteq \operatorname{Mon}(D)$ defines a restriction of the usual divisibility relation of monomials as follows. A monomial $m \in \operatorname{Mon}(D)$ is divisible by a generator $m^{\prime}$ of a cone $\left(m^{\prime}, \mu\right)$ if and only if there exists $m^{\prime \prime} \in \operatorname{Mon}(\mu)$ such that $m=m^{\prime \prime} m^{\prime}$. The disjointness of the cone decomposition entails that among cone generators there is at most one divisor.

Given a finite set $\left\{m_{1}, \ldots, m_{r}\right\}$ of monomials, there are many possible ways of how to arrange sets of multiplicative variables $\mu_{1}, \ldots, \mu_{r}$ such that $\left\{\left(m_{1}, \mu_{1}\right), \ldots,\left(m_{r}, \mu_{r}\right)\right\}$ is a set of disjoint cones. These possibilities are addressed by the notion of involutive division which was introduced
by V. P. Gerdt, Y. A. Blinkov, A. Y. Zharkov [GB98a, GB98b, ZB96], cf. also [Ape98]. Important for this exposition is only the Janet division:
Definition 2.16. For a finite subset $M$ of $\operatorname{Mon}(D)$ Janet division defines the set $\mu=\mu(m, M)$ of multiplicative variables for each $m \in M$ as follows. If $m=\partial_{1}^{i_{1}} \ldots \partial_{n}^{i_{n}}$, then for $1 \leq k \leq n$

$$
\partial_{k} \in \mu \quad: \Longleftrightarrow \quad i_{k}=\max \left\{j_{k} \mid \partial_{1}^{j_{1}} \ldots \partial_{n}^{j_{n}} \in M \text { with } j_{1}=i_{1}, j_{2}=i_{2}, \ldots, j_{k-1}=i_{k-1}\right\}
$$

This definition assumes the ordering $\partial_{1}, \partial_{2}, \ldots, \partial_{n}$ of the variables; a different ordering may be used as well. There are also other common involutive divisions. For instance, J. M. Thomas [Tho37] proposed another way of defining the multiplicative variables of cones; still another one is named after J.-F. Pommaret (cf., e.g., [Jan29, no. 58], [Pom94, p. 90], [Sei10]).
Example 2.17. For $M=\left\{\partial_{1}^{2} \partial_{2}, \partial_{1}^{2} \partial_{3}, \partial_{2}^{2} \partial_{3}, \partial_{2} \partial_{3}^{2}\right\}$ Janet division associates the sets $\mu(m, M)$ of multiplicative variables to the elements $m \in M$ as indicated in the following table, where we replace non-multiplicative variables in the set $\left\{\partial_{1}, \ldots, \partial_{n}\right\}$ with the symbol ' $*$ '.

$$
\begin{array}{ll}
\partial_{1}^{2} \partial_{2}, & \left\{\partial_{1}, \partial_{2}, \partial_{3}\right\} \\
\partial_{1}^{2} \partial_{3}, & \left\{\partial_{1}, *, \partial_{3}\right\} \\
\partial_{2}^{2} \partial_{3}, & \left\{*, \partial_{2}, \partial_{3}\right\} \\
\partial_{2} \partial_{3}^{2}, & \left\{*, *, \partial_{3}\right\}
\end{array}
$$

Definition 2.18. A finite subset $M$ of $\operatorname{Mon}(D)$ is said to be Janet complete if

$$
\bigcup_{m \in M} \operatorname{Mon}(\mu(m, M)) m=\bigcup_{m \in M} \operatorname{Mon}(D) m
$$

i.e., if every monomial that is divisible by some monomial in $M$ is obtained by multiplying a certain $m \in M$ by multiplicative variables for $m$ only. (The left hand side above is a disjoint union.)
Example 2.19. The set $M$ in Example 2.17 is not Janet complete because, e.g., the monomial $\partial_{1} \partial_{2}^{2} \partial_{3}$ is not obtained as a multiple of any $m \in M$ when multiplication is restricted to multiplicative variables for $m$. By adding this monomial and the monomial $\partial_{1} \partial_{2} \partial_{3}^{2}$ to $M$, we obtain the following Janet complete superset of $M$ in $\operatorname{Mon}(D)$.

$$
\begin{aligned}
\partial_{1}^{2} \partial_{2}, & \left\{\partial_{1}, \partial_{2}, \partial_{3}\right\} \\
\partial_{1}^{2} \partial_{3}, & \left\{\partial_{1}, *, \partial_{3}\right\} \\
\partial_{1} \partial_{2}^{2} \partial_{3}, & \left\{*, \partial_{2}, \partial_{3}\right\} \\
\partial_{1} \partial_{2} \partial_{3}^{2}, & \left\{*, *, \partial_{3}\right\} \\
\partial_{2}^{2} \partial_{3}, & \left\{*, \partial_{2}, \partial_{3}\right\} \\
\partial_{2} \partial_{3}^{2}, & \left\{*, *, \partial_{3}\right\}
\end{aligned}
$$

Remark 2.20. Every finite subset $M$ of $\operatorname{Mon}(D)$ can be enlarged to a finite Janet complete set by including multiples of certain $m \in M$ using multiplicative and non-multiplicative variables.
Exercise. Determine a finite Janet complete superset of $\left\{\partial_{1}^{3} \partial_{2}, \partial_{1} \partial_{3}^{2} \partial_{4}^{2}, \partial_{2} \partial_{3} \partial_{4}\right\}$.
Proposition 2.21. Let $I$ be a monomial ideal of $D$. Let an ordering of $\partial_{1}, \ldots, \partial_{n}$ be fixed. There exists a unique finite Janet complete generating set of monomials for $I$ which is minimal with respect to set inclusion.

Definition 2.22. Let $S \subseteq \operatorname{Mon}(D)$. We refer to the minimal Janet complete superset of $S$ as the Janet completion of $S$. Its elements are the generators of cones in a cone decomposition of the multiple-closed set generated by $S$. We call this cone decomposition the Janet decomposition of the multiple-closed set generated by $S$.

Exercise. Write an algorithm which computes the Janet completion of a finite set of monomials.
Cone decompositions of the complement of a multiple-closed set in $\operatorname{Mon}(D)$ which are defined by Janet division will be referred to as Janet decompositions as well.
Exercise. Write an algorithm which computes a Janet decomposition of the complement of a multiple-closed set of monomials in $\operatorname{Mon}(D)$.

Definition 2.23. For any set $S \subseteq \operatorname{Mon}(D)$ of monomials, the generalized Hilbert series of $S$ is the formal power series

$$
H_{S}\left(\partial_{1}, \ldots, \partial_{n}\right):=\sum_{m \in S} m \in \mathbb{Z}\left[\left[\partial_{1}, \ldots, \partial_{n}\right]\right]
$$

Remark 2.24. The Hilbert series arising in commutative algebra for subsets of homogeneous polynomials in a polynomial ring with standard grading is obtained from the generalized Hilbert series as $H_{S}(\lambda, \ldots, \lambda)$ for an indeterminate $\lambda$.

Remark 2.25. Let $C=(m, \mu)$ be a cone, where $m \in \operatorname{Mon}(D)$ and $\mu \subseteq\left\{\partial_{1}, \ldots, \partial_{n}\right\}$. We use the geometric series

$$
\frac{1}{1-x}=\sum_{i \geq 0} x^{i}
$$

to write down the generalized Hilbert series $H_{C}\left(\partial_{1}, \ldots, \partial_{n}\right)$ as follows:

$$
H_{C}\left(\partial_{1}, \ldots, \partial_{n}\right)=\frac{m}{\prod_{\partial \in \mu}(1-\partial)}
$$

More generally, every decomposition of a $\operatorname{Mon}(D)$-multiple-closed set $S$ into disjoint cones allows to compute the generalized Hilbert series of $S$ by adding the generalized Hilbert series of the cones. In an analogous way this applies to the complements of multiple-closed sets.

Example 2.26. The complement in $\operatorname{Mon}(D)$ of the multiple-closed set generated by $\partial_{x} \partial_{y}, \partial_{x}^{3} \partial_{z}$, $\partial_{x} \partial_{y}^{2} \partial_{z}, \partial_{x}^{2} \partial_{y} \partial_{z}^{2}$ in Example 2.6 admits the following Janet decomposition:

$$
\begin{aligned}
1, & \left\{*, \partial_{y}, \partial_{z}\right\} \\
\partial_{x}, & \left\{*, *, \partial_{z}\right\} \\
\partial_{x}^{2}, & \left\{*, *, \partial_{z}\right\} \\
\partial_{x}^{3}, & \left\{\partial_{x}, *, *\right\}
\end{aligned}
$$

The corresponding generalized Hilbert series is

$$
\frac{1}{\left(1-\partial_{y}\right)\left(1-\partial_{z}\right)}+\frac{\partial_{x}}{1-\partial_{z}}+\frac{\partial_{x}^{2}}{1-\partial_{z}}+\frac{\partial_{x}^{3}}{1-\partial_{x}} .
$$

2.2. Janet's algorithm. Given a system of linear PDEs, Janet's algorithm computes an equivalent system, called a Janet basis, for which it is a straightforward task to decide whether another linear PDE is a consequence of the system or not. The answer is obtained by trying to express the PDE as a linear combination of partial derivatives (of any order) of the Janet basis elements. This process is based on a multivariate polynomial division for elements of $D=K\left[\partial_{1}, \ldots, \partial_{n}\right]$, which requires a choice of most significant term in each non-zero polynomial, called leading term.

Suppose that a total ordering $>$ on $\operatorname{Mon}(D)$ is chosen which is compatible with multiplication (i.e., composition of operators). By defining leading terms of PDEs with respect to $>$, the leading terms of consequences of one PDE are predictable: the leading term of a derivative of a PDE is the derivative of the leading term of the PDE.

A total ordering > with the above property also enables us to easily determine the monomials in $\partial_{1}, \ldots, \partial_{n}$ that do not occur in leading terms of consequences of a system of linear PDEs. Hence, a Janet basis then also allows to determine all analytic solutions (around a sufficiently generic point). By choosing the total ordering $>$ appropriately, further tasks, e.g., elimination of variables, can be solved as well.

The methods to be discussed in this section can be applied in a similar way to other types of linear equations, e.g., difference equations, multidimensional discrete equations, time-delay equations and other functional equations. The coefficients of these equations may be constant or not, corresponding to commutative or non-commutative rings of operators, e.g., Ore algebras (cf., e.g., [CS98], [CQR05]). For example, singular points of differential equations may be studied in terms of $D$-modules [Kas03, Cou95], i.e., modules over Weyl algebras and related rings of differential operators.

Last but not least, Janet's algorithm applies in the same way to systems of polynomial equations, i.e., equations defining algebraic varieties. Hence, it is an alternative to Buchberger's algorithm computing Gröbner bases. In fact, every Janet basis is a Gröbner basis. Generalizations of Gröbner bases to non-commutative polynomial algebras have been studied since a couple of decades, cf., e.g., [KRW90], [Kre93], [Mor94], [Lev05], [GL11]; for rings of differential operators, cf., e.g., [CJ84], [Gal85], [IP98], [SST00]. Buchberger's algorithm was adapted to Ore algebras by F. Chyzak (cf. [Chy98], [CS98], where it is also applied to the study of special functions and combinatorial sequences). Involutive divisions were studied for the Weyl algebra in [HSS02] and were extended to non-commutative rings in [EW07].

In this section we confine ourselves to linear PDEs with constant coefficients, but these may involve $q$ unknown functions. Note that we ignore efficiency issues in favor of a concise formulation of Janet's algorithm.

Let $D=K\left[\partial_{1}, \ldots, \partial_{n}\right]$, where $K$ is a field of constants, $q \in \mathbb{N}$ and $e_{1}, \ldots, e_{q}$ the standard basis vectors of the free (left) $D$-module $D^{1 \times q}$. We define the set of monomials of $D^{1 \times q}$ to be

$$
\operatorname{Mon}\left(D^{1 \times q}\right):=\bigcup_{i=1}^{q} \operatorname{Mon}(D) e_{i}
$$

Every $p \in D^{1 \times q}$ has a unique representation

$$
\begin{equation*}
p=\sum_{k=1}^{q} \sum_{m \in \operatorname{Mon}(D)} c_{k, m} m e_{k} \tag{2.6}
\end{equation*}
$$

as linear combination of monomials in $\operatorname{Mon}\left(D^{1 \times q}\right)$ with coefficients $c_{k, m} \in K$, where only finitely many $c_{k, m}$ are non-zero.

Definition 2.27. A term ordering $>$ on $\operatorname{Mon}\left(D^{1 \times q}\right)\left(\right.$ or on $\left.D^{1 \times q}\right)$ is a total ordering on $\operatorname{Mon}\left(D^{1 \times q}\right)$ which satisfies the following two conditions.
(a) For all $1 \leq j \leq n$ and $1 \leq k \leq q$ we have $\partial_{j} e_{k}>e_{k}$.
(b) For all $m_{1} e_{k}, m_{2} e_{l} \in \operatorname{Mon}\left(D^{1 \times q}\right)$ the following implication holds:

$$
m_{1} e_{k}>m_{2} e_{l} \quad \Longrightarrow \quad \partial_{j} m_{1} e_{k}>\partial_{j} m_{2} e_{l} \quad \text { for all } j=1, \ldots, n
$$

Let a term ordering $>$ be fixed. For every $p \in D^{1 \times q} \backslash\{0\}$ the greatest monomial, with respect to $>$, occurring (with non-zero coefficient) in the representation (2.6) of $p$ is uniquely determined and is called the leading monomial of $p$, denoted by $\operatorname{lm}(p)$. The coefficient of $\operatorname{lm}(p)$ is called the leading coefficient of $p$, denoted by lc $(p)$. For any subset $S \subseteq D^{1 \times q}$ we define

$$
\operatorname{lm}(S):=\{\operatorname{lm}(p) \mid 0 \neq p \in S\}
$$

Remark 2.28. Every term ordering on $D^{1 \times q}$ is a well-ordering, i.e., every descending sequence of elements of $\operatorname{Mon}\left(D^{1 \times q}\right)$ terminates.
Example 2.29. The lexicographical ordering (lex) on $\operatorname{Mon}(D)$ (which extends the ordering $\partial_{1}>$ $\left.\partial_{2}>\ldots>\partial_{n}\right)$ is defined for monomials $m_{1}=\partial_{1}^{a_{1}} \ldots \partial_{n}^{a_{n}}, m_{2}=\partial_{1}^{b_{1}} \ldots \partial_{n}^{b_{n}} \in \operatorname{Mon}(D)$ by

$$
m_{1}>m_{2} \quad: \Longleftrightarrow m_{1} \neq m_{2} \quad \text { and } \quad a_{j}>b_{j} \quad \text { for } \quad j=\min \left\{1 \leq i \leq n \mid a_{i} \neq b_{i}\right\}
$$

Example 2.30. The degree-reverse lexicographical ordering (degrevlex) on $\operatorname{Mon}(D)$ (extending the ordering $\left.\partial_{1}>\partial_{2}>\ldots>\partial_{n}\right)$ is defined for $m_{1}=\partial_{1}^{a_{1}} \ldots \partial_{n}^{a_{n}}, m_{2}=\partial_{1}^{b_{1}} \ldots \partial_{n}^{b_{n}} \in \operatorname{Mon}(D)$ by

$$
m_{1}>m_{2} \quad \Longleftrightarrow\left\{\begin{array}{l}
\operatorname{deg}\left(m_{1}\right)>\operatorname{deg}\left(m_{2}\right) \text { or } \\
\left(\operatorname{deg}\left(m_{1}\right)=\operatorname{deg}\left(m_{2}\right) \text { and } m_{1} \neq m_{2} \text { and } a_{j}<b_{j}\right. \\
\text { for } \left.j=\max \left\{1 \leq i \leq n \mid a_{i} \neq b_{i}\right\}\right),
\end{array}\right.
$$

where deg refers to the total degree.
Example 2.31. Two ways of extending a given term ordering $>_{1}$ on $\operatorname{Mon}(D)$ to $\operatorname{Mon}\left(D^{1 \times q}\right)$ for $q>1$ are often used. The term-over-position ordering (extending $>_{1}$ and the total ordering $e_{1}>e_{2}>\ldots>e_{q}$ of the standard basis vectors) is defined for $m_{1}, m_{2} \in \operatorname{Mon}(D)$ by

$$
m_{1} e_{i}>m_{2} e_{j} \quad: \Longleftrightarrow \quad m_{1}>_{1} m_{2} \quad \text { or } \quad\left(m_{1}=m_{2} \quad \text { and } \quad i<j\right)
$$

Accordingly, the position-over-term ordering (extending $>_{1}$ and $e_{1}>e_{2}>\ldots>e_{q}$ ) is defined by

$$
m_{1} e_{i}>m_{2} e_{j} \quad: \Longleftrightarrow \quad i<j \quad \text { or } \quad\left(i=j \quad \text { and } \quad m_{1}>_{1} m_{2}\right)
$$

In what follows, we assume that a term ordering $>$ on $D^{1 \times q}$ is fixed.
Let $M$ be a submodule of $D^{1 \times q}$. Note that $\operatorname{lm}(M)$ is a $\operatorname{Mon}(D)$-multiple-closed set. More precisely, for each $k \in\{1, \ldots, q\}$, the set

$$
\left\{m \in \operatorname{Mon}(D) \mid m e_{k} \in \operatorname{lm}(M)\right\}
$$

is $\operatorname{Mon}(D)$-multiple-closed as discussed in Section 2.1.
Starting with a finite generating set $L$ of $M$, Janet's algorithm possibly removes elements from $L$ and inserts new elements of $M$ into $L$ repeatedly in order to finally achieve that the $\operatorname{Mon}(D)$ -multiple-closed set generated by $\operatorname{lm}(L)$ equals $\operatorname{lm}(M)$. An element $p \in L$ is removed if it is reduced to zero by subtraction of suitable multiples of other elements of $L$.

We denote by ${ }_{D}\langle L\rangle$ the submodule of $D^{1 \times q}$ generated by $L \subseteq D^{1 \times q}$.
For $G \subseteq \operatorname{Mon}\left(D^{1 \times q}\right)$ we denote by $[G]$ the $\operatorname{Mon}(D)$-multiple-closed set generated by $G$. If $G=\left\{m_{1}, \ldots, m_{r}\right\}$, then we also write $\left[m_{1}, \ldots, m_{r}\right]$ for $[G]$.

Definition 2.32. Let $T=\left\{\left(b_{1}, \mu_{1}\right), \ldots,\left(b_{t}, \mu_{t}\right)\right\}$, where $b_{i} \in D^{1 \times q} \backslash\{0\}$ and $\mu_{i} \subseteq\left\{\partial_{1}, \ldots, \partial_{n}\right\}$.
(a) The set $T$ is Janet complete if $\left\{\operatorname{lm}\left(b_{1}\right), \ldots, \operatorname{lm}\left(b_{t}\right)\right\}$ equals its Janet completion and, for each $i \in\{1, \ldots, t\}, \mu_{i}$ is the set of multiplicative variables of the cone with generator $\operatorname{lm}\left(b_{i}\right)$ in the Janet decomposition $\left\{\left(\operatorname{lm}\left(b_{1}\right), \mu_{1}\right), \ldots,\left(\operatorname{lm}\left(b_{t}\right), \mu_{t}\right)\right\}$ of $\left[\operatorname{lm}\left(b_{1}\right), \ldots, \operatorname{lm}\left(b_{t}\right)\right]$.
(b) An element $p \in D^{1 \times q}$ is Janet reducible modulo $T$ if there exist $(b, \mu) \in T$ and a monomial $m \in \operatorname{Mon}\left(D^{1 \times q}\right)$ which occurs in $p$ such that $m \in \operatorname{Mon}(\mu) \operatorname{lm}(b)$. In this case, $(b, \mu)$ is called a Janet divisor of $p$. If $p$ is not Janet reducible modulo $T$, then $p$ is also said to be Janet reduced modulo $T$.

The following algorithm subtracts suitable multiples of Janet divisors from a given element $p \in D^{1 \times q}$ as long as a term in $p$ is Janet reducible modulo $T$.
Algorithm 2.33 (Janet-reduce).
Input: $p \in D^{1 \times q}, T=\left\{\left(b_{1}, \mu_{1}\right), \ldots,\left(b_{t}, \mu_{t}\right)\right\}$, and a term ordering $>$ on $D^{1 \times q}$, where $T$ is Janet
complete (with respect to $>$, cf. Definition 2.32)
Output: $r \in D^{1 \times q}$ such that

$$
r+{ }_{D}\left\langle b_{1}, \ldots, b_{t}\right\rangle=p+{ }_{D}\left\langle b_{1}, \ldots, b_{t}\right\rangle
$$

and $r$ is Janet reduced modulo $T$

```
Algorithm:
    \(p^{\prime} \leftarrow p ; \quad r \leftarrow 0\)
    while \(p^{\prime} \neq 0\) do
        if \(\exists(b, \mu) \in T\) such that \(\operatorname{lm}\left(p^{\prime}\right) \in \operatorname{Mon}(\mu) \operatorname{lm}(b)\) then \(\quad / /(b, \mu)\) is a Janet divisor of \(p^{\prime}\)
        \(p^{\prime} \leftarrow p^{\prime}-\frac{\operatorname{lc}\left(p^{\prime}\right)}{\operatorname{lc}(b)} \frac{\operatorname{lm}\left(p^{\prime}\right)}{\operatorname{lm}(b)} b\)
    else
            subtract the term of \(p^{\prime}\) with monomial \(\operatorname{lm}\left(p^{\prime}\right)\) from \(p^{\prime}\) and add it to \(r\)
    end if
    end while
    return \(r\)
```

Remark 2.34. Algorithm 2.33 terminates because, as long as $p^{\prime}$ is non-zero, the leading monomial of $p^{\prime}$ decreases with respect to the term ordering $>$, which is a well-ordering. Its correctness is clear. The result $r$ is uniquely determined for the given input because every monomial has at most one Janet divisor in $T$, and also the course of Algorithm 2.33 is uniquely determined as opposed to reduction procedures which apply multivariate polynomial division without distinguishing between multiplicative and non-multiplicative variables.

Remark 2.35. Let $p_{1}, p_{2} \in D^{1 \times q}$ and $T$ be as in the input of Algorithm 2.33. In general, the equality $p_{1}+{ }_{D}\left\langle b_{1}, \ldots, b_{t}\right\rangle=p_{2}+{ }_{D}\left\langle b_{1}, \ldots, b_{t}\right\rangle$ does not imply that the results of applying Janetreduce to $p_{1}$ and $p_{2}$, respectively, are equal. However, according to Theorem 2.39 (d) below, if $T$ is a Janet basis, then the result of Janet-reduce constitutes a unique representative for every coset in $D^{1 \times q} /{ }_{D}\left\langle b_{1}, \ldots, b_{t}\right\rangle$. It is called the Janet normal form of $p_{1}$ (or $p_{2}$ ) modulo $T$.
Definition 2.36. Let $T=\left\{\left(b_{1}, \mu_{1}\right), \ldots,\left(b_{t}, \mu_{t}\right)\right\}$ be Janet complete (as in Definition 2.32 (a)). We write $\operatorname{NF}(p, T,>)$ for the result of Algorithm 2.33 (Janet-reduce) applied to $p, T,>$. The set $T$ is said to be passive if

$$
\begin{equation*}
\mathrm{NF}\left(v \cdot b_{i}, T,>\right)=0 \quad \text { for all } \quad v \in \overline{\mu_{i}}, \quad i=1, \ldots, t \tag{2.7}
\end{equation*}
$$

In this case $T$ is also called a Janet basis for ${ }_{D}\left\langle b_{1}, \ldots, b_{t}\right\rangle$ (with respect to $\rangle$ ), and $\left\{b_{1}, \ldots, b_{t}\right\}$ is often referred to as a Janet basis for ${ }_{D}\left\langle b_{1}, \ldots, b_{t}\right\rangle$ as well.
Remark 2.37. Let $M$ be a submodule of $D^{1 \times q}$. By applying Janet's algorithm to a finite generating set $L$ of $M$, an ascending chain of multiple-closed subsets of $\operatorname{lm}(M)$ is constructed. This chain terminates by Lemma 2.10. In each round, a Janet decomposition is computed for the current multiple-closed set generated by the leading monomials of a generating set for $M$. In order to obtain the minimal Janet complete set of monomials, the generating set for $M$ is first turned into an auto-reduced one, i.e., no leading monomial of a generator divides (in the conventional sense) the leading monomial of another generator.

Let $M={ }_{D}\left\langle b_{1}, \ldots, b_{t}\right\rangle$. Then every element of $M$ is a $D$-linear combination of $b_{1}, \ldots, b_{t}$. Suppose that $T$ is passive. Each summand $c_{i} m_{i} b_{i}$ in such a linear combination, where $c_{i} \in K$ and where $m_{i} \in \operatorname{Mon}(D)$ involves some variable that is non-multiplicative for $b_{i}$, can be replaced with a $K$-linear combination of elements in $\operatorname{Mon}\left(\mu_{1}\right) b_{1}, \ldots, \operatorname{Mon}\left(\mu_{t}\right) b_{t}$. Due to the passivity condition (2.7), this can be achieved by applying successively Algorithm 2.33 (Janet-reduce) to terms involving only one non-multiplicative variable. This substitution process should deal with the largest term with respect to $>$ first. Elimination of all non-multiplicative variables demonstrates that the leading monomial of every $p \in M \backslash\{0\}$ has a Janet divisor in $T$. We conclude that passivity of the Janet complete set $T$ is equivalent to $\left[\operatorname{lm}\left(b_{1}\right), \ldots, \operatorname{lm}\left(b_{t}\right)\right]=\operatorname{lm}(M)$.

Recall that for any set $S$ we denote by $\mathcal{P}(S)$ the power set of $S$.
Algorithm 2.38 (JanetBasis).
Input: A finite set $L \subseteq D^{1 \times q}$, a term ordering $>$ on $D^{1 \times q}$, and an ordering of $\partial_{1}, \ldots, \partial_{n}$ for Janet division
Output: A finite subset $J$ of $D^{1 \times q} \times \mathcal{P}\left(\left\{\partial_{1}, \ldots, \partial_{n}\right\}\right)$ such that ${ }_{D}\langle p \mid(p, \mu) \in J\rangle={ }_{D}\langle L\rangle$ (and $J=\emptyset$ if and only if $\left.{ }_{D}\langle L\rangle=\{0\}\right)$

## Algorithm:

$: G \leftarrow L$
repeat

$$
G \leftarrow \operatorname{Auto-reduce}(G,>) \quad / / \text { cf. Remark } 2.37
$$

$$
J \leftarrow \operatorname{Janet}-\operatorname{decompose}(G) \quad / / \text { cf. second exercise on } p .7
$$

$$
P \leftarrow\{\operatorname{NF}(v \cdot p, J,>) \mid(p, \mu) \in J, v \in \bar{\mu}\} \quad \text { // cf. Algorithm } 2.33
$$

$$
G \leftarrow\{p \mid(p, \mu) \in J\} \cup P
$$

until $P \subseteq\{0\}$
return $J$
Theorem 2.39 ([Rob14], Thm. 2.1.43).
(a) Algorithm 2.38 terminates and is correct.
(b) AK-basis of ${ }_{D}\langle L\rangle$ is given by $\biguplus \operatorname{Mon}(\mu) p$, where $J$ is the result of Algorithm 2.38. $(p, \mu) \in J$
In particular, every $r \in{ }_{D}\langle L\rangle$ has a unique representation

$$
r=\sum_{(p, \mu) \in J} c_{(p, \mu)} p
$$

where each $c_{(p, \mu)} \in D$ is a $K$-linear combination of elements in $\operatorname{Mon}(\mu)$.
(c) The cosets in $D^{1 \times q} /{ }_{D}\langle L\rangle$ with representatives in

$$
\operatorname{Mon}\left(D^{1 \times q}\right) \backslash[\operatorname{lm}(p) \mid(p, \mu) \in J]
$$

form a $K$-basis of $D^{1 \times q} /{ }_{D}\langle L\rangle$.
(d) For every $r_{1}, r_{2} \in D^{1 \times q}$ the following equivalence holds.

$$
r_{1}+{ }_{D}\langle L\rangle=r_{2}+{ }_{D}\langle L\rangle \Longleftrightarrow \mathrm{NF}\left(r_{1}, J,>\right)=\mathrm{NF}\left(r_{2}, J,>\right) .
$$

We present a small example illustrating the idea of Janet's algorithm.
Example 2.40. Let $D=K\left[\partial_{1}, \partial_{2}\right]$ be the commutative polynomial algebra in $\partial_{1}, \partial_{2}$ over a field $K$. We choose the degree-reverse lexicographical ordering on $\operatorname{Mon}(D)$ satisfying $\partial_{1}>\partial_{2}$ (cf. Example 2.30). Let the ideal $I$ of $D$ be generated by

$$
p_{1}:=\underline{\partial_{1}^{2}}-\partial_{2}, \quad p_{2}:=\underline{\partial_{1} \partial_{2}}-\partial_{2} .
$$

Using the ordering $\partial_{1}, \partial_{2}$ of the variables for Janet division, the Janet decomposition of the multiple-closed set which is generated by the underlined leading monomials of $p_{1}$ and $p_{2}$ is

$$
\left\{\left(\partial_{1}^{2},\left\{\partial_{1}, \partial_{2}\right\}\right),\left(\partial_{1} \partial_{2},\left\{\partial_{2}\right\}\right)\right\}
$$

This result indicates that we need to check whether $f:=\partial_{1} p_{2}$ can be written as

$$
\begin{equation*}
f=c_{1}\left(\partial_{1}^{2}-\partial_{2}\right)+c_{2}\left(\partial_{1} \partial_{2}-\partial_{2}\right), \quad c_{1} \in K\left[\partial_{1}, \partial_{2}\right], \quad c_{2} \in K\left[\partial_{2}\right] \tag{2.8}
\end{equation*}
$$

The monomials appearing in $f=\partial_{1}^{2} \partial_{2}-\partial_{1} \partial_{2} \in I$ lie in the cones $\left(\partial_{1}^{2},\left\{\partial_{1}, \partial_{2}\right\}\right)$ and $\left(\partial_{1} \partial_{2},\left\{\partial_{2}\right\}\right)$, respectively. Reduction yields $p_{3}:=\partial_{2}^{2}-\partial_{2} \in I$, which does not have a representation as in (2.8). So, we include $p_{3}$ in our list of generators, and for this example, we already arrive at the (minimal) Janet basis $\left\{\left(p_{1},\left\{\partial_{1}, \partial_{2}\right\}\right),\left(p_{2},\left\{\partial_{2}\right\}\right),\left(p_{3},\left\{\partial_{2}\right\}\right)\right\}$ for $I$.

Remark 2.41. The $K$-vector space $\mathcal{F}:=\operatorname{hom}_{K}(D, K)$ is a (left) $D$-module with action

$$
D \times \mathcal{F} \longrightarrow \mathcal{F}:(d, f) \longmapsto(a \mapsto f(a \cdot d)),
$$

and the following $K$-bilinear form is non-degenerate in both arguments:

$$
\begin{equation*}
(, \quad): D \times \mathcal{F} \longrightarrow K:(d, f) \longmapsto f(d) . \tag{2.9}
\end{equation*}
$$

Hence, $D$ and $\mathcal{F}$ are dual to each other. The linear map $D \rightarrow D$ defined by right multiplication by $d$ and the linear map $\mathcal{F} \rightarrow \mathcal{F}$ given by left multiplication by $d$ are adjoint to each other:

$$
\begin{equation*}
(a \cdot d, f)=f(a \cdot d)=(d \cdot f)(a)=(a, d \cdot f), \quad a \in D, \quad f \in \mathcal{F} \tag{2.10}
\end{equation*}
$$

Since every homomorphism $f \in \mathcal{F}$ is uniquely determined by its values for the elements of the $K$-basis $\operatorname{Mon}(D)$ of $D$, we can write $f$ in a unique way as a (not necessarily finite) formal sum

$$
\begin{equation*}
\sum_{m \in \operatorname{Mon}(D)}(m, f) m \tag{2.11}
\end{equation*}
$$

Due to (2.10), for every $d \in D$ the representation of $d \cdot f$ can be obtained from

$$
\begin{equation*}
\sum_{m \in \operatorname{Mon}(D)}(m, d \cdot f) m=\sum_{m \in \operatorname{Mon}(D)}(m \cdot d, f) m \tag{2.12}
\end{equation*}
$$

By writing the monomials in the sum (2.11) in indeterminates $z_{1}, \ldots, z_{n}$, we identify $\mathcal{F}$ with the $K$-algebra $K\left[\left[z_{1}, \ldots, z_{n}\right]\right]$ of formal power series. It follows from (2.12) that the (left) action on $\mathcal{F}$ of any monomial in $D$ effects a shift of the coefficients of the power series according to the exponent vector of the monomial, which is the same action as the one defined by partial differentiation. Therefore, the $K$-vector space bases $\left(z^{\alpha} / \alpha!\mid \alpha \in\left(\mathbb{Z}_{\geq 0}\right)^{n}\right)$ and $\left(\partial^{\beta} \mid \beta \in\left(\mathbb{Z}_{\geq 0}\right)^{n}\right)$ are dual to each other with respect to the pairing (2.9), i.e.,

$$
\left(\partial^{\beta}, \sum_{\alpha \in(\mathbb{Z} \geq 0)^{n}} c_{\alpha} \frac{z^{\alpha}}{\alpha!}\right)=c_{\beta}, \quad \beta \in\left(\mathbb{Z}_{\geq 0}\right)^{n}, \quad c_{\alpha} \in K, \quad \alpha!:=\alpha_{1}!\cdot \ldots \cdot \alpha_{n}!
$$

Suppose that a system of (homogeneous) linear PDEs with constant coefficients for one unknown function of $n$ arguments is given. We compute a Janet basis $J$ for the ideal of $D$ which is generated by the left hand sides $p$ of these equations with respect to the term ordering $>$. The differential equations are considered as linear equations for $\left(\partial^{\beta}, f\right), \beta \in\left(\mathbb{Z}_{\geq 0}\right)^{n}$, where $f \in \mathcal{F}$ is a formal power series solution, and using the term ordering $>$, we may solve each of these equations for $(\operatorname{lm}(p), f)$.

Then Janet's algorithm partitions Mon $(D)$ into a set of monomials $m$ for which $(m, f) \in K$ can be chosen arbitrarily and a set $S$ of monomials for which $(\operatorname{lm}(p), f) \in K$ is uniquely determined by these choices. The latter set is the multiple-closed subset $S:=[\operatorname{lm}(p) \mid(p, \mu) \in J]$ of $\operatorname{Mon}(D)$. In particular, the $K$-dimension of the space of formal power series solutions, if finite, can be computed as the number of monomials in the complement $C$ of $S$ in $\operatorname{Mon}(D)$. In fact, the generalized Hilbert series $H_{C}\left(\partial_{1}, \ldots, \partial_{n}\right)$ of $C$ enumerates a basis for the Taylor coefficients $\left(\partial^{\beta}, f\right)$ of $f$ whose values can be assigned freely.
M. Janet calls the monomials $\partial^{\beta}$ in $\operatorname{Mon}(D) \backslash S$ parametric derivatives because the corresponding Taylor coefficients $\left(\partial^{\beta}, f\right)$ of a formal power series solution $f$ can be chosen arbitrarily. The monomials in $S$ are called principal derivatives [Jan29, e.g., no. 22, no. 38]. The Taylor coefficients $\left(\partial^{\beta}, f\right)$ which correspond to principal derivatives $\partial^{\beta}$ are uniquely determined by $K$-linear equations in terms of the Taylor coefficients of parametric derivatives. Of course, the extension of this method of determining the formal power series solutions of a system of linear partial differential equations is extended to the case of more than one unknown function (e.g., $q$ unknown functions) in a straightforward way by using submodules of $D^{1 \times q}$ instead of ideals of $D$.

Note that convergence of series solutions is to be investigated separately.
For a treatment of partial difference equations that is similar to Remark 2.41, we refer to [OP01]. Algorithm 2.38 is applicable in an analogous way to systems of linear partial difference equations, where the differential operators $\partial_{1}, \ldots, \partial_{n}$ are replaced by shift operators. For algorithmic details and applications we refer, e.g., to [GR06], [GR10], [GR12].

Remark 2.42. The previous remark also applies to linear systems of partial differential equations whose coefficients are rational functions in the independent variables $z_{1}, \ldots, z_{n}$, i.e. $D=$ $K\left[\partial_{1}, \ldots, \partial_{n}\right]$ is replaced by $B_{n}(K):=K\left(z_{1}, \ldots, z_{n}\right)\left\langle\partial_{1}, \ldots, \partial_{n}\right\rangle$. Of course, in this case a formal power series solution is only well-defined if the left submodule $M$ of $B_{n}(K)^{1 \times q}$ which represents the left hand sides of the equations is also a left submodule of $A\left[\partial_{1}, \ldots, \partial_{n}\right]^{1 \times q}$, where $A$ is a $K$-subalgebra of $B_{n}(K)$ whose elements do not have a pole in $0 \in K^{n}$ and the Janet basis for $M$ is computed within $A\left[\partial_{1}, \ldots, \partial_{n}\right]^{1 \times q}$. In other words, a formal power series solution is only welldefined if $0 \in K^{n}$ is not a zero of any denominator occurring in the course of Janet's algorithm.

Example 2.43. [Jan29, no. 23] Let us consider the heat equation

$$
\begin{equation*}
\frac{\partial u}{\partial t}-\frac{\partial^{2} u}{\partial x^{2}}=0 \tag{2.13}
\end{equation*}
$$

for an unknown real analytic function $u$ of $t$ and $x$. The corresponding operator is $p:=\partial_{t}-\partial_{x}^{2} \in$ $D:=K\left[\partial_{t}, \partial_{x}\right]$, where $K=\mathbb{Q}$ or $\mathbb{R}$. Choosing a degree-reverse lexicographical term ordering on $D$, the leading monomial of $p$ is $\partial_{x}^{2}$. The polynomial $p$ forms a Janet basis for the ideal of $D$ it generates, and the parametric derivatives are given by $\partial_{t}^{i}, \partial_{t}^{j} \partial_{x}, i, j \in \mathbb{Z}_{\geq 0}$. Hence, any choice of formal power series in $t$ for $u(t, 0)$ and $\frac{\partial u}{\partial x}(t, 0)$ uniquely determines a formal power series solution $u$ to (2.13). In this case, every choice of convergent power series yields a convergent series solution $u$. On the other hand, using the lexicographical term ordering extending $t>x$, the parametric derivatives are given by $\partial_{x}^{i}, i \in \mathbb{Z}_{\geq 0}$. Now, the choice $u(0, x)=\sum_{i \geq 0} x^{i}$ determines a divergent series solution $u$.

Example 2.44. The (minimal) Janet basis for the system of linear PDEs in Example 2.6 is

$$
\begin{aligned}
\frac{\partial^{2} u}{\partial x \partial y}=0, & \left\{*, \partial_{y}, \partial_{z}\right\} \\
\frac{\partial^{3} u}{\partial x^{2} \partial y}=0, & \left\{*, \partial_{y}, \partial_{z}\right\} \\
\frac{\partial^{4} u}{\partial x^{3} \partial z}=0, & \left\{\partial_{x}, *, \partial_{z}\right\} \\
\frac{\partial^{4} u}{\partial x^{3} \partial y}=0, & \left\{\partial_{x}, \partial_{y}, \partial_{z}\right\}
\end{aligned}
$$

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A Janet decomposition of the set of parametric derivatives is (cf. also Example 2.26)

$$
\begin{aligned}
1, & \left\{*, \partial_{y}, \partial_{z}\right\} \\
\partial_{x}, & \left\{*, *, \partial_{z}\right\} \\
\partial_{x}^{2}, & \left\{*, *, \partial_{z}\right\} \\
\partial_{x}^{3}, & \left\{\partial_{x}, *, *\right\}
\end{aligned}
$$

The corresponding generalized Hilbert series is

$$
\frac{1}{\left(1-\partial_{y}\right)\left(1-\partial_{z}\right)}+\frac{\partial_{x}}{1-\partial_{z}}+\frac{\partial_{x}^{2}}{1-\partial_{z}}+\frac{\partial_{x}^{3}}{1-\partial_{x}} .
$$

Accordingly, a formal power series solution $u$ of (2.5) is uniquely determined as

$$
u(x, y, z)=f_{0}(y, z)+x f_{1}(z)+x^{2} f_{2}(z)+x^{3} f_{3}(x)
$$

by any choice of formal power series $f_{0}, f_{1}, f_{2}, f_{3}$ of the indicated variables.

## 3. Systems of nonlinear differential equations

The methods to be developed in this section allow to solve tasks (a), (b), (c) as stated in the Introduction for systems of partial differential equations (PDEs) that are expressed by polynomials in the unknown functions and their derivatives.

A system of partial differential equations and inequations (or simply a differential system) $S$ is given by

$$
\begin{equation*}
p_{1}=0, \quad p_{2}=0, \quad \ldots, \quad p_{s}=0, \quad q_{1} \neq 0, \quad q_{2} \neq 0, \quad \ldots, \quad q_{t} \neq 0 \tag{3.1}
\end{equation*}
$$

where $p_{1}, \ldots, p_{s}$ and $q_{1}, \ldots, q_{t}$ are polynomials in unknown functions $u_{1}, \ldots, u_{m}$ of independent variables $z_{1}, \ldots, z_{n}$ and their partial derivatives (of arbitrary order), and $s, t \in \mathbb{Z}_{\geq 0}$.

Let $\Omega$ be an open and connected subset of $\mathbb{C}^{n}$ with coordinates $z_{1}, z_{2}, \ldots, z_{n}$. Then the solution set of $S$ on $\Omega$ is

$$
\begin{aligned}
\operatorname{Sol}_{\Omega}(S):=\left\{f=\left(f_{1}, \ldots, f_{m}\right) \quad \mid\right. & f_{k}: \Omega \rightarrow \mathbb{C} \text { analytic, } k=1, \ldots, m \\
& \left.p_{i}(f)=0, q_{j}(f) \neq 0, i=1, \ldots, s, j=1, \ldots, t\right\}
\end{aligned}
$$

where $p_{i}(f)$ and $q_{j}(f)$ are obtained from $p_{i}$ and $q_{j}$, respectively, by substituting $f_{k}$ for $u_{k}$ and the partial derivatives of $f_{k}$ for the corresponding jet variables in $u_{k}$.

Example 3.1. The following differential system for one unknown function $u$ of independent variables $t$ and $x$ is a combination of the Korteweg-de Vries equation (KdV, [BC80]) and a (generalized) Wronskian determinant:

$$
\left\{\begin{aligned}
\frac{\partial u}{\partial t}-6 u \frac{\partial u}{\partial x}+\frac{\partial^{3} u}{\partial x^{3}} & =0 \\
u \frac{\partial^{2} u}{\partial t \partial x}-\frac{\partial u}{\partial t} \frac{\partial u}{\partial x} & =0
\end{aligned}\right.
$$

If we denote partial derivatives by (repeated) indices, we may also write it as

$$
\left\{\begin{aligned}
u_{t}-6 u u_{x}+u_{x, x, x} & =0, \\
u u_{t, x}-u_{t} u_{x} & =0 .
\end{aligned}\right.
$$

Example 3.2. The following Navier-Stokes equations describe the flow of an incompressible fluid, where $x, y, z$ are the spatial coordinates, $t$ is the time coordinate, $\rho$ the constant density, $p=$ $p(x, y, z, t)$ the pressure, $\left(v_{1}, v_{2}, v_{3}\right)=\left(v_{1}(x, y, z, t), v_{2}(x, y, z, t), v_{3}(x, y, z, t)\right)$ the velocity vector,
$\left(g_{1}, g_{2}, g_{3}\right)$ the gravitational acceleration and $\mu$ is the dynamic viscosity of the fluid [LL87, p. 45]:

$$
\left\{\begin{aligned}
\rho\left(\frac{\partial v_{1}}{\partial t}+v_{1} \frac{\partial v_{1}}{\partial x}+v_{2} \frac{\partial v_{1}}{\partial y}+v_{3} \frac{\partial v_{1}}{\partial z}\right) & =-\frac{\partial p}{\partial x}+\mu\left(\frac{\partial^{2} v_{1}}{\partial x^{2}}+\frac{\partial^{2} v_{1}}{\partial y^{2}}+\frac{\partial^{2} v_{1}}{\partial z^{2}}\right)+\rho g_{1} \\
\rho\left(\frac{\partial v_{2}}{\partial t}+v_{1} \frac{\partial v_{2}}{\partial x}+v_{2} \frac{\partial v_{2}}{\partial y}+v_{3} \frac{\partial v_{2}}{\partial z}\right) & =-\frac{\partial p}{\partial y}+\mu\left(\frac{\partial^{2} v_{2}}{\partial x^{2}}+\frac{\partial^{2} v_{2}}{\partial y^{2}}+\frac{\partial^{2} v_{2}}{\partial z^{2}}\right)+\rho g_{2} \\
\rho\left(\frac{\partial v_{3}}{\partial t}+v_{1} \frac{\partial v_{3}}{\partial x}+v_{2} \frac{\partial v_{3}}{\partial y}+v_{3} \frac{\partial v_{3}}{\partial z}\right) & =-\frac{\partial p}{\partial z}+\mu\left(\frac{\partial^{2} v_{3}}{\partial x^{2}}+\frac{\partial^{2} v_{3}}{\partial y^{2}}+\frac{\partial^{2} v_{3}}{\partial z^{2}}\right)+\rho g_{3} \\
\frac{\partial v_{1}}{\partial x}+\frac{\partial v_{2}}{\partial y}+\frac{\partial v_{3}}{\partial z} & =0
\end{aligned}\right.
$$

The consequences of (3.1) are the partial differential equations for $u_{1}, \ldots, u_{m}$ which are obtained in a finite number steps from the following rules:
(a) The given equations $p_{1}=0, p_{2}=0, \ldots, p_{s}=0$ are consequences of (3.1).
(b) If $p=0$ is a consequence, then any partial derivative of $p=0$ is a consequence.
(c) If $p \cdot q=0$ is a consequence and $q$ is a factor of some $q_{i}$, then $p=0$ is a consequence.
(d) If $p=0$ and $r=0$ are consequences, then $a p+b r=0$ are consequences for all polynomials $a$ and $b$ in $u_{1}, \ldots, u_{m}$ and their partial derivatives (of all orders).
Since this setup allows differential equations $p=0$ to be differentiated, we are going to work with a polynomial ring in $u_{1}, \ldots, u_{m}$ which admits these differentiations.

Definition 3.3. A differential ring $R$ with commuting derivations $\delta_{1}, \ldots, \delta_{n}$ is a commutative ring $R$ endowed with maps $\delta_{i}: R \rightarrow R$, satisfying

$$
\delta_{i}\left(r_{1}+r_{2}\right)=\delta_{i}\left(r_{1}\right)+\delta_{i}\left(r_{2}\right), \quad \delta_{i}\left(r_{1} r_{2}\right)=\delta_{i}\left(r_{1}\right) r_{2}+r_{1} \delta_{i}\left(r_{2}\right) \quad \text { for all } r_{1}, r_{2} \in R
$$

$i=1, \ldots, n$, and $\delta_{i} \circ \delta_{j}=\delta_{j} \circ \delta_{i}$ for all $1 \leq i, j \leq n$. A differential ring which is a field is called a differential field, and similarly for a differential algebra over a differential field.

In what follows we only consider differential fields $K$ of characteristic zero. Let $\partial_{1}, \ldots, \partial_{n}$ be the derivations of $K$.

Definition 3.4. The differential polynomial ring $K\left\{u_{1}, \ldots, u_{m}\right\}$ in the differential indeterminates $u_{1}, \ldots, u_{m}$ is the commutative polynomial algebra $K\left[\left(u_{k}\right)_{J} \mid 1 \leq k \leq m, J \in\left(\mathbb{Z}_{\geq 0}\right)^{n}\right]$ with infinitely many, algebraically independent indeterminates $\left(u_{k}\right)_{J}$, also called jet variables, which represent the partial derivatives

$$
\frac{\partial^{J_{1}+\ldots+J_{n}} U_{k}}{\partial z_{1}^{J_{1}} \ldots \partial z_{n}^{J_{n}}}, \quad k=1, \ldots, m, \quad J \in\left(\mathbb{Z}_{\geq 0}\right)^{n}
$$

of smooth functions $U_{1}, \ldots, U_{m}$ of independent variables $z_{1}, \ldots, z_{n}$. We identify $u_{k}$ with $\left(u_{k}\right)_{(0, \ldots, 0)}$ and we also employ the notation involving $z_{1}, \ldots, z_{n}$ as repeated indices as in Example 3.1. The ring $K\left\{u_{1}, \ldots, u_{m}\right\}$ is considered as a differential ring with commuting derivations $\delta_{1}, \ldots, \delta_{n}$ defined by extending

$$
\delta_{i}\left(u_{k}\right)_{J}:=\left(u_{k}\right)_{J+1_{i}}, \quad i=1, \ldots, n, \quad k=1, \ldots, m, \quad J \in\left(\mathbb{Z}_{\geq 0}\right)^{n}
$$

additively, respecting the product rule of differentiation, and restricting to the derivation $\partial_{i}$ on $K$. Here $1_{i}$ denotes the multi-index $(0, \ldots, 0,1,0, \ldots, 0)$ of length $n$ with 1 at position $i$. More generally, the differential polynomial ring may be constructed with coefficients in a differential ring rather than in a differential field in the same way.

Recall that we consider an open and connected subset $\Omega$ of $\mathbb{C}^{n}$. The set of (complex) meromorphic functions on $\Omega$ form a field $K$, and together with the partial differential operators with respect to $z_{1}, \ldots, z_{n}$ the field $K$ is a differential field.

A suitable choice of differential polynomial ring $R=K\left\{u_{1}, \ldots, u_{m}\right\}$ allows to consider the left hand sides $p_{1}, \ldots, p_{s}, q_{1}, \ldots, q_{t}$ in the system of nonlinear PDEs (3.1) as elements of $R$. Moreover, the left hand sides of all consequences of the system are elements of $R$ as well. In fact, we may consider the differential ideal $I$ of $R$ which is generated by $p_{1}, \ldots, p_{s}$, i.e., the smallest ideal of $R$ which contains $p_{1}, \ldots, p_{s}$ and all their derivatives (of all orders). This is only a first step, because
in general $I$ does not contain all consequences of (3.1). Before we study these ideas further, we first deal with algebraic systems (e.g., in $\left(u_{k}\right)_{J}$, where $\left.1 \leq k \leq m, J \in\left(\mathbb{Z}_{\geq 0}\right)^{n}\right)$.
3.1. Thomas decomposition of algebraic systems. In this subsection let $K$ be a field of characteristic zero and $R=K\left[x_{1}, \ldots, x_{n}\right]$ the commutative polynomial algebra with indeterminates $x_{1}, \ldots, x_{n}$ over $K$. We denote by $\bar{K}$ an algebraic closure of $K$.

Definition 3.5. An algebraic system $S$, defined over $R$, is given by finitely many equations and inequations

$$
\begin{equation*}
p_{1}=0, \quad p_{2}=0, \quad \ldots, \quad p_{s}=0, \quad q_{1} \neq 0, \quad q_{2} \neq 0, \quad \ldots, \quad q_{t} \neq 0 \tag{3.2}
\end{equation*}
$$

where $p_{1}, \ldots, p_{s}, q_{1}, \ldots, q_{t} \in R$ and $s, t \in \mathbb{Z}_{\geq 0}$. The solution set of $S$ in $\bar{K}^{n}$ is

$$
\operatorname{Sol}_{\bar{K}}(S):=\left\{a \in \bar{K}^{n} \mid p_{i}(a)=0 \text { and } q_{j}(a) \neq 0 \text { for all } 1 \leq i \leq s, 1 \leq j \leq t\right\}
$$

We fix a total ordering $>$ on the set $\left\{x_{1}, \ldots, x_{n}\right\}$ allowing us to consider every non-constant element $p$ of $R$ as a univariate polynomial in the greatest variable with respect to $>$ which occurs in $p$, with coefficients which are themselves univariate polynomials in lower ranked variables, etc. Without loss of generality we may assume that

$$
x_{1}>x_{2}>\ldots>x_{n}
$$

The choice of $>$ corresponds to a choice of projections

$$
\begin{aligned}
& \pi_{1}: \bar{K}^{n} \longrightarrow \bar{K}^{n-1}: \quad\left(a_{1}, a_{2}, \ldots, a_{n}\right) \longmapsto\left(a_{2}, a_{3}, a_{4}, \ldots, a_{n}\right), \\
& \pi_{2}: \bar{K}^{n} \longrightarrow \bar{K}^{n-2}: \quad\left(a_{1}, a_{2}, \ldots, a_{n}\right) \longmapsto \quad\left(a_{3}, a_{4}, \ldots, a_{n}\right), \\
& \pi_{n-1}: \bar{K}^{n} \longrightarrow \bar{K}: \quad\left(a_{1}, a_{2}, \ldots, a_{n}\right) \longmapsto \quad a_{n} .
\end{aligned}
$$

According to this choice, the recursive representation of polynomials is motivated by considering the $(k-1)$-st projection $\pi_{k-1}\left(\operatorname{Sol}_{\bar{K}}(S)\right)$ of the solution set as fibered over the $k$-th projection $\pi_{k}\left(\operatorname{Sol}_{\bar{K}}(S)\right)$, for $k=1, \ldots, n-1$, where we define $\pi_{0}:=\operatorname{id}_{\bar{K}^{n}}$ (cf. also [Ple09a]). The purpose of a Thomas decomposition of $\operatorname{Sol}_{\bar{K}}(S)$, to be defined below, is to clarify this fibration structure. The solution set $\operatorname{Sol}_{\bar{K}}(S)$ is partitioned into subsets $\operatorname{Sol}_{\bar{K}}\left(S_{1}\right), \ldots, \operatorname{Sol}_{\bar{K}}\left(S_{r}\right)$ in such a way that, for each $i=1, \ldots, r$ and $k=1, \ldots, n-1$, the fiber cardinality $\left|\pi_{k}^{-1}(\{a\})\right|$ does not depend on the choice of $a \in \pi_{k}\left(\operatorname{Sol}_{\bar{K}}\left(S_{i}\right)\right)$. In terms of the defining equations and inequations in (3.2), the fundamental obstructions to this uniform behavior are zeros of the leading coefficients of $p_{i}$ or $q_{j}$ and zeros of $p_{i}$ or $q_{j}$ of multiplicity greater than one.

Definition 3.6. Let $p \in R \backslash K$.
(a) The greatest indeterminate with respect to $>$ which occurs in $p$ is referred to as the leader of $p$ and is denoted by $\operatorname{ld}(p)$.
(b) For $v=\operatorname{ld}(p)$ we denote by $\operatorname{deg}_{v}(p)$ the degree of $p$ in $v$.
(c) The coefficient of the highest power of $\operatorname{ld}(p)$ occurring in $p$ is called the initial of $p$ and is denoted by $\operatorname{init}(p)$.
(d) The discriminant of $p$ is defined as

$$
\operatorname{disc}(p):=(-1)^{d(d-1) / 2} \operatorname{res}\left(p, \frac{\partial p}{\partial \operatorname{ld}(p)}, \operatorname{ld}(p)\right) / \operatorname{init}(p), \quad d=\operatorname{deg}_{\operatorname{ld}(p)}(p)
$$

where $\operatorname{res}(p, q, v)$ is the resultant of $p$ and $q$ with respect to the variable $v$.
(Note that $\operatorname{disc}(p)$ is a polynomial because $\operatorname{init}(p) \operatorname{divides} \operatorname{res}(p, \partial p / \partial \operatorname{ld}(p), \operatorname{ld}(p))$ : for $p=$

$$
\begin{aligned}
& c_{d} v^{d}+c_{d-1} v^{d-1}+\ldots+c_{1} v+c_{0} \text { the }(2 d-1) \times(2 d-1) \text { matrix in } \\
& \operatorname{res}(p, \partial p / \partial \operatorname{ld}(p), \operatorname{ld}(p)):=\operatorname{det}\left(\begin{array}{cccccccc}
c_{d} & \ldots & \ldots & c_{1} & c_{0} & & & \\
& c_{d} & \ldots & \ldots & c_{1} & c_{0} & & \\
& & \ddots & & & & \ddots & \\
& & & c_{d} & \ldots & \ldots & c_{1} & c_{0} \\
d c_{d} & \ldots & 2 c_{2} & c_{1} & & & & \\
& d c_{d} & \ldots & 2 c_{2} & c_{1} & & & \\
& & \ddots & & & \ddots & & \\
& & & d c_{d} & \ldots & 2 c_{2} & c_{1} & \\
& & & & d c_{d} & \ldots & 2 c_{2} & c_{1}
\end{array}\right)
\end{aligned}
$$

has a column all of whose entries are divisible by $c_{d}=\operatorname{init}(p)$.)
$\operatorname{Both} \operatorname{init}(p)$ and $\operatorname{disc}(p)$ are elements of the polynomial algebra $K[x \mid \operatorname{ld}(p)>x]$. The zeros of a univariate polynomial with multiplicity greater than one are the common zeros of the polynomial and its derivative. The solutions of $\operatorname{disc}(p)=0$ in $\bar{K}^{n-k}$, where $\operatorname{ld}(p)=x_{k}$, are therefore those tuples $\left(a_{k+1}, a_{k+2}, \ldots, a_{n}\right)$ for which the substitution

$$
x_{k+1}=a_{k+1}, \quad x_{k+2}=a_{k+2}, \quad \ldots, \quad x_{n}=a_{n}
$$

in $p$ results in a univariate polynomial with a zero of multiplicity greater than one.
Definition 3.7. An algebraic system $S$, defined over $R$, as in (3.2) is said to be simple (with respect to $>$ ) if the following three conditions hold.
(a) For all $i=1, \ldots, s$ and $j=1, \ldots, t$ we have $p_{i} \notin K$ and $q_{j} \notin K$.
(b) The leaders of the left hand sides of the equations and inequations in $S$ are pairwise distinct, i.e., $\left|\left\{\operatorname{ld}\left(p_{1}\right), \ldots, \operatorname{ld}\left(p_{s}\right), \operatorname{ld}\left(q_{1}\right), \ldots, \operatorname{ld}\left(q_{t}\right)\right\}\right|=s+t$.
(c) For every $r \in\left\{p_{1}, \ldots, p_{s}, q_{1}, \ldots, q_{t}\right\}$, if $\operatorname{ld}(r)=x_{k}$, then neither the equation $\operatorname{init}(r)=0$ nor the equation $\operatorname{disc}(r)=0$ has a solution $\left(a_{k+1}, a_{k+2}, \ldots, a_{n}\right)$ in $\pi_{k}\left(\operatorname{Sol}_{\bar{K}}(S)\right)$.
Subsets of non-constant polynomials in $R$ with pairwise distinct leaders (i.e., satisfying (a) and (b)) are also referred to as triangular sets (cf., e.g., [ALMM99], [Hub03a, Hub03b], [Wan01]).

Remark 3.8. A simple algebraic system $S$ admits the following solution procedure, which also shows that its solution set is not empty. Let $S_{<k}$ be the subset of $S$ consisting of the equations $p=0$ and inequations $q \neq 0$ with $x_{k}>\operatorname{ld}(p)$ and $x_{k}>\operatorname{ld}(q)$. The fibration structure implied by (c) ensures that, for every $k=1, \ldots, n-1$, every solution

$$
\left(a_{k+1}, a_{k+2}, \ldots, a_{n}\right) \in \bar{K}^{n-k}
$$

in $\pi_{k}\left(\operatorname{Sol}_{\bar{K}}(S)\right)=\pi_{k}\left(\operatorname{Sol}_{\bar{K}}\left(S_{<k}\right)\right)$ can be extended to a solution

$$
\left(a_{k}, a_{k+1}, \ldots, a_{n}\right) \in \bar{K}^{n-(k-1)}
$$

in $\pi_{k-1}\left(\operatorname{Sol}_{\bar{K}}(S)\right)$. If $S$ contains an equation $p=0$ with leader $x_{k}$, then there exist exactly $\operatorname{deg}_{x_{k}}(p)$ such extensions (because zeros with multiplicity greater than one are excluded by the non-vanishing discriminant). If $S$ contains an inequation $q \neq 0$ with leader $x_{k}$, all $a_{k} \in \bar{K}$ except $\operatorname{deg}_{x_{k}}(q)$ elements define a tuple $\left(a_{k}, a_{k+1}, \ldots, a_{n}\right)$ as above. If no equation and no inequation in $S$ has leader $x_{k}$, then $a_{k} \in \bar{K}$ can be chosen arbitrarily.

Definition 3.9. Let $S$ be an algebraic system, defined over $R$. A Thomas decomposition of $S$ (or $\mathrm{Sol}_{\bar{K}}(S)$ ) with respect to $>$ is a collection of finitely many algebraic systems $S_{1}, \ldots, S_{r}$, each of which is defined over $R$ and is simple, such that $\operatorname{Sol} \bar{K}(S)$ is the disjoint union of the solution sets $\operatorname{Sol}_{\bar{K}}\left(S_{1}\right), \ldots, \operatorname{Sol}_{\bar{K}}\left(S_{r}\right)$.

We outline a method for computing a Thomas decomposition of algebraic systems.
Remark 3.10. Given $S$ as in (3.2) and a total ordering $>$ on $\left\{x_{1}, \ldots, x_{n}\right\}$, a Thomas decomposition of $S$ with respect to $>$ can be constructed by combining Euclid's algorithm with a splitting strategy.

First of all, if $S$ contains an equation $c=0$ with $0 \neq c \in K$ or the inequation $0 \neq 0$, then $S$ is discarded because it has no solutions. Moreover, from now on the equation $0=0$ and inequations $c \neq 0$ with $0 \neq c \in K$ are supposed to be removed from $S$.

An elementary step of the algorithm applies a pseudo-division to a pair $p_{1}, p_{2}$ of non-constant polynomials in $R$ with the same leader $x_{k}$ and $\operatorname{deg}_{x_{k}}\left(p_{1}\right) \geq \operatorname{deg}_{x_{k}}\left(p_{2}\right)$. The result is a pseudo-remainder

$$
\begin{equation*}
r=c_{1} \cdot p_{1}-c_{2} \cdot p_{2}, \quad \text { where } c_{1}, c_{2} \in R \tag{3.3}
\end{equation*}
$$

and $r$ is constant or has leader less than $x_{k}$ or has leader $x_{k}$ and $\operatorname{deg}_{x_{k}}(r)<\operatorname{deg}_{x_{k}}\left(p_{1}\right)$. Since the coefficients of $p_{1}$ and $p_{2}$ are polynomials in lower ranked variables, multiplication of $p_{1}$ by a non-constant polynomial $c_{1}$ may be necessary in general to perform the reduction in $R$ (and not in its field of fractions). Choosing $c_{1}$ as a suitable power of $\operatorname{init}\left(p_{2}\right)$ always achieves this.

In order to turn $S$ into a triangular set, the algorithm deals with three kinds of subsets of $S$ of cardinality two. Firstly, each pair of equations $p_{1}=0, p_{2}=0$ in $S$ with $\operatorname{ld}\left(p_{1}\right)=\operatorname{ld}\left(p_{2}\right)$ is replaced with the single equation $r=0$, where $r$ is the result of applying Euclid's algorithm to $p_{1}$ and $p_{2}$, considered as univariate polynomials in their leader, using the above pseudo-division. (If this computation was stable under substitution of values for lower ranked variables in $p_{1}$ and $p_{2}$, then $r$ would be the greatest common divisor of the specialized polynomials.)

The solution set of the system is supposed not to change, when the equation $p_{1}=0$ is replaced with the equation $r=0$ given by the pseudo-reduction (3.3). Therefore, we assume that the polynomial $c_{1}$, and hence $\operatorname{init}\left(p_{2}\right)$, does not vanish on the solution set of the system. In order to ensure this condition, a preparatory step splits the system into two, if necessary, and adds the inequation $\operatorname{init}\left(p_{2}\right) \neq 0$ to one of them and the equation $\operatorname{init}\left(p_{2}\right)=0$ to the other. The algorithm then deals with both systems separately. These case distinctions also allow to arrange for the part of condition (c) in Definition 3.7 which concerns initials.

Secondly, let $p=0, q \neq 0$ be in $S$ with $\operatorname{ld}(p)=\operatorname{ld}(q)=x_{k}$. If $\operatorname{deg}_{x_{k}}(p) \leq \operatorname{deg}_{x_{k}}(q)$, then $q \neq 0$ is replaced with $r \neq 0$, where $r$ is the result of applying the pseudo-division (3.3) to $q$ and $p$. Otherwise, Euclid's algorithm is applied to $p$ and $q$, keeping track of the coefficients used for the reductions as in (3.3). Given the result $r$, the system is then split into two, adding the conditions $r \neq 0$ and $r=0$, respectively. The inequation $q \neq 0$ is removed from the first new system, because $p=0$ and $q \neq 0$ have no common solution in that case. The assumption $r=0$ and the bookkeeping allow to divide $p$ by the common factor of $p$ and $q$ (modulo left hand sides of equations with smaller leader). The left hand side of $p=0$ is replaced with that quotient in the second new system. Some particular cases admit an accelerated treatment. For instance, if $p$ divides $q$, then the solution set of $S$ is empty and $S$ is discarded.

Thirdly, for a pair $q_{1} \neq 0, q_{2} \neq 0$ in $S$ with $\operatorname{ld}\left(q_{1}\right)=\operatorname{ld}\left(q_{2}\right)$, Euclid's algorithm is applied to $q_{1}$ and $q_{2}$ in the same way as above. Keeping track of the coefficients used in intermediate steps allows to determine the least common multiple $m$ of $q_{1}$ and $q_{2}$, which again depends on distinguishing the cases whether the result of Euclid's algorithm vanishes or not. The pair $q_{1} \neq 0, q_{2} \neq 0$ is then replaced with the single inequation $m \neq 0$.

The part of condition (c) in Definition 3.7 regarding discriminants is taken care of by applying Euclid's algorithm as above to $p$ and $\partial p / \partial \operatorname{ld}(p)$, where $p$ is the left hand side of an equation or inequation. Bookkeeping allows to determine the square-free part of $p$, which depends again on case distinctions.

Expressions tend to grow very quickly when performing these reductions, so that an appropriate strategy is essential for dealing with non-trivial systems. Apart from dividing by the content (in $K$ ) of polynomials, in intermediate steps of Euclid's algorithm the coefficients should be reduced modulo equations in the system with lower ranked leaders. In practice, subresultant computations (cf., e.g., [Mis93]) allow to diminish the growth of coefficients significantly.

Termination of the procedure sketched above depends on the organization of its steps. One possible strategy is to maintain an intermediate triangular set, reduce new equations and inequations modulo the equations in the triangular set, and select among these results the one with smallest leader and least degree, preferably an equation, for insertion into the triangular set. If the set already contains an equation or inequation with the same leader, then the pair is treated as discussed above. Since equations are replaced with equations of smaller degree and inequations are


Figure 3.1: Stereographic projection from a circle
replaced with equations if possible or with the least common multiple of inequations, this strategy terminates after finitely many steps.

For more details on the algebraic part of Thomas' algorithm, we refer to [BGL+12], [Bäc14], and [Rob14, Subsect. 2.2.1].

An implementation of Thomas' algorithm for algebraic systems was developed by T. Bächler at RWTH Aachen University as Maple package AlgebraicThomas [BLH].

In what follows, variables are underlined to emphasize that they are leaders of polynomials with respect to the fixed total ordering $>$.

Example 3.11. Let us compute a Thomas decomposition of the algebraic system

$$
x^{2}+y^{2}-1=0
$$

consisting of one equation, defined over $R=\mathbb{Q}[x, y]$, with respect to $x>y$. First we set $p_{1}:=$ $x^{2}+y^{2}-1$. Then we have $\operatorname{ld}\left(p_{1}\right)=x$ and $\operatorname{init}\left(p_{1}\right)=1$ and

$$
\operatorname{disc}\left(p_{1}\right)=-4 y^{2}+4
$$

We distinguish the cases whether or not $p_{1}=0$ has a solution which is also a zero of $\operatorname{disc}\left(p_{1}\right)$, or equivalently, of $y^{2}-1$. In other words, we replace the original algebraic system with two algebraic systems which are obtained by adding the inequation $y^{2}-1 \neq 0$ or the equation $y^{2}-1=0$. The first system is readily seen to be simple, whereas the second one is transformed into a simple system by taking the difference of the two equations and computing a square-free part. Clearly, the solution sets of the two resulting simple systems form a partition of the solution set of $p_{1}=0$. We obtain the Thomas decomposition

$$
\begin{aligned}
\underline{x}^{2}+y^{2}-1 & =0 \\
\underline{y}^{2}-1 & \neq 0
\end{aligned}
$$

$$
\begin{aligned}
& \underline{x}=0 \\
& \underline{y}^{2}-1=0 \\
& \hline
\end{aligned}
$$

In this example, all points of $\operatorname{Sol}_{\bar{K}}\left(\left\{p_{1}=0\right\}\right)$ for which the projection $\pi_{1}$ onto the $y$-axis has fibers of an exceptional cardinality have real coordinates, and the significance of the above case distinction can be confirmed graphically.

As a further illustration let us augment the original system by the equation which expresses the coordinate $t$ of the point of intersection of the line through the two points $(0,1)$ and $(x, y)$ on the circle with the $x$-axis (stereographic projection, cf. Figure 3.1):

$$
\left\{\begin{array}{l}
x^{2}+y^{2}-1=0 \\
(1-y) t-x=0
\end{array}\right.
$$

A Thomas decomposition with respect to the ordering $x>y>t$ is obtained as follows. We set $p_{2}:=x+t y-t$. Since $\operatorname{ld}\left(p_{1}\right)=\operatorname{ld}\left(p_{2}\right)$, we apply polynomial division:

$$
p_{1}-(x-t y+t) p_{2}=\left(1+t^{2}\right) \underline{y}^{2}-2 t^{2} \underline{y}+t^{2}-1=(\underline{y}-1)\left(\left(1+t^{2}\right) \underline{y}-t^{2}+1\right) .
$$

Replacing $p_{1}$ with the remainder of this division does not alter the solution set of the algebraic system. It is convenient (but not necessary) to split the system into two systems according to the
factorization of the remainder:

$$
\left\{\begin{array} { r l } 
{ \underline { x } + t y - t } & { = 0 } \\
{ ( 1 + t ^ { 2 } ) \underline { y } - t ^ { 2 } + 1 } & { = 0 } \\
{ \underline { y } - 1 } & { \neq 0 }
\end{array} \quad \left\{\begin{array}{r}
\underline{x}+t y-t=0 \\
\underline{y}-1=0
\end{array}\right.\right.
$$

Another polynomial division reveals that the equation and the inequation with leader $y$ in the first system have no common solutions. Hence, the inequation can be omitted from that system. For the investigation of the initial of the equation we note that the assumption $1+t^{2}=0$ leads to a contradiction. Finally, the equation with leader $y$ can be used to eliminate $y$ in the first equation:

$$
\left(1+t^{2}\right)(x+t y-t)-t\left(\left(1+t^{2}\right) y-t^{2}+1\right)=\left(1+t^{2}\right) \underline{x}-2 t
$$

A similar simplification can be applied to the second system. We obtain the Thomas decomposition

$$
\begin{align*}
\left(1+t^{2}\right) \underline{x}-2 t & =0 \\
\left(1+t^{2}\right) \underline{y}-t^{2}+1 & =0 \\
\underline{t}^{2}+1 & \neq 0
\end{align*}
$$

$$
\begin{array}{r}
\underline{x}=0 \\
\underline{y}-1=0
\end{array}
$$

from which a rational parametrization of the circle can be read off.
Remark 3.12. A Thomas decomposition of an algebraic system is not uniquely determined. It depends on the chosen total ordering $>$, the order in which intermediate systems are dealt with and other choices, such as whether factorizations of left hand sides of equations are taken into account or not.

The solution sets $V$ in $\bar{K}^{n}$ of systems of polynomial equations in $x_{1}, \ldots, x_{n}$, defined over $R$, are in one-to-one correspondence with their vanishing ideals in $R$

$$
\mathcal{I}_{R}(V):=\{p \in R \mid p(a)=0 \text { for all } a \in V\}
$$

and these are the radical ideals of $R$, i.e., the ideals $I$ of $R$ which equal their radicals

$$
\sqrt{I}:=\left\{p \in R \mid p^{r} \in I \text { for some } r \in \mathbb{Z}_{\geq 0}\right\}
$$

(Hilbert's Nullstellensatz; cf., e.g., [Eis95]). The solution sets $V$ can then be considered as closed subsets of $\bar{K}^{n}$ with respect to the Zariski topology.

The fibration structure of a simple algebraic system $S$ allows to deduce that the polynomials in $R$ which vanish on $\operatorname{Sol}_{\bar{K}}(S)$ are precisely those polynomials in $R$ whose pseudo-remainders modulo $p_{1}, \ldots, p_{s}$ are zero, where $p_{1}=0, \ldots, p_{s}=0$ are the equations in $S$. If $E$ is the ideal of $R$ generated by $p_{1}, \ldots, p_{s}$ and $q$ the product of all $\operatorname{init}\left(p_{i}\right)$, then these polynomials form the saturation ideal

$$
E: q^{\infty}:=\left\{p \in R \mid q^{r} \cdot p \in E \text { for some } r \in \mathbb{Z}_{\geq 0}\right\}
$$

In particular, simple algebraic systems admit an effective way to decide membership of a polynomial to the associated radical ideal (cf. also Proposition 3.32 below).
Proposition 3.13 ([Rob14], Prop. 2.2.7). Let the algebraic system $S$ given by

$$
p_{1}=0, \quad p_{2}=0, \quad \ldots, \quad p_{s}=0, \quad q_{1} \neq 0, \quad q_{2} \neq 0, \quad \ldots, \quad q_{t} \neq 0
$$

be simple. Moreover, let $E$ be the ideal of $R$ generated by $p_{1}, p_{2}, \ldots, p_{s}$, and $q$ the product of all $\operatorname{init}\left(p_{i}\right)$. Then $E: q^{\infty}$ consists of all polynomials in $R$ which vanish on $\operatorname{Sol}_{\bar{K}}(S)$. In particular, $E: q^{\infty}$ is a radical ideal. Given $p \in R$, we have $p \in E: q^{\infty}$ if and only if the pseudo-remainder of $p$ modulo $p_{1}, \ldots, p_{s}$ is zero.

Example 3.14. Continuing Example 3.11, let $E$ be the ideal of $R$ generated by the left hand sides of the equations of the simple algebraic system

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

and define $q=1+t^{2}$. Moreover, let $p=\left(1-t^{2}\right) \underline{x}+2 t y \in R$. The pseudo-remainder of $p$ modulo the equations of the first simple algebraic system displayed at the end of Example 3.11 is computed as follows. First we have

$$
p^{\prime}:=\left(1+t^{2}\right) p-\left(1-t^{2}\right)\left[\left(1+t^{2}\right) x-2 t\right]=2\left(1+t^{2}\right) t \underline{y}+2\left(1-t^{2}\right) t
$$

Then we have

$$
r:=p^{\prime}-2 t\left[\left(1+t^{2}\right) y-t^{2}+1\right]=0
$$

Since the pseudo-remainder $r$ is zero, we conclude that $p \in E: q^{\infty}$.
3.2. Thomas decomposition of differential systems. Let $K$ be the differential field of meromorphic functions on an open and connected subset $\Omega$ of $\mathbb{C}^{n}$ with coordinates $z_{1}, \ldots, z_{n}$. We define the differential polynomial ring $R=K\left\{u_{1}, \ldots, u_{m}\right\}$ with commuting derivations $\partial_{1}, \ldots, \partial_{n}$ and we set $\Delta:=\left\{\partial_{1}, \ldots, \partial_{n}\right\}$.
Definition 3.15. A differential system $S$, defined over $R=K\left\{u_{1}, \ldots, u_{m}\right\}$, is given by finitely many equations and inequations

$$
\begin{equation*}
p_{1}=0, \quad p_{2}=0, \quad \ldots, \quad p_{s}=0, \quad q_{1} \neq 0, \quad q_{2} \neq 0, \quad \ldots, \quad q_{t} \neq 0 \tag{3.4}
\end{equation*}
$$

where $p_{1}, \ldots, p_{s}, q_{1}, \ldots, q_{t} \in R$ and $s, t \in \mathbb{Z}_{\geq 0}$. The solution set of $S$ is

$$
\begin{aligned}
\operatorname{Sol}_{\Omega}(S):=\left\{f=\left(f_{1}, \ldots, f_{m}\right) \quad \mid\right. & f_{k}: \Omega \rightarrow \mathbb{C} \text { analytic, } k=1, \ldots, m \\
& \left.p_{i}(f)=0, q_{j}(f) \neq 0, i=1, \ldots, s, j=1, \ldots, t\right\}
\end{aligned}
$$

Remark 3.16. Since each component $f_{k}$ of a solution of (3.4) is assumed to be analytic, the equations $p_{i}=0$ and inequations $q_{j} \neq 0$ (and their consequences) can be translated into algebraic conditions on the Taylor coefficients of power series expansions of $f_{1}, \ldots, f_{m}$ (around a point in $\Omega$ ). An inequation $q \neq 0$ then turns into a disjunction of algebraic inequations for all coefficients which result from substitution of power series expansions for $u_{1}, \ldots, u_{m}$ in $q$.

An appropriate choice of $\Omega \subseteq \mathbb{C}^{n}$ can often only be made after the formal treatment of a given differential system by the formal methods discussed in these notes (as, e.g., singularities of coefficients in differential consequences will only be detected during that process). In general, we assume that $\Omega$ is chosen in such a way that the given systems have analytic solutions on $\Omega$.

Clearly, by neglecting the derivations on $R=K\left\{u_{1}, \ldots, u_{m}\right\}$, a differential system can be considered as an algebraic system in the finitely many variables $\left(u_{k}\right)_{J}$ which occur in the equations and inequations. The same recursive representation of polynomials as in the algebraic case is employed, but the total ordering on the set of variables $\left(u_{k}\right)_{J}$ is supposed to respect the action of the derivations. Then the methods of the previous section on algebraic systems are applicable.

Definition 3.17. A ranking $>$ on $R=K\left\{u_{1}, \ldots, u_{m}\right\}$ is a total ordering on the set

$$
\operatorname{Mon}(\Delta) u:=\left\{\left(u_{k}\right)_{J} \mid 1 \leq k \leq m, J \in\left(\mathbb{Z}_{\geq 0}\right)^{n}\right\}
$$

such that for all $j \in\{1, \ldots, n\}, k, k_{1}, k_{2} \in\{1, \ldots, m\}, J_{1}, J_{2} \in\left(\mathbb{Z}_{\geq 0}\right)^{n}$ we have
(a) $\partial_{j} u_{k}>u_{k}$ and
(b) $\left(u_{k_{1}}\right)_{J_{1}}>\left(u_{k_{2}}\right)_{J_{2}}$ implies $\partial_{j}\left(u_{k_{1}}\right)_{J_{1}}>\partial_{j}\left(u_{k_{2}}\right)_{J_{2}}$.

A ranking $>$ is said to be orderly if

$$
\left|J_{1}\right|>\left|J_{2}\right| \quad \Longrightarrow \quad\left(u_{k_{1}}\right)_{J_{1}}>\left(u_{k_{2}}\right)_{J_{2}} \quad \text { for all } 1 \leq k_{1}, k_{2} \leq m, J_{1}, J_{2} \in\left(\mathbb{Z}_{\geq 0}\right)^{n} .
$$

Remark 3.18. Every ranking $>$ on $R$ is a well-ordering (cf., e.g., [Kol73, Ch. 0, Sect. 17, Lemma 15]), i.e., every descending sequence of elements of $\operatorname{Mon}(\Delta) u$ terminates.
Example 3.19. On the differential polynomial ring $K\{u\}$ (i.e., where $m=1$ ) with commuting derivations $\partial_{1}, \ldots, \partial_{n}$ the degree-reverse lexicographical ranking (with $\partial_{1} u>\partial_{2} u>\ldots>\partial_{n} u$ ) is defined for $u_{J}, u_{J^{\prime}}, J=\left(j_{1}, \ldots, j_{n}\right), J^{\prime}=\left(j_{1}^{\prime}, \ldots, j_{n}^{\prime}\right) \in\left(\mathbb{Z}_{\geq 0}\right)^{n}$, by

$$
u_{J}>u_{J^{\prime}} \quad \Longleftrightarrow\left\{\begin{array}{l}
j_{1}+\ldots+j_{n}>j_{1}^{\prime}+\ldots+j_{n}^{\prime} \quad \text { or } \\
\left(j_{1}+\ldots+j_{n}=j_{1}^{\prime}+\ldots+j_{n}^{\prime} \quad \text { and } \quad J \neq J^{\prime} \quad\right. \text { and } \\
\left.j_{i}<j_{i}^{\prime} \text { for } i=\max \left\{1 \leq k \leq n \mid j_{k} \neq j_{k}^{\prime}\right\}\right)
\end{array}\right.
$$

For instance, if $n=3$, we have $u_{(1,2,1)}>u_{(1,2,0)}>u_{(2,0,1)}$.

In what follows, we assume that a ranking $>$ on $R=K\left\{u_{1}, \ldots, u_{m}\right\}$ is fixed. For each $p \in R \backslash K$, the leader $\operatorname{ld}(p)$ and the initial $\operatorname{init}(p)$ are defined as in the previous section on algebraic systems. With the aim of introducing simple differential systems (Definition 3.24) we discuss pseudo-division for differential polynomials first.

Remark 3.20. Let $p_{1}, p_{2} \in R$ be two non-constant differential polynomials. If $p_{1}$ and $p_{2}$ have the same leader $\left(u_{k}\right)_{J}$ and the degree of $p_{1}$ in $\left(u_{k}\right)_{J}$ is greater than or equal to the degree of $p_{2}$ in $\left(u_{k}\right)_{J}$, then the same pseudo-division as in (3.3) yields a remainder which is either zero, or has leader less than $\left(u_{k}\right)_{J}$, or has leader $\left(u_{k}\right)_{J}$ and smaller degree in $\left(u_{k}\right)_{J}$ than $p_{1}$.

More generally, if $\operatorname{ld}\left(p_{1}\right)=\theta \operatorname{ld}\left(p_{2}\right)$ for some $\theta \in \operatorname{Mon}(\Delta)$, then this pseudo-division can be applied with $p_{2}$ replaced with $\theta p_{2}$. Note that, by condition (b) of the definition of a ranking, we have $\operatorname{ld}\left(\theta p_{2}\right)=\theta \operatorname{ld}\left(p_{2}\right)$, and that, if $\theta \neq 1$, the degree of $\theta p_{2}$ in $\theta \operatorname{ld}\left(p_{2}\right)$ is one, so that the reduction can be applied without assumption on the degree of $p_{2}$ in $\operatorname{ld}\left(p_{2}\right)$. Then $c_{1}$ in (3.3) is again chosen as a suitable power of $\operatorname{init}\left(\theta p_{2}\right)$. In case $\theta \neq 1$ we have

$$
\operatorname{init}\left(\theta p_{2}\right)=\frac{\partial p_{2}}{\partial \operatorname{ld}\left(p_{2}\right)}=: \operatorname{sep}\left(p_{2}\right)
$$

and this differential polynomial is referred to as the separant of $p_{2}$.
In order not to change the solution set of a differential system, when $p_{1}=0$ is replaced with $r=0$, where $r$ is the result of a reduction of $p_{1}$ modulo $p_{2}$ or $\theta p_{2}$ as above, it is assumed that $\operatorname{init}\left(p_{2}\right)$ and $\operatorname{sep}\left(p_{2}\right)$ do not vanish on the solution set of the system. By definition of the separant and the discriminant (cf. Definition $3.6(\mathrm{~d})$ ), non-vanishing of $\operatorname{sep}\left(p_{2}\right)$ follows from non-vanishing of $\operatorname{disc}\left(p_{2}\right)$, as ensured by the algebraic part of Thomas' algorithm (cf. Remark 3.10).

We assume now that the given differential system is simple as an algebraic system (cf. Definition 3.7); it could be one of the systems resulting from the algebraic part of Thomas' algorithm.
Remark 3.21. The symmetry of the second derivatives $\partial_{i} \partial_{j} u_{k}=\partial_{j} \partial_{i} u_{k}$ (and similarly for higher order derivatives) imposes necessary conditions on the solvability of a system of partial differential equations. Taking identities like these into account and forming linear combinations of (derivatives of) the given equations may produce differential consequences with lower ranked leaders. As already discussed in the case of systems of linear PDEs above, in order to obtain a complete set of algebraic conditions on the Taylor coefficients of an analytic solution, the system has to include these integrability conditions. If a system of partial differential equations admits a translation into algebraic conditions on the Taylor coefficients such that no further integrability conditions have to be taken into account, then it is said to be formally integrable.
Definition 3.22. Each equation $p_{i}=0$ in a differential system is assigned the set of admissible derivations $\mu\left(\theta_{i}, M_{k}\right)$, where $\operatorname{ld}\left(p_{i}\right)=\theta_{i} u_{k}$ and

$$
\begin{equation*}
M_{k}:=\left\{\theta \in \operatorname{Mon}(\Delta) \mid \theta u_{k} \in\left\{\operatorname{ld}\left(p_{1}\right), \ldots, \operatorname{ld}\left(p_{s}\right)\right\}\right\} \tag{3.5}
\end{equation*}
$$

is the set of all monomials which define leaders $\operatorname{ld}\left(p_{i}\right)$ involving the same differential indeterminate $u_{k}$. We refer to $d p_{i}$ for $d \in \operatorname{Mon}\left(\mu\left(\theta_{i}, M_{k}\right)\right)$ as the admissible derivatives of $p_{i}$.

Formal integrability of a differential system is then decided by applying to each equation $p_{i}=0$ every of its non-admissible derivations $d \in \bar{\mu}\left(\theta_{i}, M_{k}\right)$ and computing the pseudo-remainder of $d p_{i}$ modulo $p_{1}, \ldots, p_{s}$ and their admissible derivatives. The restriction of the pseudo-division to admissible derivatives requires $M_{k}$ to be Janet complete (cf. Definition 2.18). If one of these pseu-do-remainders is non-zero, then it is added as a new equation to the system, and the augmented system has to be treated by the algebraic part of Thomas' algorithm again.
Definition 3.23. A system $\left\{p_{1}=0, \ldots, p_{s}=0\right\}$ of PDEs, where $p_{1}, \ldots, p_{s} \in R \backslash K$, is said to be passive if the following two conditions hold for $\operatorname{ld}\left(p_{1}\right)=\theta_{1} u_{k_{1}}, \ldots, \operatorname{ld}\left(p_{s}\right)=\theta_{s} u_{k_{s}}$, where $\theta_{i} \in \operatorname{Mon}(\Delta), k_{i} \in\{1, \ldots, m\}$.
(a) For all $k \in\{1, \ldots, m\}$, the set $M_{k}$ defined in (3.5) is Janet complete.
(b) For all $i \in\{1, \ldots, s\}$ and all non-admissible derivations $d \in \bar{\mu}\left(\theta_{i}, M_{k_{i}}\right)$, the pseudo-remainder of $d p_{i}$ modulo $p_{1}, \ldots, p_{s}$ and their admissible derivatives is zero.

Definition 3.24. A differential system $S$, defined over $R$, as in (3.4) is said to be simple (with respect to the ranking $>$ ) if the following three conditions hold.
(a) The system $S$ is simple as an algebraic system (in the finitely many variables $\left(u_{k}\right)_{J}$ which occur in the equations and inequations of $S$, totally ordered by $>$ ).
(b) The system $\left\{p_{1}=0, \ldots, p_{s}=0\right\}$ is passive.
(c) The left hand sides of the inequations $q_{1} \neq 0, \ldots, q_{t} \neq 0$ equal their pseudo-remainders modulo $p_{1}, \ldots, p_{s}$ and their derivatives.
Definition 3.25. Let $S$ be a differential system, defined over $R$. A Thomas decomposition of $S$ (or of $\operatorname{Sol}_{\Omega}(S)$ ) with respect to the ranking $>$ is a collection of finitely many simple differential systems $S_{1}, \ldots, S_{r}$, defined over $R$, such that the solution set $\operatorname{Sol}_{\Omega}(S)$ of $S$ is the disjoint union of the solution sets $\operatorname{Sol}_{\Omega}\left(S_{1}\right), \ldots, \operatorname{Sol}_{\Omega}\left(S_{r}\right)$.
Remark 3.26. Given $S$ as in (3.4) and a ranking on $R$, a Thomas decomposition of $S$ with respect to $>$ can be computed by interweaving the algebraic part discussed in Subsection 3.1 and differential reduction and completion with respect to Janet division.

First of all, a Thomas decomposition of $S$, considered as an algebraic system, is computed. Each of the resulting simple algebraic systems is then treated as follows. Differential pseudo-division is applied to pairs of distinct equations with leaders $\theta_{1} u_{k}$ and $\theta_{2} u_{k}$ such that $\theta_{1} \mid \theta_{2}$ until either a non-zero pseudo-remainder is obtained or no such further reductions are possible. Non-zero pseu-do-remainders are added to the system, the algebraic part of Thomas' algorithm is applied again, and the process is repeated. Once the system is auto-reduced in this sense, then it is possibly augmented with certain derivatives of equations so that the sets $M_{k}$ defined in (3.5) are Janet complete. Then it is checked whether the system is passive. If a non-zero remainder is obtained by a pseudo-division of a non-admissible derivative modulo the equations and their admissible derivatives, then the algebraic part of Thomas' algorithm is applied again to the augmented system. Otherwise, the system is passive. Finally, the left hand side of each inequation is replaced with its pseudo-remainder modulo the equations and their derivatives, in order to ensure condition (c) of Definition 3.24. The main reason why this procedure terminates is Dickson's Lemma, which shows that the ascending sequence of ideals of the semigroup $\operatorname{Mon}(\Delta)$ formed by the monomials $\theta$ defining leaders of equations (for each differential indeterminate) becomes stationary after finitely many steps.

For more details on the differential part of Thomas' algorithm, we refer to [BGL+12], [LH14], and [Rob14, Subsect. 2.2.2].

An implementation of Thomas' algorithm for differential systems was developed by M. LangeHegermann at RWTH Aachen University as Maple package DifferentialThomas [BLH], [GLR19].

When displaying a simple differential system we indicate next to each equation its set of admissible derivations.
Example 3.27. Let us consider the ODE (discussed in [Inc56, Example in Sect. 4.7])

$$
\left(\frac{\partial u}{\partial x}\right)^{3}-4 x u(x) \frac{\partial u}{\partial x}+8 u(x)^{2}=0 .
$$

The left hand side is represented by the element $p:=u_{x}^{3}-4 x u u_{x}+8 u^{2}$ of the differential polynomial ring $R=K\{u\}$ with one derivation $\partial_{x}$, where $K=\mathbb{Q}(x)$ is the field of rational functions in $x$, endowed with differentiation with respect to $x$.

The initial of $p$ is constant, the separant of $p$ is $3 u_{x}^{2}-4 x u$. The algebraic part of Thomas' algorithm only distinguishes the cases whether the discriminant of $p$ vanishes or not. We have

$$
\operatorname{disc}(p)=-\operatorname{res}\left(p, \operatorname{sep}(p), u_{x}\right)=-64 \underline{u}^{3}\left(27 \underline{u}-4 x^{3}\right) .
$$

This case distinction leads to the Thomas decomposition

$$
\begin{aligned}
\underline{u x}^{3}-4 x u \underline{u_{x}}+8 u^{2} & =0, \quad\left\{\partial_{x}\right\} \\
\left(27 \underline{u}-4 x^{3}\right) \underline{u} & \neq 0
\end{aligned}
$$

$$
\left(27 \underline{u}-4 x^{3}\right) \underline{u}=0, \quad\left\{\partial_{x}\right\}
$$

Since both systems contain only one equation, no differential reductions are necessary. The second simple system could be split into two with equations $27 u-4 x^{3}=0$ and $u=0$, respectively. The solutions of the first simple system are given by $u(x)=c(x-c)^{2}$, where $c$ is an arbitrary non-zero constant. The solutions $u(x)=0$ and $u(x)=\frac{4}{27} x^{3}$ of the second simple system are called singular solutions, the latter one being an envelope of the general solution.

More about singular solutions can be found, e.g., in [Dar73], [Ham93], [Rit36], [Hub97].
Example 3.28. Let us compute a Thomas decomposition of the system of (nonlinear) PDEs

$$
\left\{\begin{aligned}
\frac{\partial^{2} u}{\partial x^{2}}-\frac{\partial^{2} u}{\partial y^{2}} & =0 \\
\frac{\partial u}{\partial x}-u^{2} & =0
\end{aligned}\right.
$$

for one unknown function $u(x, y)$. We define the elements $p_{1}:=u_{x, x}-u_{y, y}$ and $p_{2}:=u_{x}-u^{2}$ of the differential polynomial ring $R=\mathbb{Q}\{u\}$ with commuting derivations $\partial_{x}, \partial_{y}$. We choose the degree-reverse lexicographical ranking $>$ on $R$ with $\partial_{x} u>\partial_{y} u$ (cf. Example 3.19).

Since the monomial $\partial_{x}$ defining the leader of $p_{2}$ divides the monomial $\partial_{x}^{2}$ defining the leader of $p_{1}$, differential pseudo-division is applied and $p_{1}$ is replaced with

$$
p_{3}:=p_{1}-\partial_{x} p_{2}-2 u p_{2}=-u_{y, y}+2 u^{3} .
$$

Janet division associates the sets of admissible derivations to the equations as follows:

$$
\left\{\begin{array}{rl}
\underline{u_{x}}-u^{2} & =0,
\end{array} \quad\left\{\partial_{x}, \partial_{y}\right\}, ~=0, \quad\left\{*, \partial_{y}\right\}\right.
$$

The set of monomials $\left\{\partial_{x}, \partial_{y}^{2}\right\}$ defining the leaders $u_{x}$ and $u_{y, y}$ is Janet complete. The check whether the above system is passive involves the following reduction:

$$
\partial_{x} p_{3}+\partial_{y}^{2} p_{2}-6 u^{2} p_{2}-2 u p_{3}=-2\left(\underline{u_{y}}+u^{2}\right)\left(\underline{u_{y}}-u^{2}\right)
$$

This non-zero remainder is a differential consequence which is added as an equation to the system. In fact, the system can be split into two systems according to the given factorization. For both systems a differential reduction of $p_{3}$ modulo the chosen factor is applied because the monomial $\partial_{y}$ defining the new leader divides the monomial $\partial_{y}^{2}$ defining $\operatorname{ld}\left(p_{3}\right)$. In both cases the remainder is zero, the sets of monomials defining leaders are Janet complete, and the passivity checks confirm formal integrability. We obtain the Thomas decomposition

$$
\begin{aligned}
& \underline{u_{x}}-u^{2}=0, \quad\left\{\partial_{x}, \partial_{y}\right\} \\
& \underline{u_{y}}+u^{2}=0, \quad\left\{*, \partial_{y}\right\}
\end{aligned}
$$

$$
\begin{aligned}
\underline{u_{x}}-u^{2} & =0, \quad\left\{\partial_{x}, \partial_{y}\right\} \\
\underline{u_{y}}-u^{2} & =0, \quad\left\{*, \partial_{y}\right\} \\
\underline{u} & \neq 0
\end{aligned}
$$

If the above factorization is ignored, then the discriminant of $p_{4}:=u_{y}^{2}-u^{4}$ needs to be considered, which implies vanishing or non-vanishing of the separant $2 u_{y}$. This case distinction leads to a different Thomas decomposition.

Exercise. Complete the computation of the alternative Thomas decomposition at the end of the previous example.
Remark 3.29. A Thomas decomposition of a differential system is not uniquely determined in general (cf. also Remark 3.12). In the special case of a system $S$ of linear partial differential equations no case distinctions are necessary, and the single simple system in any Thomas decomposition of $S$ is a Janet basis for $S$.
Exercise. Compute a Thomas decomposition of the differential system given in Example 3.1. For example, with respect to the degree-reverse lexicographical ranking satisfying $\partial_{t} u>\partial_{x} u$, a Thomas decomposition of that differential system is given by (cf. also [Rob14, Ex. 2.2.61])

| $u=0, \quad\left\{\partial_{t}, \partial_{x}\right\}$ | $\begin{aligned} \underline{u_{t}}-6 u u_{x} & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\ u_{x, x} & =0, \quad\left\{*, \partial_{x}\right\} \\ u & \neq 0, \end{aligned}$ | $\begin{aligned} u_{t} & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\ \underline{u_{x, x, x}}-6 u u_{x} & =0, \quad\left\{*, \partial_{x}\right\} \\ u_{x, x} & \neq 0 \\ u & \neq 0 \end{aligned}$ |
| :---: | :---: | :---: |

Determine the analytic solutions of each of these simple differential systems.

Example 3.30 ([GLR19], Ex. 14). For computing a Thomas decomposition of the Navier-Stokes equations given in Example 3.2 (with $\left.\rho \equiv 1,\left(g_{1}, g_{2}, g_{3}\right) \equiv(0,0,0)\right)$ we choose the ranking $>$ on $\mathbb{Q}(\mu)\left\{v_{1}, v_{2}, v_{3}, p\right\}$ with commuting derivations $\partial_{t}, \partial_{x}, \partial_{y}, \partial_{z}$ as follows. First the differential operators corresponding to the jet variables in question are compared with respect to the degreereverse lexicographical ranking satisyfing $\partial_{t}>\partial_{x}>\partial_{y}>\partial_{z}$. In case of equality of the differential operators the respective differential indeterminates are compared according to $v_{1}>v_{2}>v_{3}>p$ (cf. also the definition of the term-over-position ordering in Example 2.31). Then a Thomas decomposition with respect to $>$ consists of one simple differential system:

$$
\begin{aligned}
& \underline{\left(v_{1}\right)_{x}}+\left(v_{2}\right)_{y}+\left(v_{3}\right)_{z}=0, \quad\left\{\partial_{t}, \partial_{x}, \partial_{y}, \partial_{z}\right\} \\
& \mu\left(v_{1}\right)_{y, y}+\mu\left(v_{1}\right)_{z, z}-\mu \frac{\left(v_{2}\right)_{x, x}+\mu\left(v_{2}\right)_{y, y}+\mu\left(v_{2}\right)_{z, z}-v_{1}\left(v_{2}\right)_{x, y}}{}-\mu\left(v_{3}\right)_{x, z}+v_{1}\left(v_{2}\right)_{y}+v_{1}\left(v_{3}\right)_{z} \\
&-v_{2}\left(v_{2}\right)_{y}-v_{3}\left(v_{2}\right)_{z}-p_{y}-\left(v_{2}\right)_{t}=0, \quad\left\{\partial_{t}, \partial_{x}, \partial_{y}, \partial_{z}\right\} \\
&-v_{2}\left(v_{1}\right)_{y}-v_{3}\left(v_{1}\right)_{z}-p_{x}-\left(v_{1}\right)_{t}=0, \quad\left\{\partial_{t}, *, \partial_{y}, \partial_{x}\right\} \\
& \mu \underline{\left(v_{3}\right)_{x, x}+\mu\left(v_{3}\right)_{y, y}+\mu\left(v_{3}\right)_{z, z}-v_{1}\left(v_{3}\right)_{x}} \\
&-v_{2}\left(v_{3}\right)_{y}-v_{3}\left(v_{3}\right)_{z}-p_{z}-\left(v_{3}\right)_{t}=0, \quad\left\{\partial_{t}, \partial_{x}, \partial_{y}, \partial_{z}\right\} \\
& \underline{p_{x, x}}+p_{y, y}+p_{z, z}+2\left(v_{1}\right)_{y}\left(v_{2}\right)_{x}+2\left(v_{1}\right)_{z}\left(v_{3}\right)_{x}+2\left(v_{2}\right)_{y}^{2} \\
&+2\left(v_{2}\right)_{y}\left(v_{3}\right)_{z}+2\left(v_{2}\right)_{z}\left(v_{3}\right)_{y}+2\left(v_{3}\right)_{z}^{2}=0 . \quad\left\{\partial_{t}, \partial_{x}, \partial_{y}, \partial_{z}\right\}
\end{aligned}
$$

The last equation is obtained as

$$
\partial_{x} A_{1}+\partial_{y} A_{2}+\partial_{z} A_{3}+\left(\mu \Delta-\partial_{t}-v_{1} \partial_{x}-v_{2} \partial_{y}-v_{3} \partial_{z}+2\left(v_{2}\right)_{y}+2\left(v_{3}\right)_{z}-A_{4}\right) A_{4}
$$

where $A_{1}, A_{2}, A_{3}, A_{4}$ are the differential polynomials given by the equations in Example 3.2 and $\Delta$ is the Laplace operator. (Modulo the other equations in the system the last equation is the Poisson pressure equation.)

This simple system of the Thomas decomposition allows to enumerate the Taylor coefficients of $v_{1}(t, x, y, z), v_{2}(t, x, y, z), v_{3}(t, x, y, z), p(t, x, y, z)$ whose values can be chosen arbitrarily in a power series solution to the Navier-Stokes equations (similarly to Example 2.44). Janet decompositions of the sets of parametric derivatives for the differential indeterminates $v_{1}, v_{2}, v_{3}, p$ are given by

$$
\begin{aligned}
& v_{1}:\left\{\begin{aligned}
1, & \left\{\partial_{t}, *, \partial_{y}, \partial_{z}\right\}
\end{aligned}\right. \\
& v_{2}:
\end{aligned}\left\{\begin{aligned}
1, & \left\{\partial_{t}, *, \partial_{y}, \partial_{z}\right\} \\
\partial_{x}, & \left\{\partial_{t}, *, *, \partial_{z}\right\}
\end{aligned}\right\}
$$

Hence, the corresponding generalized Hilbert series are

$$
\frac{1}{\left(1-\partial_{t}\right)\left(1-\partial_{y}\right)\left(1-\partial_{z}\right)}
$$

for $v_{1}(t, x, y, z)$,

$$
\frac{1}{\left(1-\partial_{t}\right)\left(1-\partial_{y}\right)\left(1-\partial_{z}\right)}+\frac{\partial_{x}}{\left(1-\partial_{t}\right)\left(1-\partial_{z}\right)}
$$

for $v_{2}(t, x, t, z)$, and for $v_{3}(t, x, y, z)$ and for $p(t, x, y, z)$

$$
\frac{1}{\left(1-\partial_{t}\right)\left(1-\partial_{y}\right)\left(1-\partial_{z}\right)}+\frac{\partial_{x}}{\left(1-\partial_{t}\right)\left(1-\partial_{y}\right)\left(1-\partial_{z}\right)}
$$

Therefore, extending the Cauchy-Kovalevskaya Theorem, a Cauchy problem for the Navier-Stokes equations around an arbitrary point $\left(t_{0}, x_{0}, y_{0}, z_{0}\right)$ may be posed as follows:

$$
\left\{\begin{aligned}
v_{1}\left(t, x_{0}, y, z\right) & =f_{1}(t, y, z) \\
v_{2}\left(t, x_{0}, y, z\right) & =f_{2}(t, y, z) \\
\frac{\partial v_{2}}{\partial x}\left(t, x_{0}, y_{0}, z\right) & =f_{3}(t, z) \\
v_{3}\left(t, x_{0}, y, z\right) & =f_{4}(t, y, z) \\
\frac{\partial v_{3}}{\partial x}\left(t, x_{0}, y, z\right) & =f_{5}(t, y, z) \\
p\left(t, x_{0}, y, z\right) & =f_{6}(t, y, z) \\
\frac{\partial p}{\partial x}\left(t, x_{0}, y, z\right) & =f_{7}(t, y, z)
\end{aligned}\right.
$$

where $f_{1}, f_{2}, \ldots, f_{7}$ are arbitrary functions of their arguments which are analytic around the point $\left(t_{0}, x_{0}, y_{0}, z_{0}\right)$. The arbitrariness of analytic solutions is determined by $f_{1}, f_{2}, \ldots, f_{7}$.

Pseudo-reduction of a differential polynomial modulo the equations of a simple differential system and their derivatives decides membership to the corresponding saturation ideal (cf. also Proposition 3.13).
Proposition 3.31 ([Rob14], Prop. 2.2.50). Let $S$ be a simple differential system, defined over $R$, with equations $p_{1}=0, p_{2}=0, \ldots, p_{s}=0$. Moreover, let $E$ be the differential ideal of $R$ which is generated by $p_{1}, \ldots, p_{s}$ and define the product $q$ of the initials and separants of all $p_{1}, \ldots, p_{s}$. Then $E: q^{\infty}$ is a radical differential ideal. Given $p \in R$, we have $p \in E: q^{\infty}$ if and only if the pseudo-remainder of $p$ modulo $p_{1}, \ldots, p_{s}$ and their derivatives is zero.

Similarly to the algebraic case treated in the previous section, the Nullstellensatz for analytic functions (due to J. F. Ritt and H. W. Raudenbush, cf. [Rit50, Sects. II.7-11, IX.27]) establishes a one-to-one correspondence of solution sets $V:=\operatorname{Sol}_{\Omega}(S)$ of systems of partial differential equations $S=\left\{p_{1}=0, \ldots, p_{s}=0\right\}$ for $m$ unknown functions, defined over $R$, and their vanishing ideals in $R=K\left\{u_{1}, \ldots, u_{m}\right\}$

$$
\mathcal{I}_{R}(V):=\{p \in R \mid p(f)=0 \text { for all } f \in V\}
$$

These are the radical differential ideals of $R$. The Nullstellensatz implies that, with the notation of Proposition 3.31, we have $\mathcal{I}_{R}\left(\operatorname{Sol}_{\Omega}(S)\right)=E: q^{\infty}$.

The following proposition allows to decide whether or not a given differential equation $p=0$ is a consequence of a (not necessarily simple) differential system $S$ by applying pseudo-reductions to $p$ modulo each of the simple systems in a Thomas decomposition of $S$.

Proposition 3.32 ([Rob14], Prop. 2.2.72). Let a (not necessarily simple) differential system $S$ be given by $p_{1}=0, p_{2}=0, \ldots, p_{s}=0, q_{1} \neq 0, q_{2} \neq 0, \ldots, q_{t} \neq 0$, and let $S_{1}, \ldots, S_{r}$ be a Thomas decomposition of $S$ with respect to any ranking on $R$. Moreover, let $E$ be the differential ideal of $R$ generated by $p_{1}, \ldots, p_{s}$ and define the product $q$ of $q_{1}, \ldots, q_{t}$. For $i \in\{1, \ldots, r\}$, let $E^{(i)}$ be the differential ideal of $R$ generated by the equations in $S_{i}$ and define the product $q^{(i)}$ of the initials and separants of all these equations in $S_{i}$. Then we have

$$
\sqrt{E: q^{\infty}}=\left(E^{(1)}:\left(q^{(1)}\right)^{\infty}\right) \cap \ldots \cap\left(E^{(r)}:\left(q^{(r)}\right)^{\infty}\right)
$$

An important class of rankings can be defined as follows (following C. Riquier [Riq10, no. 102]).
Remark 3.33. Let the map $\varphi: \operatorname{Mon}(\Delta) u \rightarrow \mathbb{Q}^{(n+m) \times 1}=\mathbb{Q}^{n \times 1} \oplus \mathbb{Q}^{m \times 1}$ be defined by

$$
\partial^{I} u_{k} \longmapsto\left(I, e_{k}\right)^{\top}, \quad I \in\left(\mathbb{Z}_{\geq 0}\right)^{n}, \quad k=1, \ldots, m
$$

where $e_{1}, \ldots, e_{m}$ are the standard basis vectors of $\mathbb{Q}^{m \times 1}$. Then every matrix $M \in \mathbb{Q}^{r \times(n+m)}$ defines an irreflexive and transitive relation $>$ on $\operatorname{Mon}(\Delta) u$ by

$$
\begin{equation*}
v>w \quad: \Longleftrightarrow \quad M \varphi(v)>M \varphi(w), \quad v, w \in \operatorname{Mon}(\Delta) u \tag{3.6}
\end{equation*}
$$

where vectors on the right hand side are compared lexicographically. Assume that $M$ admits a left inverse (in particular, we have $r \geq n+m$ ). Then the linear map $\mathbb{Q}^{(n+m) \times 1} \rightarrow \mathbb{Q}^{r \times 1}$ induced by $M$ is injective, and $>$ is a total ordering on $\operatorname{Mon}(\Delta) u$. Linearity of matrix multiplication implies that $>$ satisfies condition (b) of Definition 3.17, p. 21, of a ranking. Moreover, condition (a) of the same definition holds if and only if, for each $j=1, \ldots, n$, the first non-zero entry of the $j$-th column of $M$ is positive. Every ranking $>$ defined by (3.6) is a Riquier ranking, i.e.,

$$
\theta_{1} u_{i}>\theta_{2} u_{i} \quad \Longleftrightarrow \quad \theta_{1} u_{j}>\theta_{2} u_{j}
$$

holds for all $\theta_{1}, \theta_{2} \in \operatorname{Mon}(\Delta), 1 \leq i, j \leq m$.
In every equation $p=0$ of a simple differential system $S$ we can solve for the term containing the highest power of the leader $\operatorname{ld}(p)$ to obtain an equivalent equation

$$
\operatorname{init}(p) \operatorname{ld}(p)^{k}=r
$$

where $r$ consists of terms which involve lower powers of $\operatorname{ld}(p)$ than the one on the left hand side or whose leaders are ranked lower than $\operatorname{ld}(p)$. Moreover, the differential polynomial init $(p)$ does not vanish for any solution of the simple system $S$. We obtain a generalization of the Cauchy-Kovalevskaya Theorem (cf. Theorem 1.1); cf. also [Tho28], [Tho34], [Ger09], [RRW99], [GLR19].
Corollary 3.34. Let $S$ be a simple differential system as in (3.4). Suppose that $\left(z_{1}^{0}, \ldots, z_{n}^{0}\right)$ is a point where all $p_{1}, \ldots, p_{s}$ and all $q_{1}, \ldots, q_{t}$ are defined and such that no initial or separant of any of these differential polynomials vanishes. Let formal power series around $\left(z_{1}^{0}, \ldots, z_{n}^{0}\right)$ be defined by

$$
f_{k}:=\sum_{J \in\left(\mathbb{Z}_{\geq 0}\right)^{n}} c_{k, J} \frac{\left(z_{1}-z_{1}^{0}\right)^{J_{1}}}{J_{1}!} \ldots \frac{\left(z_{n}-z_{n}^{0}\right)^{J_{n}}}{J_{n}!}, \quad k=1, \ldots, m
$$

with Taylor coefficients $c_{k, J}$. Then any assignment of values to $c_{k, J}$ for all $J \in\left(\mathbb{Z}_{\geq 0}\right)^{n}$ such that $\partial^{J} u_{k}$ is not a principal derivative and

$$
\left.q_{j}\left(f_{1}, \ldots, f_{m}\right)\right|_{\left(z_{1}, \ldots, z_{n}\right)=\left(z_{1}^{0}, \ldots, z_{n}^{0}\right)} \neq 0, \quad j=1, \ldots, t
$$

gives rise to formal power series solutions

$$
u_{1}\left(z_{1}, \ldots, z_{n}\right)=f_{1}, \quad \ldots, \quad u_{m}\left(z_{1}, \ldots, z_{n}\right)=f_{m}
$$

of $S$ around $\left(z_{1}^{0}, \ldots, z_{n}^{0}\right)$ determined by the consistent system of algebraic equations for $c_{k, J}$

$$
\left.\left(\partial^{J} p_{i}\right)\left(f_{1}, \ldots, f_{m}\right)\right|_{\left(z_{1}, \ldots, z_{n}\right)=\left(z_{1}^{0}, \ldots, z_{n}^{0}\right)}=0, \quad J \in\left(\mathbb{Z}_{\geq 0}\right)^{n}, \quad i=1, \ldots, s
$$

and conversely, every formal power series solution of $S$ around $\left(z_{1}^{0}, \ldots, z_{n}^{0}\right)$ stems from such an assignment. If $>$ is an orderly Riquier ranking, then sufficiently generic initial conditions determined by convergent power series yield convergent power series solutions.
3.3. Elimination. Thomas' algorithm can be used to solve various differential elimination problems. This subsection presents results on certain rankings on the differential polynomial ring $R=K\left\{u_{1}, \ldots, u_{m}\right\}$ which allow to compute all differential consequences of a given differential system involving only a specified subset of the differential indeterminates $u_{1}, \ldots, u_{m}$. In other words, this technique allows to determine all differential equations which are satisfied by certain components of the solution tuples. We adopt the notation from the previous subsection.

Definition 3.35. Let $I_{1}, I_{2}, \ldots, I_{k}$ form a partition of $\{1,2, \ldots, m\}$ such that $i_{1} \in I_{j_{1}}, i_{2} \in I_{j_{2}}$, $i_{1} \leq i_{2}$ implies $j_{1} \leq j_{2}$. Let $B_{j}:=\left\{u_{i} \mid i \in I_{j}\right\}, j=1, \ldots, k$. Moreover, fix some degree-reverse lexicographical ordering $>$ on $\operatorname{Mon}(\Delta)$. Then the block ranking on $R$ with blocks $B_{1}, \ldots, B_{k}$ (with $\left.u_{1}>u_{2}>\ldots>u_{m}\right)$ is defined for $\theta_{1} u_{i_{1}}, \theta_{2} u_{i_{2}} \in \operatorname{Mon}(\Delta) u$, where $u_{i_{1}} \in B_{j_{1}}, u_{i_{2}} \in B_{j_{2}}$, by

$$
\theta_{1} u_{i_{1}}>\theta_{2} u_{i_{2}} \quad \Longleftrightarrow\left\{\begin{array}{c}
j_{1}<j_{2} \quad \text { or }\left(j _ { 1 } = j _ { 2 } \text { and } \left(\theta_{1}>\theta_{2}\right.\right. \text { or } \\
\left.\left.\left(\theta_{1}=\theta_{2} \quad \text { and } i_{1}<i_{2}\right)\right)\right)
\end{array}\right.
$$

Such a ranking is said to satisfy $B_{1} \gg B_{2} \gg \ldots \gg B_{k}$.
Example 3.36. With respect to the block ranking on $K\left\{u_{1}, u_{2}, u_{3}\right\}$ with blocks $\left\{u_{1}\right\},\left\{u_{2}, u_{3}\right\}$ (and $\left.u_{1}>u_{2}>u_{3}\right)$ we have $\left(u_{1}\right)_{(0,1)}>u_{1}>\left(u_{2}\right)_{(1,2)}>\left(u_{3}\right)_{(1,2)}>\left(u_{2}\right)_{(0,1)}$.

In the situation of the previous definition, for every $i \in\{1, \ldots, k\}$, we consider

$$
K\left\{B_{i}, \ldots, B_{k}\right\}:=K\left\{u \mid u \in B_{i} \cup \ldots \cup B_{k}\right\}
$$

as a differential subring of $R=K\left\{u_{1}, \ldots, u_{m}\right\}$, endowed with the restrictions of the derivations $\partial_{1}, \ldots, \partial_{n}$ to $K\left\{B_{i}, \ldots, B_{k}\right\}$.

For any algebraic or differential system $S$ we denote by $S^{=}$(resp. $S^{\neq}$) the set of the left hand sides of all equations (resp. inequations) in $S$.

Proposition 3.37 ([Rob14], Prop. 3.1.36). Let $S$ be a simple differential system, defined over $R$, with respect to a block ranking with blocks $B_{1}, \ldots, B_{k}$. Moreover, let $E$ be the differential ideal of $R$ generated by $S^{=}$and $q$ the product of the initials and separants of all elements of $S^{=}$. For every $i \in\{1, \ldots, k\}$, let $E_{i}$ be the differential ideal of $K\left\{B_{i}, \ldots, B_{k}\right\}$ generated by $P_{i}:=$ $S=\cap K\left\{B_{i}, \ldots, B_{k}\right\}$ and let $q_{i}$ be the product of the initials and separants of all elements of $P_{i}$. Then, for every $i \in\{1, \ldots, k\}$, we have

$$
\left(E: q^{\infty}\right) \cap K\left\{B_{i}, \ldots, B_{k}\right\}=E_{i}: q_{i}^{\infty} .
$$

In other words, the differential equations implied by $S$ which involve only the differential indeterminates in $B_{i} \cup \ldots \cup B_{k}$ are precisely those whose pseudo-remainders modulo the elements of $S^{=} \cap K\left\{B_{i}, \ldots, B_{k}\right\}$ and their derivatives are zero.

The following well-known example could also be dealt with using Janet bases because the PDEs are linear.

Example 3.38. The Cauchy-Riemann equations for a complex function of $z=x+i y$ with real part $u$ and imaginary part $v$ are

$$
\left\{\begin{array}{l}
\frac{\partial u}{\partial x}-\frac{\partial v}{\partial y}=0 \\
\frac{\partial u}{\partial y}+\frac{\partial v}{\partial x}=0
\end{array}\right.
$$

The left hand sides are represented by the elements $p_{1}:=u_{x}-v_{y}$ and $p_{2}:=u_{y}+v_{x}$ of the differential polynomial ring $R=\mathbb{Q}\{u, v\}$ with derivations $\partial_{x}$ and $\partial_{y}$. Choosing a block ranking on $R$ satisfying $\{u\} \gg\{v\}$, the passivity check yields the equation

$$
\partial_{x} p_{2}-\partial_{y} p_{1}=v_{x, x}+v_{y, y}=0
$$

Similarly, the choice of a block ranking on the differential polynomial ring $R$ satisfying $\{v\} \gg\{u\}$ yields the consequence $u_{x, x}+u_{y, y}=0$. These computations confirm that the real and imaginary parts of a holomorphic function are harmonic functions.

Similarly to Proposition 3.32 we obtain a corollary to Proposition 3.37 for not necessarily simple differential systems.

Corollary 3.39 ([Rob14], Cor. 3.1.37). Let $S$ be a (not necessarily simple) differential system, defined over $R$, and $S_{1}, \ldots, S_{r}$ a Thomas decomposition of $S$ with respect to a block ranking with blocks $B_{1}, \ldots, B_{k}$. Moreover, let $E$ be the differential ideal of $R$ generated by $S^{=}$and $q$ the product of all elements of $S^{\neq}$. Let $i \in\{1, \ldots, k\}$ be fixed. For every $j \in\{1, \ldots, r\}$, let $E^{(j)}$ be the differential ideal of $K\left\{B_{i}, \ldots, B_{k}\right\}$ generated by $P_{j}:=S_{j}^{=} \cap K\left\{B_{i}, \ldots, B_{k}\right\}$ and let $q^{(j)}$ be the product of the initials and separants of all elements of $P_{j}$. Then we have

$$
\sqrt{E: q^{\infty}} \cap K\left\{B_{i}, \ldots, B_{k}\right\}=\left(E^{(1)}:\left(q^{(1)}\right)^{\infty}\right) \cap \ldots \cap\left(E^{(r)}:\left(q^{(r)}\right)^{\infty}\right) .
$$

We finish with two examples demonstrating applications of differential elimination.
Example 3.40 ([LHR], Ex. 10). A model of a 2-D crane is given by the following system of ordinary differential equations (cf. [FLMR95, Sect. 4.1] and the references therein), where $x(t)$ and $z(t)$ are the coordinates of the load of mass $m, \theta(t)$ is the angle between the rope and the $z$-axis (which points in the same direction as the gravitational force), $d(t)$ the trolley position, $T(t)$ the
tension of the rope, $R(t)$ the rope length, and $g$ the gravitational constant.

$$
\left\{\begin{aligned}
m \ddot{x} & =-T \sin \theta \\
m \ddot{z} & =-T \cos \theta+m g \\
x & =R \sin \theta+d \\
z & =R \cos \theta
\end{aligned}\right.
$$

We represent $\cos \theta$ and $\sin \theta$ by differential indeterminates $c$ and $s$ and add the generating relation $c^{2}+s^{2}=1$ to the differential system. Then this equation and the above equations define elements of the differential polynomial ring $\mathbb{Q}(m, g)\{T, c, s, d, R, x, z\}$ with derivation $\partial_{t}$, for which $m$ and $g$ are constants.

The task is to decide whether or not $\{x, z\}$ is a flat output of the system, i.e., whether or not
(a) no consequence of the differential system is a (differential) equation for $x$ and $z$ only, and
(b) for each of the other differential indeterminates $u \in\{T, c, s, d, R\}$ there exists a consequence $p=0$ of the differential system, where $p$ is a polynomial in that indeterminate $u$ (of differential order zero) with coefficients in $\mathbb{Q}(m, g)\{x, z\}$ and such that neither the leading coefficient of $p$ nor the derivative $\partial p / \partial u$ vanishes on the solution set of the system.

Condition (a) means that the residue classes of $x$ and $z$ modulo the differential ideal defined by the system are differentially algebraically independent over $\mathbb{Q}(m, g)$, whereas condition (b) implies that, for each of the differential indeterminates $T, c, s, d, R$, by the implicit function theorem, any polynomial $p$ as required can be solved locally for the indeterminate $u$ in the sense that the trajectory of $u$ can locally be expressed as an analytic function in terms of the trajectories of $x$ and $z$ (and their derivatives). For more details we refer to [LHR13], [LHR].

Choosing the block ranking $>$ on the differential polynomial ring $\mathbb{Q}(m, g)\{T, c, s, d, R, x, z\}$ which satisfies $\{T, c, s, d, R\} \gg\{x, z\}$ and $T>c>s>d>R$ and $x>z$, a Thomas decomposition of the given differential system with respect to $>$ is given by the following seven simple differential systems (where factorizations of left hand sides of equations have been taken into account).

$$
\begin{aligned}
& z \underline{T}+m z_{t, t} R-m g R=0, \quad\left\{\partial_{t}\right\} \\
& R \underline{c}-z=0, \quad\left\{\partial_{t}\right\} \\
& z_{t, t} R \underline{s}-g R \underline{s}-z x_{t, t}=0, \quad\left\{\partial_{t}\right\} \\
& z_{t, t} \underline{d}-g \underline{d}+z x_{t, t}-x z_{t, t}+g x=0, \quad\left\{\partial_{t}\right\} \\
& z_{t, t}^{2} \underline{R}^{2}-2 g z_{t, t} \underline{R}^{2}+g^{2} \underline{R}^{2}-z^{2} x_{t, t}^{2}-z^{2} z_{t, t}^{2}+2 g z^{2} z_{t, t}-g^{2} z^{2}=0, \quad\left\{\partial_{t}\right\} \\
& z\left(z_{t, t}-g\right) x_{t, t}\left(x_{t, t}^{2}-2 g z_{t, t}+g^{2}+z_{t, t}^{2}\right) \neq 0 \\
& \begin{array}{rlrl}
\underline{T} & =0, & \left\{\partial_{t}\right\} \\
R \underline{c}-z & =0, & \left\{\partial_{t}\right\} \\
R \underline{s}+d-x & =0, \quad\left\{\partial_{t}\right\} \\
\underline{d}^{2}-2 x \underline{d}+x^{2}-R^{2}+z^{2} & =0, \quad\left\{\partial_{t}\right\} \\
\underline{x_{t, t}} & =0, \quad\left\{\partial_{t}\right\} \\
\underline{z_{t, t}}-g & =0, \quad\left\{\partial_{t}\right\} \\
z \underline{R}(\underline{R}+z)(\underline{R}-z) & \neq 0
\end{array} \\
& \underline{T}-m z_{t, t}+m g=0,\left\{\partial_{t}\right\} \\
& \underline{c}+1=0, \quad\left\{\partial_{t}\right\} \\
& \underline{s}=0, \quad\left\{\partial_{t}\right\} \\
& \underline{d}-x=0,\left\{\partial_{t}\right\} \\
& \underline{R}+z=0, \quad\left\{\partial_{t}\right\} \\
& \underline{x_{t, t}}=0, \quad\left\{\partial_{t}\right\} \\
& \underline{z} \neq 0
\end{aligned}
$$

$$
\left.\begin{array}{rl}
\underline{T}+m z_{t, t}-m g & =0, \quad\left\{\partial_{t}\right\} \\
\underline{c}-1 & =0, \quad\left\{\partial_{t}\right\} \\
\underline{s} & =0, \quad\left\{\partial_{t}\right\} \\
\underline{d}-x & =0, \quad\left\{\partial_{t}\right\} \\
\underline{R}-z & =0, \quad\left\{\partial_{t}\right\} \\
\underline{x_{t, t}} & =0, \quad\left\{\partial_{t}\right\} \\
\underline{z} & \neq 0
\end{array}\right] \begin{aligned}
c \underline{T}-m g & =0, \quad\left\{\partial_{t}\right\} \\
\underline{c}+1 & =0, \quad\left\{\partial_{t}\right\} \\
\underline{s} & =0, \quad\left\{\partial_{t}\right\} \\
\underline{d}-x & =0, \quad\left\{\partial_{t}\right\} \\
\underline{R} & =0, \quad\left\{\partial_{t}\right\} \\
\frac{x_{t, t}}{} & =0, \quad\left\{\partial_{t}\right\} \\
\underline{z} & =0, \quad\left\{\partial_{t}\right\}
\end{aligned}
$$

We note that the first simple system $S_{1}$ contains no equation involving derivatives of $x$ and $z$ only, which shows that (a) is satisfied. Moreover, the other equations in $S_{1}$ show that (b) is satisfied as well. Hence, $\{x, z\}$ is a flat output of $S_{1}$. The remaining six simple differential systems describe particular configurations for which $\{x, z\}$ is not a flat output. In fact, the movement of the load is restricted by some constraint in these cases (e.g., $x_{t, t}=0$ or $z=0$, one reason being, e.g., that vanishing rope tension implies constant acceleration of the load, another being a constant rope length of zero allowing no vertical movement of the load).

Example 3.41. We consider the following system of nonlinear PDEs for $u(t, x), v(t, x), w(t, x)$ :

$$
\left\{\begin{align*}
u_{t, t}+w_{t, t}-v_{t}^{2} w_{x} & =0  \tag{3.7}\\
2 u w_{x} v_{t, t}+2 u_{t} v_{t} w_{x}+2 u v_{t} w_{t, x}+w_{t, x} & =0 \\
u_{t, t}+v_{t, x}+2 u v_{t} v_{t, x}+u_{x} v_{t}^{2} & =0
\end{align*}\right.
$$

These are the Euler-Lagrange equations obtained for a certain Lagrangian density. For more details on this example and the approach we refer to [GR16].

We would like to determine all consequences of (3.7) that are equations involving partial differentiation with respect to $t$ of at most order 1 . To this end we choose a ranking $>$ on the differential polynomial ring $\mathbb{Q}\{u, v, w\}$ with commuting derivations $\partial_{t}, \partial_{x}$ such that

$$
\partial_{t} \varphi>\partial_{x}^{i} \psi \quad \text { for all } \quad i \in \mathbb{Z}_{\geq 0}, \quad \varphi, \psi \in\{u, v, w\}
$$

For example, we may choose the ranking $>$ which first compares the differential operators corresponding to the jet variables in question with respect to the lexicographical ordering satisfying $\partial_{t}>\partial_{x}$ (cf. also Example 2.29, p. 9), and which, in case of equality of the differential operators, compares the respective differential indeterminates according to $u>v>w$. A Thomas decomposition of the given differential system with respect to $>$ is given by the following six simple differential systems, each of which allows to determine its consequences as required in a straightforward way.

| $\underline{u_{t, t}}+\left(2 u v_{t}+1\right) v_{t, x}+u_{x} v_{t}^{2}$ | $=0, \quad\left\{\partial_{t}, \partial_{x}\right\}$ |
| ---: | :--- |
| $2 u w_{x} \underline{v_{t, t}}+\left(2 u v_{t}+1\right) w_{t, x}+2 w_{x} v_{t} u_{t}$ | $=0$, |
| $\left(2 u v_{t}+1\right) v_{t, x}-\underline{w_{t, t}}+\left(u_{x}+w_{x}\right) v_{t}^{2}$ | $=0, \quad\left\{\partial_{x}\right\}$ |
| $w_{x}$ | $\neq 0$, |
| $u$ | $\neq 0$ |,$\quad$| $u$ | $=0, \quad\left\{\partial_{t}, \partial_{x}\right\}$ |
| ---: | :--- |
| $v_{t}$ | $=0, \quad\left\{\partial_{t}, \partial_{x}\right\}$ |
| $w_{t, t}$ | $=0, \quad\left\{\partial_{t}, \partial_{x}\right\}$ |
| $w_{t, x}$ | $=0, \quad\left\{*, \partial_{x}\right\}$ |
| $w_{x}$ | $\neq 0$ |


| $\begin{aligned} \underline{u_{t, t}}+\left(2 u v_{t}+1\right) v_{t, x}+u_{x} v_{t}^{2} & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\ \left(2 u v_{t}+1\right) \underline{v_{t, x, x}}+2 u v_{t, x}^{2}+4 u_{x} v_{t} v_{t, x}+v_{t}^{2} u_{x, x} & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\ \left(2 u v_{t}+1\right) v_{t, x}+u_{x} v_{t}^{2}-\underline{w_{t, t}} & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\ w_{t, x} & =0, \quad\left\{*, \partial_{x}\right\} \\ w_{x} & =0, \quad\left\{*, \partial_{x}\right\} \\ u & \neq 0, \\ 2 u \underline{v_{t}}+1 & \neq 0 \end{aligned}$ | $\begin{array}{rlrl} u=0, & \left\{\partial_{t}, \partial_{x}\right\} \\ v_{t, x} & =0, & \left\{\partial_{t}, \partial_{x}\right\} \\ w_{t, t}=0, & \left\{\partial_{t}, \partial_{x}\right\} \\ w_{t, x}=0, & \left\{*, \partial_{x}\right\} \\ w_{x}=0 & \left\{*, \partial_{x}\right\} \end{array}$ |
| :---: | :---: |

$$
\begin{aligned}
u_{t, t} & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\
u_{t, x} & =0, \quad\left\{*, \partial_{x}\right\} \\
u_{x} & =0, \quad\left\{*, \partial_{x}\right\} \\
2 u \underline{v_{t}}+1 & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\
w_{t, t} & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\
w_{t, x} & =0, \quad\left\{*, \partial_{x}\right\} \\
w_{x} & =0, \quad\left\{*, \partial_{x}\right\} \\
u & \neq 0
\end{aligned}
$$

$$
\begin{aligned}
u & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\
v_{t, x} & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\
\frac{w_{t, t}}{}-w_{x} v_{t}^{2} & =0, \quad\left\{\partial_{t}, \partial_{x}\right\} \\
w_{t, x} & =0, \quad\left\{*, \partial_{x}\right\} \\
w_{x, x} & =0, \quad\left\{*, \partial_{x}\right\} \\
w_{x} & \neq 0, \\
v_{t} & \neq 0
\end{aligned}
$$

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