THE CAUCHY-KOVALEVSKAYA THEOREM

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ABSTRACT. After presenting the basic analytic theory of ordinary differential equations, we discuss the Cauchy-Kovalevskaya theorem, characteristic surfaces, and the notion of well posedness. We include the fundamentals of analytic functions in the appendices.

Contents

1.	First order ordinary differential equations	1
2.	Systems of ordinary differential equations	4
3.	Linear partial differential equations	6
4.	Characteristic surfaces	9
5.	Well posedness	12
6.	The general Cauchy-Kovalevskaya theorem	14
Ap	ppendix A. Multivariate power series	17
Ap	ppendix B. Analytic functions	19

1. FIRST ORDER ORDINARY DIFFERENTIAL EQUATIONS

The Cauchy-Kovalevskaya theorem is a result on local existence of analytic solutions to a very general class of PDEs. However, it is best to start with the ODE case, which is simpler yet contains half the main ideas. Consider the problem

$$u' = f(u), \qquad u(0) = 0,$$
 (1)

where f is a given function analytic at 0, and u is the unknown function. Cauchy's theorem, proved by him during 1831-35, guarantees that a *unique solution exists that is analytic at* 0.

Remark 1. In (1) we may think of u as either a function of a real variable $x \in \mathbb{R}$, or a function of a complex variable $z \in \mathbb{C}$. If f is a real analytic function in a neighbourhood of $0 \in \mathbb{R}$, then it can be uniquely extended to a complex analytic function \tilde{f} in a neighbourhood of $0 \in \mathbb{C}$. Moreover, if u is a complex analytic solution of (1) in a neighbourhood of $0 \in \mathbb{C}$, with f replaced by \tilde{f} , then the restriction of u to \mathbb{R} will be real analytic and satisfy (1) with the original f. Thus we see that the complex analytic setting is more general, which we will assume henceforth.

Remark 2. Another simple observation is that the initial condition u(0) = 0 in (1) is not a loss of generality, since the initial value problem

$$v' = g(v), \qquad v(0) = v_0,$$
 (2)

is equivalent to (1), under the substitutions $v = v_0 + u$ and $f(u) = g(v_0 + u)$.

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With the intent of finding the Maclaurin series coefficients of u, we can repeatedly differentiate u' = f(u) to get

$$u'' = [f(u)]' = f'(u)u', \qquad u''' = [f(u)]'' = f''(u)(u')^2 + f'(u)u'', \qquad \dots$$
(3)

The Faà di Bruno formula would give the precise expression for $[f(u)]^{(m)}$, but without having to look up or derive that formula, just from the considerations (3) it is clear that

$$u^{(k)} = [f(u)]^{(k-1)} = q_k(f(u), \dots, f^{(k-1)}(u), u', \dots, u^{(k-1)}),$$
(4)

where q_k is a multivariate polynomial with *nonnegative* integer coefficients. We evaluate this at z = 0, and use u(0) = 0, to get

$$u^{(k)}(0) = q_k(f(0), \dots, f^{(k-1)}(0), u'(0), \dots, u^{(k-1)}(0)).$$
(5)

Now we repeatedly apply the same formula (with k having values k - 1, k - 2, etc.) to eliminate all $u^{(m)}(0)$ from the right hand side, inferring

$$u^{(k)}(0) = Q_k(f(0), \dots, f^{(k-1)}(0)),$$
(6)

with another multivariate polynomial Q_k having *nonnegative* integer coefficients. This incidentally proves uniqueness of analytic solutions to (1), since (6) fixes their Maclaurin series coefficients at 0. Moreover, provided that the Maclaurin series

$$u(z) = \sum_{n=0}^{\infty} \frac{u^{(n)}(0)}{n!} z^n,$$
(7)

with $u^{(n)}(0)$ given by (6) converges in a neighbourhood of 0, the function v = u' - f(u) is analytic at 0, and by construction, its Maclaurin series is identically zero. Hence, by the identity theorem v must vanish wherever it is defined, meaning that u' = f(u) there. Now it remains only to show that the series above converges in a neighbourhood of 0.

The heart of Cauchy's proof is his *method of majorants*, which is an ingenious and a very peculiar way of exploiting the positivity of the coefficients of Q_k against the underlying analytic setting. For two functions g and G, both infinitely differentiable at $c \in \mathbb{C}$, we say that G majorizes g at c, if

$$|g^{(k)}(c)| \le G^{(k)}(c), \qquad k = 0, 1, \dots$$
 (8)

In other words, the Taylor series coefficients of g at c is bounded in magnitude by the corresponding coefficients of G. Since our right hand side f is analytic at 0, there exist constants $M < \infty$ and r > 0 such that

$$\frac{|f^{(k)}(0)|}{k!} \le \frac{M}{r^k}, \qquad k = 0, 1, \dots$$
(9)

Then certainly the function

$$F(z) = \frac{M}{1 - z/r} = M + \frac{M}{r}z + \dots + \frac{M}{r^k}x^k + \dots,$$
 (10)

majorizes f at 0. Let us consider the initial value problem

$$U' = F(U), \qquad U(0) = 0.$$
 (11)

Then by (6) we have

$$U^{(k)}(0) = Q_k(F(0), \dots, F^{(k-1)}(0)),$$
(12)

and

$$|u^{(k)}(0)| = |Q_k(f(0), \dots, f^{(k-1)}(0))|$$

$$\leq Q_k(|f(0)|, \dots, |f^{(k-1)}(0)|)$$

$$\leq Q_k(F(0), \dots, F^{(k-1)}(0))$$

$$= U^{(k)}(0).$$
(13)

where we have used the nonnegativity of the coefficients of Q_k in the second an third lines, and the majorant property of F in the third line. The conclusion is that the solution u of the original problem (1) is majorized by the solution U of (11) at 0. Hence, if (11) has an analytic solution, u is automatically analytic. But (11) is easily solvable, with

$$U(z) = r(1 - \sqrt{1 - Mz/r}) = Mz/2 + \dots,$$
(14)

whose Taylor series around 0 has nonnegative coefficients. We have proved the following.

Theorem 3. The initial value problem

$$u' = f(u), \qquad u(0) = \eta_0,$$
(15)

with $f : \mathbb{C} \to \mathbb{C}$ analytic at η_0 , has a unique solution u that is analytic at 0.

Remark 4. From the majorant (14), the radius of convergence of the solution u can be estimated as $R \ge r/M$. Recalling that M > 0 and r > 0 are constants from the bounds (9), and recalling Cauchy's estimates (Cauchy 1831)

$$|f^{(n)}(0)| \le \frac{n!}{r^n} \sup_{|z|=r} |f(z)|, \tag{16}$$

that is valid if f is analytic on the disk $|z| \leq r$, we can estimate $M \leq \sup_{|z|=r} |f(z)|$. But there are

functions such as r/(r - Mz) that saturate Cauchy's estimates, meaning that M is essentially the magnitude of f in its domain of analyticity. Now, the magnitude of f is equal to the "speed" |u'|, hence the "time" it takes for u to become of magnitude r is roughly r/M. If we assume that f ceases to be analytic outside the disk $|z| \leq r$, these considerations imply that the convergence radius of u is roughly of order r/M, which cannot be improved in general.

Exercise 5. Give examples of f that makes the above statement precise.

Example 6. Let us consider the initial value problem

$$(u')^2 - u = 0, \qquad u(0) = \eta_0.$$
 (17)

If we want to write this in the form u' = f(u), we face the following obstacle: For each $\eta \in \mathbb{C} \setminus \{0\}$ there are two values of $p \in \mathbb{C}$ such that $p^2 - \eta = 0$. In other words, the function $\eta \mapsto \eta^{\frac{1}{2}}$ is many-valued. A way around this obstacle is to simply modify the formulation of the problem so that a branch of the function $\eta^{\frac{1}{2}}$ near the initial datum $\eta = \eta_0$ is also specified (This can be thought of as part of a generalized form of the initial condition). It can be done as follows. Notice that for any $\eta_0 \in \mathbb{C} \setminus \{0\}$ and $p_0 \in \mathbb{C}$ satisfying $p_0^2 - \eta_0 = 0$, there exists a unique (analytic) function f defined in a neighbourhood U of η_0 , such that $f(\eta)^2 - \eta = 0$ for $\eta \in U$ and $f(\eta_0) = p_0$. We assume that $u_0 \in \mathbb{C} \setminus \{0\}$ and $p_0 \in \mathbb{C}$ are given, and that they satisfy $p_0^2 - \eta_0 = 0$. The quantity p_0 is our device to specify the desired branch of $\eta^{\frac{1}{2}}$, and the condition $\eta_0 \neq 0$ is to make sure that such a branch is unique. Then instead of (17), we consider the problem

$$(u')^2 - u = 0, \qquad u(0) = \eta_0, \qquad u'(0) = p_0,$$
(18)

which, by Theorem 3, has a unique analytic solution in a neighbourhood of 0.

Theorem 7. Consider the initial value problem

$$F(u, u') = 0,$$
 $u(0) = \eta_0,$ $u'(0) = p_0,$ (19)

where $F : \mathbb{C} \times \mathbb{C} \to \mathbb{C}$ is analytic at $(\eta_0, p_0), F(\eta_0, p_0) = 0$, and

$$\frac{\partial F}{\partial p}(\eta_0, p_0) \neq 0. \tag{20}$$

Then there is a unique solution u that is analytic at 0.

Proof. By the analytic implicit function theorem, there exists a unique analytic function f defined in a neighbourhood U of η_0 , such that $F(\eta, f(\eta)) = 0$ for $\eta \in U$ and $f(\eta_0) = p_0$. By Theorem 3, there is a unique analytic solution v in a neighbourhood of 0, to the problem v' = f(v) and $v(0) = \eta_0$. This solution also satisfies (19) in a (possibly smaller) neighbourhood of 0, because $v'(0) = f(v(0)) = f(\eta_0) = p_0$. On the other hand, any analytic solution of (19) in a neighbourhood of 0 satisfies $u(0) = \eta_0$ and u' = f(u) in a (possibly smaller) neighbourhood of 0. By uniqueness of v, we must have $u \equiv v$.

2. Systems of ordinary differential equations

Our next step towards the Cauchy-Kovalevskaya theorem is Cauchy's existence theorem for the system:

$$u'_{j} = f_{j}(z, u_{1}, \dots, u_{m}), \qquad u_{j}(0) = 0, \qquad j = 1, \dots, m.$$
 (21)

We could have eliminated the dependence of f on z by introducing the new variable u_{m+1} with the equation $u'_{m+1} = 1$, but we intentionally leave it there in anticipation of the PDE case that is considered in the next section. The above equation can be written compactly as

$$u' = f(z, u), \qquad u(0) = 0,$$
(22)

with u and f having values in \mathbb{C}^m . We assume that the right hand side f is analytic in its m + 1 arguments. To determine the higher derivatives of u, we start differentiating the equation $u'_j = f_j(z, u)$ as

$$[f_j(z,u)]' = \partial_z f_j + \partial_{u_i} f_j \cdot u'_i, \qquad [f_j(z,u)]'' = \partial_z^2 f_j + 2\partial_z \partial_{u_i} f_j \cdot u'_i + \partial_{u_i} \partial_{u_\ell} f_j \cdot u'_i u'_\ell + \partial_{u_i} f_j \cdot u''_i,$$

where summation is taken over repeated indices. From here it is clear that

$$u_j^{(k)} = [f_j(z, u)]^{(k-1)} = q_k(\{\partial^\beta f_j(u)\}, \{u^{(\ell)}\}),$$
(23)

where q_k is a multivariate polynomial with nonnegative coefficients, and it is understood that the arguments of q_k are all $\partial^{\beta} f_j(u)$ with $|\beta| \leq k - 1$, and all components of all $u^{(\ell)}$ with $\ell \leq k - 1$. We evaluate this at z = 0, and use u(0) = 0, to get

$$u_j^{(k)}(0) = q_k(\{\partial^\beta f_j(0)\}, \{u^{(\ell)}(0)\}) = Q_{j,k}(\{\partial^\beta f_\ell(0)\}),$$
(24)

with $Q_{j,k}$ a multivariate polynomial having nonnegative coefficients. Note that the arguments of $Q_{j,k}$ are all $\partial^{\beta} f_{\ell}(0)$ with $|\beta| \leq k-1$ and $\ell = 1, \ldots, m$.

Having found that the derivatives of u at 0 is given by a positive coefficient polynomial of the derivatives of f at 0, we would like to replace f by a simpler majorant of it. We say G majorizes g at $c \in \mathbb{C}^n$, if

$$|\partial^{\alpha}g(c)| \le \partial^{\alpha}G(c), \quad \text{for all } \alpha.$$
 (25)

In other words, the Taylor series coefficients of g at c is bounded in magnitude by the corresponding coefficients of G. Since our right hand side f is componentwise analytic at 0, there exist constants $M < \infty$ and r > 0 such that

$$\frac{|\partial^{\alpha} f_j(0)|}{\alpha!} \le \frac{M}{r^{|\alpha|}}, \quad \text{for all } \alpha, \quad \text{and all } j.$$
(26)

Exercise 8. Show that any of the functions

$$F_{j}(z,v) = \frac{M}{(1-z/r)(1-v_{1}/r)\cdots(1-v_{m}/r)},$$

$$F_{j}(z,v) = \frac{M}{(1-z/r)(1-(v_{1}+\ldots+v_{m})/r)},$$

$$F_{j}(z,v) = \frac{M}{1-(z+v_{1}+\ldots+v_{m})/r},$$

$$F_{j}(z,v) = \frac{M}{1-(z/\rho+v_{1}+\ldots+v_{m})/r},$$
with a constant $\rho \in (0,1],$
(27)

majorizes f_j at 0.

Let us consider the system

$$U'_{j} = F_{j}(z, U), \qquad U_{j}(0) = 0, \qquad j = 1, \dots, m,$$
(28)

with F_j being a majorant of f_j . Using the positivity of the coefficients of $Q_{j,k}$, we get

$$|\partial^{\alpha} u_{j}(0)| = |Q_{j,k}(\{\partial^{\beta} f_{\ell}(0)\})| \le Q_{j,k}(\{|\partial^{\beta} f_{\ell}(0)|\}) \le Q_{j,k}(\{\partial^{\beta} F_{\ell}(0)\}) = \partial^{\alpha} U_{j}(0),$$
(29)

i.e., U_j majorizes u_j at 0. To establish analyticity of u_j , it only remains to solve (28) in analytic functions. Given the supply of majorants (27), it is not hard. For example, choosing F_j to be the second function in (27), and putting $U_1 = \ldots = U_m$, we get

$$U_j(z) = \frac{r}{m} \left(1 - \sqrt{1 + 2mM \log\left(1 - \frac{z}{r}\right)} \right),\tag{30}$$

which is analytic at least in the region $|z| < (1 - e^{-\frac{1}{2mM}})r$. We have proved the following.

Theorem 9 (Basic form). Consider the initial value problem

$$u'_{j} = f_{j}(z, u), \qquad u_{j}(0) = 0, \qquad j = 1, \dots, m.$$
 (31)

Let $f_j : \mathbb{C}^{n+m} \to \mathbb{C}$ be analytic at 0, for each j. Then there exists a unique solution u that is analytic at 0.

This theorem can easily be generalized to higher order quasilinear equations, and in a certain sense to any ODE system that can be solved. The most general form of an ODE system for the unknown function $u : \mathbb{C} \to \mathbb{C}^m$ can be written as

$$F_i(z, u, u', u'', \ldots) = 0, \qquad i = 1, \ldots, m.$$
 (32)

Suppose that this can be written in the form

$$u_i^{(q_i)} = f_i(z, u, u', \ldots), \qquad i = 1, \ldots, m,$$
(33)

where for each i and j, the function f_i depends on the derivatives of u_j only up to order $q_j - 1$. In other words, we solve for the highest over derivatives of each component of u. Note that $u_i^{(q_i)}$ in (33) is *not* necessarily the highest order derivative of u_i in (32).

Example 10. Consider the system

$$u''(v''+1) = 0, \qquad u' = v' + z. \tag{34}$$

It looks like the system is second order in both u and v, so that the general solution involves 4 arbitrary constants. But one cannot be sure. Namely, differentiating u' = v' + z gives u'' = v'' + 1, and substituting this into u''(v'' + 1) = 0 we get u'' = 0 and v'' = -1. So the system is equivalent to

$$u'' = 0, \qquad v' = u' - z, \tag{35}$$

which makes it clear that we need only 3 arbitrary constants.

It was observed by Carl Gustav Jacob Jacobi that the key to bringing order to the general system (32) is to somehow write it in the form (33), which was called by him the *normal* form. General methods to do such transformations that often work include using the implicit function theorem (as we have done in Theorem 7), and differentiating the equation (32) with respect to the independent variable. As soon as one brings the system into the normal form, we have local existence.

Theorem 11 (Normal form). Consider the initial value problem

$$u_i^{(q_i)} = f_i(z,\eta), \qquad i = 1, \dots, m, \qquad \eta(0) = \eta_0,$$
(36)

where $\eta \in \mathbb{C}^N$ denotes the collection of derivatives $\{u_i^{(k)} : k = 0, \dots, q_i - 1; i = 1, \dots, m\}$, with $N = \sum_i q_i$. Suppose that for each $i, f_i : \mathbb{C}^{1+N} \to \mathbb{C}^m$ is analytic at $(0, \eta_0)$. Then there exists a unique solution u that is analytic at 0.

In the spirit of Theorem 7, we can also state a fully nonlinear version that relies on the implicit functions theorem.

Theorem 12 (Fully nonlinear form). Consider the initial value problem

$$F(z, \eta, p) = 0, \qquad \eta(0) = \eta_0, \qquad p(0) = p_0,$$
(37)

where $\eta : \mathbb{C} \to \mathbb{C}^N$ denotes the collection $\{u_i^{(k)} : k = 0, \dots, q_i - 1; i = 1, \dots, m\}$, and $p : \mathbb{C} \to \mathbb{C}^m$ denotes the collection $\{u_i^{(q_i)} : i = 1, \dots, m\}$. We suppose that $F : \mathbb{C}^{1+N+m} \to \mathbb{C}^m$ is analytic at $(0, \eta_0, p_0)$, that $F(0, \eta_0, p_0) = 0$, and that the Jacobian matrix $\frac{\partial F}{\partial p}$ is invertible at $(0, \eta_0, p_0)$. Then there exists a unique solution u that is analytic at 0.

Exercise 13. Prove Theorem 11 and Theorem 12.

3. LINEAR PARTIAL DIFFERENTIAL EQUATIONS

Along the same lines, one can establish the local existence of analytic solutions to a very general class of systems of partial differential equations. Such a result was proved by Augustin-Louis Cauchy in 1842 on first order quasilinear evolution equations, and formulated in its most general form by Sofia Vasilyevna Kovalevskaya in 1874. At about the same time, Gaston Darboux also reached similar results, although with less generality than Kovalevskaya's work. Both Kovalevskaya's and Darboux's papers were published in 1875, and the proof was later streamlined by Édouard Jean-Baptiste Goursat in his influential calculus texts around 1900. Nowadays these results are collectively known as the *Cauchy-Kovalevskaya theorem*.

Theorem 14 (Basic form). Consider the Cauchy (or initial value) problem

$$\partial_n u = A_1(z)\partial_1 u + \ldots + A_{n-1}(z)\partial_{n-1} u + A_0(z)u + a(z),$$

$$u|_{\{z_n=0\}} = 0,$$
(38)

where $A_j : \mathbb{C}^n \to \mathbb{C}^{m \times m}$, (j = 0, ..., n - 1), and $a : \mathbb{C}^n \to \mathbb{C}^m$ are analytic at 0. Then there exist an open set $\Omega \subset \mathbb{C}^n$ with $0 \in \Omega$, and an analytic solution $u : \Omega \to \mathbb{C}^m$ that is unique among the functions from $C^{\omega}(\Omega, \mathbb{C}^m)$.

Remark 15. The zero initial condition in (38) is not a restriction, since the Cauchy problem

$$\partial_n v = B_1(z)\partial_1 v + \ldots + B_{n-1}(z)\partial_{n-1}v + B_0(z)v + b(z),$$
(39)

$$v|_{\{z_n=0\}} = \phi|_{\{z_n=0\}},$$

with ϕ analytic near 0, is equivalent to (38) under the substitutions $v = u + \phi$, $A_j = B_j$, (j = 0, ..., n - 1), and

$$a = -\partial_n \phi + B_1 \partial_1 \phi + \ldots + B_{n-1} \partial_{n-1} \phi + B_0 \phi + b.$$

$$\tag{40}$$

Proof of Theorem 14. Since the initial condition is identically zero, we have

$$\partial^{\alpha} u(0) = 0, \quad \text{if} \quad \alpha_n = 0.$$
 (41)

The derivatives $\partial^{\alpha} u$ with $\alpha_n > 0$ can be found by differentiating the equation (38). For example, we have

$$\partial_k \partial_n u = \partial_k a + (\partial_k A_0) u + (\partial_k A_1) \partial_1 u + \dots + (\partial_k A_{n-1}) \partial_{n-1} u + A_0 \partial_k u + A_1 \partial_1 \partial_k u + \dots + A_{n-1} \partial_{n-1} \partial_k u.$$

$$\tag{42}$$

In general, for α with $\alpha_n > 0$, we have

$$\partial^{\alpha} u_j = q_{j,\alpha}(\{\partial^{\beta} A_k\}, \{\partial^{\beta} a\}, \{\partial^{\gamma} u\}), \tag{43}$$

where $q_{i,\alpha}$ is a polynomial with nonnegative coefficients, depending on the individual components of $\partial^{\beta} A_k$ with $|\beta| \leq |\alpha| - 1$ and $k = 0, \ldots, n - 1$, of $\partial^{\beta} a$ with $|\beta| \leq |\alpha| - 1$, and of $\partial^{\gamma} u$ with $|\gamma| \leq |\alpha|$ and $\gamma_n \leq \alpha_n - 1$. Exactly as before, we can eliminate the terms $\partial^{\gamma} u$ and evaluate the result at 0 to get

$$\partial^{\alpha} u_{j}(0) = q_{j,\alpha}(\{\partial^{\beta} A_{k}(0)\}, \{\partial^{\beta} a(0)\}, \{\partial^{\gamma} u(0)\}) = Q_{j,\alpha}(\{\partial^{\beta} A_{k}(0)\}, \{\partial^{\beta} a(0)\}),$$
(44)

where $Q_{j,\alpha}$ is a polynomial with nonnegative coefficients, depending on the individual components of $\partial^{\beta} A_k(0)$ with $|\beta| \leq |\alpha| - 1$ and $k = 0, \ldots, n-1$, and of $\partial^{\beta} a(0)$ with $|\beta| \leq |\alpha| - 1$. Now we consider the system

$$\partial_n U = B_1(z)\partial_1 U + \ldots + B_{n-1}(z)\partial_{n-1} U + B_0(z)U + b(z),$$

$$U|_{\{z_n=0\}} = 0,$$
(45)

with B_j majorizing A_j componentwise at 0 for each $j = 0, \ldots, n-1$, and b majorizing a componentwise at 0. Then for all multi-indices α with $\alpha_n > 0$, we have

$$\begin{aligned} |\partial^{\alpha} u_{j}(0)| &= |Q_{j,\alpha}(\{\partial^{\beta} A_{k}(0)\}, \{\partial^{\beta} a(0)\})| \leq Q_{j,\alpha}(\{|\partial^{\beta} A_{k}(0)|\}, \{|\partial^{\beta} a(0)|\}) \\ &\leq Q_{j,\alpha}(\{\partial^{\beta} B_{k}(0)\}, \{\partial^{\beta} b(0)\}) = \partial^{\alpha} U_{j}(0), \end{aligned}$$
(46)

where it is understood that the absolute values of matrices and vectors are taken componentwise. This shows that U_j majorizes u_j as a function of $z \in \mathbb{C}^n$ at 0.

Since A_k (k = 0, ..., n - 1) and a are componentwise analytic at 0, there exist constants $M < \infty$ and r > 0 such that

$$|\partial^{\beta}A_k(0)|_{\infty} \le M \frac{\beta!}{r^{|\beta|}}, \quad \text{and} \quad |\partial^{\beta}a(0)|_{\infty} \le M \frac{\beta!}{r^{|\beta|}}, \quad \text{for all } \beta,$$
(47)

where $|X|_{\infty}$ denotes the magnitude of the largest component (in magnitude) of X. Then we set each component of B_k (k = 0, ..., n - 1) and of b to be equal to

$$f(z) = \frac{M'}{1 - \frac{z_1 + \dots + z_n}{r}},\tag{48}$$

with $M' \ge M$ to be chosen in a moment. Putting $s = z_1 + \ldots + z_n$ and $v = U_1 = \ldots = U_m$, the equation (45) reduces to

$$v'(s) = mf(s)v(s) + (n-1)mf(s)v'(s) + f(s),$$
(49)

where

$$f(s) = \frac{M'r}{r-s}.$$
(50)

Now we choose $M' = \max\{M, \frac{3}{m(n-1)}\}$, so that (n-1)m|f(s)| > 2 for $|s| < \frac{r}{2}$, and hence

$$v'(s) = \frac{f(s)(mv(s)+1)}{1-(n-1)mf(s)} = \frac{M'r(mv(s)+1)}{(1-(n-1)mM')r-s},$$
(51)

¹This argument was suggested by Mikhail Karpukhin.

has an analytic solution in a neighbourhood of 0, with v(0) = 0. It is also clear that the radius of convergence of the Maclaurin series for v depends only on M, r, m, and n.

It is easy to generalize the preceding theorem to higher order equations, and to solve them near an open subset of the hyperplane $\{z_n = 0\}$.

Theorem 16 (Linear equations). Consider the linear equation

$$Au \equiv \sum_{|\alpha| \le q} a_{\alpha} \partial^{\alpha} u = f, \tag{52}$$

where a_{α} and f are analytic at $0 \in \mathbb{C}^n$, with $a_{(0,\ldots,0,q)}(0) \neq 0$. Furthermore, consider the Cauchy problem of finding a solution to (52) with the prescribed initial values

$$\partial_n^k u(\zeta, 0) = \psi_k(\zeta), \quad \zeta \in \mathbb{C}^{n-1}, \quad k = 0, \dots, q-1,$$
(53)

where the functions ψ_k are analytic at $0 \in \mathbb{C}^{n-1}$. Then there exists a unique solution u to the Cauchy problem that is analytic at $0 \in \mathbb{C}^n$.

Proof. First, let us transform the equation into a first order system amenable to Theorem 14, by introducing the new variables $u_{\alpha} = \partial^{\alpha} u$ for $|\alpha| \leq q - 1$. We need equations for $\partial_n u_{\alpha}$. For $|\alpha| \leq q - 2$, we simply use the definitions

$$\partial_n u_\alpha = u_{(\alpha_1, \dots, \alpha_{n-1}, \alpha_n + 1)}.$$
(54)

(- ·)

For $|\alpha| = q - 1$ with $\alpha_n < q - 1$, we necessarily have $\alpha_\ell > 0$ for some $\ell \neq n$. So we can use the equation

$$\partial_n u_\alpha = \partial_\ell u_{(\dots,\alpha_\ell - 1,\dots,\alpha_n + 1)}.\tag{55}$$

Let us denote the dependence of ℓ on α by $\ell = \ell(\alpha)$. What remains is an equation for $\partial_n u_{(0,\ldots,0,q-1)}$. Of course, this is supplied by (52), which becomes

$$a_{(0,\dots,0,q)}\partial_{n}u_{(0,\dots,0,q-1)} = f - \sum_{|\alpha| \le q-1} a_{\alpha}u_{\alpha} - \sum_{|\alpha|=q} a_{\alpha}\partial_{\ell}u_{(\dots,\alpha_{\ell}-1,\dots)},$$
(56)

where in the last sum, ℓ is chosen, depending on α , so that $\alpha_{\ell} > 0$ and $\ell \neq n$.

The initial condition for u_{α} is obtained from that of $\partial_n^{\alpha_n} u$ by applying the "spatial" differential operator $\partial_1^{\alpha_1} \cdots \partial_{n-1}^{\alpha_{n-1}}$. Hence the equation is reduced to a first order system, and Theorem 14 guarantees a nonempty ball $B \subset \mathbb{C}^n$ centred at 0, such that a unique analytic solution $\{u_{\alpha}\}$ exists on B.

What remains to be shown is $\partial^{\alpha} u_0 = u_{\alpha}$ for $|\alpha| \leq q - 1$, so that $u := u_0$ would satisfy the equation (52) together with the initial conditions (53). First, we observe that (54) implies

$$\partial_n^k u_0 = u_{(0,\dots,0,k)}, \qquad k = 0,\dots, q-1.$$
 (57)

Then, (55) gives $\partial_n u_{\alpha} = \partial_\ell u_{(\dots,\alpha_\ell-1,\dots,\alpha_n+1)} = \partial_n \partial_\ell u_{(\dots,\alpha_\ell-1,\dots,\alpha_n)}$, meaning that

$$\partial_n [u_\alpha - \partial_\ell u_{(\dots,\alpha_\ell - 1,\dots,\alpha_n)}] = 0.$$
(58)

Since $u_{\alpha} = \partial_{\ell} u_{(\dots,\alpha_{\ell}-1,\dots,\alpha_n)}$ at $\{z_n = 0\}$, we have $u_{\alpha,i} \equiv \partial_{\ell} u_{(\dots,\alpha_{\ell}-1,\dots,\alpha_n)}$. This in turn gives us $\partial_n u_{(\dots,\alpha_n-1)} = \partial_n \partial_{\ell} u_{(\dots,\alpha_{\ell}-1,\dots,\alpha_n-1)}$, leading to $u_{(\dots,\alpha_n-1)} \equiv \partial_{\ell} u_{(\dots,\alpha_{\ell}-1,\dots,\alpha_n-1)}$ etc., we get

$$u_{(\alpha_1\dots\alpha_{n-1},k)} \equiv \partial_\ell u_{(\alpha_1\dots\alpha_{\ell-1},\alpha_{\ell}-1,\alpha_{\ell+1}\dots,\alpha_{n-1},k)}, \qquad k = 0,\dots,\alpha_n.$$
⁽⁵⁹⁾

Now, let α be a multi-index with $0 < |\alpha| \le q-1$. If $\alpha_n = |\alpha|$, then we have $\partial^{\alpha} u_0 = u_{\alpha}$ by (57). Suppose that $\alpha_n = 0$. Then there is a multi-index α^* such that $\alpha_k^* = \alpha_k$ for all $k = 1, \ldots, n-1$, and $|\alpha^*| = q-1$. Namely, we have $\alpha_n^* = q-1 - |\alpha|$, which satisfies $0 \le \alpha_n^* < q-1$. Invoking (55) and (59), we conclude that there is an index β with $|\beta| = |\alpha| - 1$ and $\beta_n = 0$, such that $\partial^{\alpha-\beta}u_{\beta} = u_{\alpha}$. If $\beta \ne 0$, then we repeat this procedure for β , and so on, until we reach the multi-index 0. This implies that $\partial^{\alpha}u_0 = u_{\alpha}$ for α with $\alpha_n = 0$. For general α , we combine this result with (57), to get $\partial^{\alpha}u_0 = u_{\alpha}$, completing the proof. In the preceding theorem, apart from the analyticity conditions which are natural, there is the condition on the allowed derivatives of u appearing in the right hand side. The following counterexample due to Kovalevskaya illustrates that this condition is necessary.

Example 17. Consider the heat equation

$$\partial_t u = \partial_r^2 u,\tag{60}$$

to be solved in a neighbourhood of the origin in $(x, t) \in \mathbb{C}^2$, with an analytic initial datum u(x, 0) prescribed on the line $\{t = 0\}$. Differentiating the equation with respect to t gives

$$\partial_t \partial_t u = \partial_t \partial_x^2 u = \partial_x^2 \partial_t u = \partial_x^4 u, \tag{61}$$

and by repeated differentiations, we get

$$\partial_t^k u = \partial_x^{2k} u, \qquad \Rightarrow \qquad \partial_t^k u(0,0) = \partial_x^{2k} u(0,0). \tag{62}$$

The strongest bounds on $\partial_x^{2k}u(0,0)$ for general analytic initial data u(x,0) are of the form $M(2k)!/r^k$. On the other hand, in order for u to be analytic in the *t*-direction, the derivatives $\partial_t^k u(0,0)$ must necessarily have a bound of the form $Ck!/\rho^k$. The moral of the story is that by equating more spatial derivatives on the right hand side with less time derivatives on the left hand side, we generate faster growth in the right hand side than is allowed for the left hand side to be analytic.

Exercise 18. Cook up an initial datum u(x, 0) for the heat equation that is analytic for all x such that the function u(0, t) is not analytic at t = 0.

4. Characteristic surfaces

In this section we discuss how one can adapt the Cauchy-Kovalevskaya theorem if one were to specify Cauchy data on a general analytic surface. Since the theorem concerned is a local result, local considerations will suffice. So locally, an analytic surface is the zero level set of an analytic function. More precisely, $\Gamma \subset \mathbb{C}^n$ is an analytic surface if there is an analytic function $\varphi: U \to \mathbb{C}$ with U an open subset of \mathbb{C}^n , such that $\Gamma = \{z \in U : \varphi(z) = 0\}$ and $\partial \varphi = (\partial_1 \varphi, \dots, \partial_n \varphi)$ is nonzero on Γ . In order to specify Cauchy data on Γ , we assume that there is an analytic injection $w: U \to \mathbb{C}^n$ with $w_n \equiv \varphi$, i.e., that there is an analytic coordinate system (w_1, \ldots, w_n) in a neighbourhood of Γ , that makes $\Gamma = \{w_n = 0\}$. This is always possible locally at any given point $z \in \Gamma$, by shrinking the neighbourhood U if necessary. For example, it suffices to take a rectilinear coordinate system with its n-th axis having the same direction as the normal of Γ at z, and then adjust the n-th coordinate so that Γ becomes $\{w_n = 0\}$. The approach we take in this section is to specify the Cauchy data on Γ in the *w*-coordinate system. Then since in the *w*-coordinates Γ is just $\{w_n = 0\}$, the Cauchy-Kovalevskaya theorem readily applies, provided that the equation can be solved for the term $\partial_{w_n}^q u$. Looking at what this tells us in the original z-coordinates, we will obtain an important insight on what type of initial surfaces the equation "prefers".

For simplicity, let us consider the q-th order linear equation

$$Au \equiv \sum_{|\alpha| \le q} a_{\alpha} \partial^{\alpha} u = f.$$
(63)

Denote by b_{α} the coefficients of the q-th order derivatives in w-coordinates, i.e.,

$$\sum_{|\alpha| \le q} a_{\alpha}(z) \partial_{z}^{\alpha} u = \sum_{|\alpha| \le q} b_{\alpha}(w) \partial_{w}^{\alpha} u.$$
(64)

What is important for us is the particular coefficient b_{α^*} with $\alpha^* = (0, \ldots, 0, q)$, because if, say, $b_{\alpha^*}(w) \neq 0$, we can solve for the term $\partial_{w_n}^q u$ in a neighbourhood of w, and therefore can

apply the Cauchy-Kovalevskaya theorem at w. Considerations such as

$$\frac{\partial u}{\partial z_i} = \frac{\partial u}{\partial w_k} \frac{\partial w_k}{\partial z_i}, \qquad \frac{\partial^2 u}{\partial z_i \partial z_j} = \frac{\partial^2 u}{\partial w_k \partial w_l} \frac{\partial w_k}{\partial z_i} \frac{\partial w_l}{\partial z_j} + \frac{\partial u}{\partial w_k} \frac{\partial^2 w_k}{\partial z_i \partial z_j}, \tag{65}$$

imply that

$$b_{\alpha^*} = \sum_{|\alpha|=q} a_{\alpha} \left(\frac{\partial w_n}{\partial z_1}\right)^{\alpha_1} \cdots \left(\frac{\partial w_n}{\partial z_n}\right)^{\alpha_n}.$$
 (66)

Shifting back to z-coordinates and using the notation $\varphi \equiv w_n$, we see that if

$$\sum_{|\alpha|=q} a_{\alpha} \left(\frac{\partial \varphi}{\partial z_1}\right)^{\alpha_1} \cdots \left(\frac{\partial \varphi}{\partial z_n}\right)^{\alpha_n} \neq 0, \tag{67}$$

at some point $z \in \Gamma$, then the Cauchy problem with initial data on Γ is locally solvable at z. It is a good time to introduce some terminologies.

Definition 19. The function

$$C_A(z,\xi) = \sum_{|\alpha|=q} a_\alpha(z)\xi^\alpha \equiv \sum_{|\alpha|=q} a_\alpha(z)\xi_1^{\alpha_1}\cdots\xi_n^{\alpha_n},$$
(68)

defined for $z \in U$ and $\xi \in \mathbb{C}^n$, is called the *characteristic form* of the operator A.

The characteristic form is a homogeneous function of degree q in ξ , i.e.,

$$C_A(z,\lambda\xi) = \lambda^q C(z,\xi), \qquad \lambda \in \mathbb{C}.$$
(69)

In terms of the characteristic form, the condition (67) becomes

$$C_A(z,\partial_z\varphi(z)) \neq 0. \tag{70}$$

Definition 20. If $C_A(z, \partial_z \varphi(z)) = 0$ for $z \in \Gamma$, then Γ is said to be *characteristic at z* to the equation (63). If Γ is characteristic at each of its points, it is called a *characteristic surface*.

The Cauchy problem for (63) has a unique analytic solution near a hypersurface Γ , if Γ is nowhere characteristic.

Exercise 21. Let Γ be an analytic hypersurface in \mathbb{R}^n , meaning that for each $x \in \Gamma$, there exists a neighbourhood $U \subset \mathbb{R}^n$ of x, and an analytic function $\phi: U \to \mathbb{R}$ with $\nabla \phi \neq 0$ on Γ , such that $\Gamma \cap U = \phi^{-1}(\{0\})$. Let $\Omega \subset \mathbb{R}^n$ be an open set containing Γ , and let $X: \Omega \to \mathbb{R}^n$ be an analytic vector field, satisfying $X(x) \notin T_x\Gamma$ for every $x \in \Gamma$. We consider the linear differential equation

$$Au \equiv \sum_{|\alpha| \le q} a_{\alpha} \partial^{\alpha} u = f, \tag{\dagger}$$

where a_{α} ($|\alpha| \leq q$) and f are analytic functions in Ω , and we assume that Γ is nowhere characteristic for A. Prove that there exists an open set $\omega \subset \mathbb{R}^n$ containing Γ , such that there exists a unique analytic solution $u : \omega \to \mathbb{R}$ to the Cauchy problem for (\dagger) with the Cauchy data

 $X^k u = g_k, \quad k = 1, \dots, q - 1, \qquad \text{on } \Gamma,$

where g_k (k = 1, ..., q-1) are analytic functions in Ω . Note that if we denote the components of X by $X_1, ..., X_n$, the differential operator X^k is understood to be

$$X^k = (X_1\partial_1 + \ldots + X_n\partial_n)^k.$$

Definition 22. We can also introduce the *characteristic cone at* $z \in U$ as

$$\operatorname{Char}_{z} A = \{ \xi : C(z,\xi) = 0 \}.$$
 (71)

Then a surface is characteristic at a point if the normal to the surface at that point belongs to the characteristic cone at the same point. **Example 23.** The characteristic form of the Laplace operator is

$$C_{\Delta}(z,\xi) = \sum_{i=1}^{n} \xi_i^2.$$

There is no nonzero real vector $\xi \in \mathbb{R}^n$ that makes $C_{\Delta}(z,\xi) = 0$, so the generators of the characteristic cones cannot be parallel to any real vector. Let us denote the characteristic cone by $\operatorname{Char}_z \Delta_n$, which is of course independent of z. To reiterate, we have $\operatorname{Char}_z \Delta_n \cap \mathbb{R}^n = \{0\}$, and any real surface cannot be characteristic to the Laplace equation. Equations without real characteristic surfaces are called *elliptic equations*.

Example 24. The cone $\operatorname{Char}_{z}\Delta_{n}$ can easily be described as a whole as an object in \mathbb{C}^{n} , but the most relevant to us is the behaviour of the cone on the hyperplanes \mathbb{R}^{n} and $\mathbb{R}^{n-1} \times i\mathbb{R}$. The former is trivial and has just been discussed. For the latter, it is convenient to make the substitution $z_{n} \mapsto iz_{n}$, called the *Wick rotation*, under which the Laplace equation becomes the wave equation, and the set $\mathbb{R}^{n-1} \times i\mathbb{R}$ becomes \mathbb{R}^{n} . For the wave equation, we have

$$C_{\Box}(z,\xi) = -\xi_n^2 + \sum_{i=1}^{n-1} \xi_i^2.$$

Restricting every variable to the reals, the characteristic cone in this case is called the *light* cone, and any surface whose normal makes an angle $\pi/4$ with the direction of z_n is a characteristic surface.

Example 25. The heat and Schrödinger equations transform into each other by Wick rotations. The both equations have

$$C(z,\xi) = \sum_{i=1}^{n-1} \xi_i^2,$$

as their characteristic form, and the characteristic cone is exactly $\operatorname{Char}_{z}\Delta_{n-1} \times \mathbb{C}$. Restricted to the reals, the characteristic cone is the vertical line $\{\xi : \xi_1 = \ldots = \xi_{n-1} = 0\}$, and so the characteristic surfaces are the horizontal planes $\{x : x_n = \text{const}\}$.

Exercise 26. For each of the following cases, determine the characteristic cones and characteristic surfaces, restricted to the reals.

- a) Tricomi-type equation: $u_{xx} + yu_{yy} = 0$.
- b) Wave equation with wave speed c > 0: $u_{xx} + u_{yy} + u_{zz} = c^{-2}u_{tt}$. How many regions does the characteristic cone divide \mathbb{R}^4 into?
- c) Ultrahyperbolic "wave" equation: $u_{xx} + u_{yy} = u_{zz} + u_{tt}$. How many regions does the characteristic cone divide \mathbb{R}^4 into?
- d) Linear transport equation: $\sum_{i=1}^{n} \alpha_i(x) \partial_i u = 0.$

Example 27. As a prototypical example of what happens when one tries to prescribe initial data on a characteristic surface, let us look at the linear transport equation from part d) of the preceding exercise. There it is found that if Γ is characteristic at $x \in \Gamma$, then the vector $\alpha(x)$ is tangent to Γ . Let us assume that Γ is everywhere characteristic. Then all our transport equation tells us is the behaviour of u along Γ , and what u does in the transversal direction is completely "free". This means that the existence is lost unless the initial condition on Γ satisfies certain constraints, and if a solution exists, it will not be unique. The situation is entirely analogous to solving the linear system Ax = b with a non-invertible square matrix A. Now let us forget about specifying initial conditions and take a slightly different point of view. Imagine that the graphs of several solutions to the transport equation are drawn in \mathbb{R}^{n+1} , and imagine also several surfaces in \mathbb{R}^n , which is to be understood as the base of the space \mathbb{R}^{n+1} in which the graphs live. Then we see that the characteristic surfaces are the only

surfaces along which two different solutions can touch each other, for if two solutions are the same on a non-characteristic surface, by uniqueness they must coincide in a neighbourhood of the surface.

5. Well posedness

The complex analytic setting is completely natural for the Cauchy-Kovalevskaya theorem. This is because any real analytic function uniquely extends to a complex analytic one in a neighbourhood of \mathbb{R}^n considered as a subset of \mathbb{C}^n , and more importantly this point of view offers a better insight on the behaviour of analytic functions. Hence the complex analytic treatment contains the real analytic case as a special case. However, it is known that if we allowed only analytic solutions, we would be missing out on most of the interesting properties of partial differential equations. For instance, since analytic functions are completely determined by its values on any open set however small, it would be extremely cumbersome, if not impossible, to describe phenomena like wave propagation, in which initial data on a region of the initial surface are supposed to influence only a specific part of space-time. A much more natural setting for a differential equation would be to require its solutions to have just enough regularity for the equation to make sense. For example, the Laplace equation $\Delta u = 0$ already makes sense for twice differentiable functions. Actually, the solutions to the Laplace equation, i.e., harmonic functions, are automatically analytic, which has a deep mathematical reason that could not be revealed if we restricted ourselves to analytic solutions from the beginning. In fact, the solutions to the Cauchy-Riemann equations, i.e., holomorphic functions, are analytic by the same underlying reason, and complex analytic functions are nothing but functions satisfying the Cauchy-Riemann equations. From this point of view, looking for analytic solutions to a PDE in \mathbb{R}^n would mean coupling the PDE with the Cauchy-Riemann equations and solving them simultaneously in \mathbb{R}^{2n} . In other words, if we are not assuming analyticity, \mathbb{C}^n is better thought of as \mathbb{R}^{2n} with an additional algebraic structure. Hence the real case is more general than the complex one, and from now on, we will be working explicitly in *real* spaces such as \mathbb{R}^n , unless indicated otherwise.

As soon as we allow non-analytic data and/or solutions, many interesting questions arise surrounding the Cauchy-Kovalevskaya theorem. First, assuming a setting to which the Cauchy-Kovalevskaya theorem can be applied, we can ask if there exists any (necessarily non-analytic) solution other than the solution given by the Cauchy-Kovalevskaya theorem. In other words, is the uniqueness part of the Cauchy-Kovalevskaya theorem still valid if we now allow non-analytic solutions? For linear equations an affirmative answer is given by *Holm*gren's uniqueness theorem. Moreover, uniqueness holds for first order equations, but fails in general for higher order equations and systems. Such a uniqueness result can also be thought of as a regularity theorem, in the sense that if u is a solution then it would be automatically analytic by uniqueness.

The second question is whether existence holds for non-analytic data, and again the answer is negative in general. A large class of counterexamples can be constructed, by using the fact that some equations, such as the Laplace and the Cauchy-Riemann equations, have only analytic solutions, therefore their initial data, as restrictions of the solutions to an analytic hypersurface, cannot be non-analytic. Hence such equations with non-analytic initial data do not have solutions. In some cases, this can be interpreted as one having "too many" initial conditions that make the problem overdetermined, since in those cases the situation can be remedied by removing some of the initial conditions. For example, with sufficiently regular closed surfaces as initial surfaces, one can remove either one of the two Cauchy data in the Laplace equation, arriving at the Dirichlet or Neumann problem. Starting with Hans Lewy's celebrated counterexample of 1957, more complicated constructions along similar lines have been made that ensure the inhomogeneous part of a linear equation to be analytic, thus exhibiting examples of linear equations with no solutions when the inhomogeneous part is nonanalytic, regardless of initial data. The lesson to be learned from these examples is that the existence theory in a non-analytic setting is much more complicated than the corresponding analytic theory, and in particular one has to carefully decide on what would constitute the initial data for the particular equation.

Indeed, there is an illuminating way to detect the poor behaviour of some equations discussed in the previous paragraph with regard to the Cauchy problem, entirely from within the analytic setting, that runs as follows. Suppose that in the analytic setting, for a generic initial datum ψ it is associated the solution $u = S(\psi)$ of the equation under consideration, where $S: \psi \mapsto u$ is the solution map. Now suppose that the datum ψ is non-analytic, say, only continuous. Then by the Weierstrass approximation theorem, for any $\varepsilon > 0$ there is a polynomial ψ_{ε} that is within an ε distance from ψ . Taking some sequence $\varepsilon \to 0$, if the solutions $u_{\varepsilon} = S(\psi_{\varepsilon})$ converge locally uniformly to a function u, we could reasonably argue that u is a solution (in a generalized sense) of our equation with the (non-analytic) datum ψ . The counterexamples from the preceding paragraph suggest that in those cases the sequence u_{ε} cannot converge. Actually, the situation is much worse, as the following example due to Jacques Salomon Hadamard shows.

Example 28. Consider the Cauchy problem for the Laplace equation

$$u_{tt} + u_{xx} = 0, \qquad u(x,0) = a_{\nu} \sin \nu x, \qquad u_t(x,0) = b_{\nu} \sin \nu x, \tag{72}$$

whose solution is given by

$$u(x,t) = (a_{\nu} \cosh \nu t + \frac{b_{\nu}}{\nu} \sinh \nu t) \sin \nu x.$$
(73)

Choosing, e.g., $a_{\nu} = 1/\nu$ and $b_{\nu} = 0$ with ν large, we see that the solution grows arbitrarily fast, although the initial data are arbitrarily small. In a certain sense, the relation between the solution and the Cauchy data becomes more and more difficult to invert as we add higher and higher frequencies. The initial data could be, for instance, the error of an approximation of non-analytic initial data in the uniform norm, with the approximation getting better as $\nu \to \infty$. Then the solutions with initial data given by the approximations diverge unless a_{ν} and b_{ν} decay faster than exponential. But functions that can be approximated by analytic functions with such small errors form a severely restricted class, being between the smooth functions C^{∞} and the analytic functions C^{ω} .

Exercise 29. Consider the problem

$$u_{tt} + u_{xx} = 0,$$
 $u(x,0) = \phi(x),$ $u_t(x,0) = \psi(x).$

For given $\varepsilon > 0$ and an integer k > 0, construct initial data ϕ and ψ such that

$$\|\phi\|_{\infty} + \ldots + \|\phi^{(k)}\|_{\infty} + \|\psi\|_{\infty} + \ldots + \|\psi^{(k)}\|_{\infty} < \varepsilon,$$

and

$$||u(\cdot,\varepsilon)||_{\infty} > \frac{1}{\varepsilon}.$$

Let us contrast the preceding example with the following.

Example 30. Consider the Cauchy problem for the wave equation

$$u_{tt} - u_{xx} = 0,$$
 $u(x,0) = \phi(x),$ $u_t(x,0) = \psi(x),$ (74)

whose solution is given by the d'Alambert formula

$$u(x,t) = \frac{\phi(x-t) + \phi(x+t)}{2} + \frac{1}{2} \int_{x-t}^{x+t} \psi(y) \,\mathrm{d}y, \tag{75}$$

where if t < 0, the integral over (x - t, x + t) is understood to be the minus of the integral over (x + t, x - t). From this, it is easy to deduce the bound

$$|u(x,t)| \le \sup_{y \in [x-t,x+t]} |\phi(y)| + |t| \sup_{y \in [x-t,x+t]} |\psi(y)|,$$
(76)

making it clear that small initial data lead to small solutions. Moreover, one can show uniqueness by an energy argument, meaning that the solution given by the d'Alambert formula is the only one.

Triggered by considerations such as the preceding ones, Hadamard introduced the concept of *well-posedness* of a problem. To define this concept abstractly, we assume a set D, that represents all possible *data* in the problem, a second set S, that represents all possible *solutions*, and finally a *relation* $R(f, u) \in \{0, ...\}$, defined for $f \in D$ and $u \in S$. Then we consider the following *problem*: Given $f \in D$, find $u \in S$ such that R(f, u) = 0. This problem is said to be *well-posed* if

- For any $f \in D$ there exists a unique solution $u \in S$, and
- Varying f a bit results in a small variation of u, i.e., u depends on f continuously.

In order make the second point precise, we need to define what we mean by continuity of maps $\sigma: D \to S$, i.e., we need to choose topologies for the sets D and S. We can introduce some flexibility on the choice of the sets D and S too, leading to the meta-problem: Find "reasonable" topological spaces D and S such that the problem R(f, u) = 0 with $f \in D$ and $u \in S$ is well-posed. Usually, the "correct" topologies on D and S are suggested by the structure of the problem itself, or what is essentially the same, by the real world or mathematical phenomenon the problem is supposed to model. The concept of well-posedness has proved to be very useful in revealing the true nature of the equations, especially in identifying the "correct" initial and/or boundary conditions. Of course, one important motivation of the well-posedness concept is that in modelling of real world phenomena, the problem data always have some measurement or computational error in it, so without well-posedness, we cannot say that the solution corresponding to imprecise data is anywhere near the solution we are trying to capture. Thus, a necessary condition for a physics theory to have any predictive power is that it must produce well-posed problems. One might wonder if a counterexample to this statement can be exhibited by mentioning the fact that in practice people routinely solve what are normally considered as *ill-posed* problems, i.e., problems that are not well-posed. However, in those situations "solving a problem" has a broader meaning, and as part of this process one replaces the original ill-posed problem by a well-posed one, with the aid of a reqularization procedure. For example, from Hadamard's example we have seen that essentially the "trouble makers" are initial data that oscillate rapidly in space, and a bit more analysis shows that if the initial data has frequencies not exceeding ν , then the Cauchy problem can be solved without trouble for time of order $1/\nu$. This offers a good theory provided that in the particular situation under consideration, we know for sure there will not be frequencies higher than ν present in any realistic initial data, and we do not need to solve the Cauchy problem for time intervals much longer than $1/\nu$.

6. The general Cauchy-Kovalevskaya theorem

For completeness, we include here a proof of the Cauchy-Kovalevskaya theorem for general (nonlinear) partial differential equations. The following is the most basic form of the theorem, from which all more general forms can be deduced.

Theorem 31 (Basic form). Consider the Cauchy (or initial value) problem

$$\partial_n u_j = f_j(z, u, \partial_1 u, \dots, \partial_{n-1} u), \qquad u_j(\zeta, 0) = 0, \quad \zeta \in \mathbb{C}^{n-1}, \qquad j = 1, \dots, m.$$
(77)

Let $f_j : \mathbb{C}^{n+m+(n-1)\times m} \to \mathbb{C}$ be analytic at 0, for all j. Then there exists a unique solution u that is analytic at 0.

Proof. Without loss of generality, we can assume that $f_j(0) = 0$ for all j, by replacing $u_j(z)$ with $u_j(z) - z_n f_j(0)$. Also, it will be convenient to label by p_{ik} the slot of f_j that takes $\partial_i u_k$ as its argument, i.e., $f_j = f_j(z, u, p)$ with $z \in \mathbb{C}^n$, $u \in \mathbb{C}^m$, and $p \in \mathbb{C}^{(n-1) \times m}$. Since the initial condition is identically zero, we have

$$\partial^{\alpha} u(0) = 0, \qquad \text{if} \quad \alpha_n = 0. \tag{78}$$

The derivatives $\partial^{\alpha} u$ with $\alpha_n > 0$ can be found by differentiating the equation (77). For example, we have

$$\partial_k \partial_n u_j = \frac{\partial f_j}{\partial z_k} + \frac{\partial f_j}{\partial u_q} \frac{\partial u_q}{\partial z_k} + \frac{\partial f_j}{\partial p_{iq}} \frac{\partial^2 u_q}{\partial z_i \partial z_k},$$

where implicit summations over q = 1, ..., m, and i = 1, ..., n - 1, are assumed in the terms they appear. In general, for α with $\alpha_n > 0$, we have

$$\partial^{\alpha} u_j = q_{\alpha} (\partial^{\beta} f_j, \partial^{\gamma} u), \tag{79}$$

where q_{α} is a polynomial with nonnegative coefficients, depending on $\partial^{\beta} f_{j}$ with $|\beta| \leq |\alpha| - 1$, and $\partial^{\gamma} u$ with $|\gamma| \leq |\alpha|$ and $\gamma_{n} \leq \alpha_{n} - 1$. Exactly as before, we can eliminate the terms $\partial^{\gamma} u$ and evaluate the result at 0 to get

$$\partial^{\alpha} u_j(0) = q_{\alpha}(\partial^{\beta} f_j(0), \partial^{\gamma} u(0)) = Q_{j,\alpha}(\partial^{\beta} f(0)), \tag{80}$$

where $Q_{j,\alpha}$ is a polynomial with nonnegative coefficients, depending on $\partial^{\beta} f_j$ with $|\beta| \leq |\alpha| - 1$. Now it is time to consider the system

$$\partial_n U_j = F_j(z, U, \partial_1 U, \dots, \partial_{n-1} U), \qquad j = 1, \dots, m,$$
(81)

with F_j majorizing f_j at 0 for each j. Then for all multi-indices α with $\alpha_n > 0$, we have

$$|\partial^{\alpha} u_{j}(0)| = |Q_{j,\alpha}(\partial^{\beta} f(0))| \le Q_{j,\alpha}(|\partial^{\beta} f(0)|) \le Q_{j,\alpha}(\partial^{\beta} F(0)) = \partial^{\alpha} U_{j}(0).$$
(82)

If in addition, $U_j|_{z_n=0}$ majorizes 0 as a function of $(z_1, \ldots, z_{n-1}) \in \mathbb{C}^{n-1}$ at 0, then U_j majorizes u_j as a function of $z \in \mathbb{C}^n$ at 0.

Supposing that f_j satisfies the bound (26), let us try the following majorant of f_j .

$$F_j(z, u, p) = \frac{M}{\left(1 - \frac{z_1 + \dots + z_{n-1} + z_n/\rho + u_1 + \dots + u_m}{r}\right) \left(1 - \frac{1}{r} \sum_{i,k} p_{ik}\right)} - M,$$
(83)

where $\rho \in (0, 1]$ is a constant whose value is to be adjusted later. Put $s = z_1 + \ldots + z_{n-1}$, $t = z_n$, and $v := U_1 = \ldots = U_m$, to get

$$\partial_t v = \frac{M}{\left(1 - \frac{s + t/\rho + mv}{r}\right) \left(1 - \frac{(n-1)m}{r} \partial_s v\right)} - M.$$
(84)

Defining the new variable $\sigma = t + \rho s$, and assuming v depends only on σ , this becomes

$$\partial_{\sigma} v = \frac{M}{\left(1 - \frac{\sigma/\rho + mv}{r}\right) \left(1 - \frac{(n-1)m\rho}{r} \partial_{\sigma} v\right)} - M,\tag{85}$$

or, after rearranging

$$\left(1 - \frac{(n-1)mM\rho}{r}\right)\partial_{\sigma}v - \frac{(n-1)m\rho}{r}(\partial_{\sigma}v)^2 = \frac{M}{1 - \frac{\sigma/\rho + mv}{r}} - M.$$
(86)

We choose $\rho \in (0,1]$ so small that $1 - \frac{(n-1)mM\rho}{r} > 0$. Then the preceding equation can be solved for $\partial_{\sigma} v$, in the power series

$$\partial_{\sigma}v = c_1(\sigma/\rho + mv) + c_2(\sigma/\rho + mv)^2 + \dots,$$
(87)

convergent for some $\sigma/\rho + mv \neq 0$, with all coefficients nonnegative: $c_k \geq 0$. In other words, there is a function g analytic at 0, with nonnegative Maclaurin series coefficients and with g(0) = 0, such that

$$\partial_{\sigma} v = g(\sigma/\rho + mv). \tag{88}$$

Now we can apply Cauchy's theorem for analytic ODEs from the previous section to infer that the above equation has a solution v analytic at 0, satisfying v(0) = 0, whose Maclaurin series coefficients are nonnegative. Rewinding everything, the vector function U with components

$$U_j(z) = v(\rho(z_1 + \ldots + z_{n-1}) + z_n),$$
(89)

solves (81). Since the Maclaurin series coefficients of v are nonnegative, the same holds for U_j , implying that $U_j|_{z_n=0}$ majorizes 0 at 0. This establishes the proof.

Remark 32. The choice of our majorant in this proof is one introduced by Goursat, that allows us to treat the right hand side f directly. We refer to Folland's book for an alternative proof that is based on transforming the equations into a quasilinear form.

It is easy to generalize the preceding theorem to higher order equations, and to solve them near an open subset of the hyperplane $\{z_n = 0\}$.

Theorem 33 (Standard form). Consider the equations

$$\partial_n^{q_i} u_i = f_i(z, \{\partial^\beta u\}), \qquad i = 1, \dots, m, \tag{90}$$

where for each i and j, the function f_i depends on the derivatives of u_j only up to order q_j , is independent of $\partial_n^{q_j} u_j$, and analytic in all of its arguments. Furthermore, consider the Cauchy problem of finding a solution to (90) with the prescribed initial values

$$\partial_n^k u_i(\zeta, 0) = \psi_{i,k}(\zeta), \quad \zeta \in \Omega, \quad k = 0, \dots, q_i - 1, \qquad i = 1, \dots, m, \tag{91}$$

where $\Omega \subset \mathbb{C}^{n-1}$ is open, and all $\psi_{i,k}$ are analytic on Ω . Then there exists a unique solution u to the Cauchy problem that is analytic in an open set of \mathbb{C}^n containing $\Omega \times \{0\}$.

Exercise 34. Prove this theorem.

Example 35. Let us consider the Cauchy problem

$$(u_t)^2 - u_x = 0, \qquad u(x,0) = \psi(x),$$
(92)

to be solved near the origin in $(x,t) \in \mathbb{C}^2$. This is of course a variation of Example 6. If there is an analytic solution to this problem, in a small neighbourhood of the origin, u_x will be close to $\psi_x(0)$. Therefore, we can specify a branch of $\eta^{\frac{1}{2}}$ in $u_t = (u_x)^{\frac{1}{2}}$ by choosing the value of this branch at $\psi_x(0)$. In other words, instead of (92), we consider the problem

$$(u_t)^2 - u_x = 0, \qquad u(x,0) = \psi(x), \qquad u_t(0) = p_0,$$
(93)

where we assume that $\psi_x(0) \neq 0$, and that $p_0 \in \mathbb{C}$ is a given number satisfying $p_0^2 = \psi_x(0)$. Since there is a unique analytic function f defined in a neighbourhood U of $\psi_x(0)$ satisfying $f(\eta)^2 = \eta$ for $\eta \in U$ and $f(\psi_x(0)) = p_0$, Theorem 31 guarantees a unique analytic solution in a neighbourhood of 0.

Theorem 36 (Fully nonlinear form). Consider the Cauchy problem

$$F(z,\eta,p) = 0, \qquad \eta(\zeta,0) = \eta_0(\zeta), \qquad p(0) = p_0,$$
(94)

where $\eta : \mathbb{C}^n \to \mathbb{C}^N$ denotes the collection $\{\partial^{\alpha} u_i : |\alpha| \leq q_i - 1; \alpha_n \leq q_i - 1; i = 1, ..., m\}$, and $p : \mathbb{C}^n \to \mathbb{C}^m$ denotes the collection $\{\partial_n^{q_i} u_i : i = 1, ..., m\}$. Suppose that $F : \mathbb{C}^{n+N+m} \to \mathbb{C}^m$ is

analytic at $(0, \eta_0(0), p_0)$, that $F(0, \eta_0(0), p_0) = 0$, and that the Jacobian matrix $\frac{\partial F}{\partial p}$ is invertible at $(0, \eta_0(0), p_0)$. Then there exists a unique solution u that is analytic at 0.

Exercise 37. Prove this theorem.

Remark 38. Let us note that the notion of characteristic surfaces can be extended without much difficulty to nonlinear equations, where whether or not a surface is characteristic now may depend on what function we plug into the differential operator, just as in the case of first order nonlinear equations. For more details we recommend Jeffrey Rauch's book.

APPENDIX A. MULTIVARIATE POWER SERIES

Let $\mathbb{K} = \mathbb{R}$ or $\mathbb{K} = \mathbb{C}$, and depending on this choice, let X be a real or complex Banach space. In what follows, assuming $X = \mathbb{K}$ would not lose generality, and might simplify reading. An *n*-variable power series with values in X is an expression of the form

$$f(z) = \sum_{\alpha_1=0}^{\infty} \dots \sum_{\alpha_n=0}^{\infty} a_{\alpha_1,\dots,\alpha_n} (z_1 - c_1)^{\alpha_1} \dots (z_n - c_n)^{\alpha_n},$$
(95)

with the coefficients $a_{\alpha_1,\ldots,\alpha_n} \in X$, and the centre $c \in \mathbb{K}^n$. Introducing the *multi-index* $\alpha = (\alpha_1,\ldots,\alpha_n) \in \mathbb{N}_0^n$, and the convention $z^{\alpha} = z_1^{\alpha_1} \cdots z_n^{\alpha_n}$ for $z \in \mathbb{K}^n$, this series can also be written as

$$f(z) = \sum_{\alpha} a_{\alpha} (z - c)^{\alpha}.$$
(96)

We talk about real power series in case $\mathbb{K} = \mathbb{R}$, and complex power series in case $\mathbb{K} = \mathbb{C}$.

If the preceding series converges for some z and for some linear ordering of the multi-indices α , then obviously there is a constant $M < \infty$, such that $||a_{\alpha}|| |z_1 - c_1|^{\alpha_1} \cdots |z_n - c_n|^{\alpha_n} \leq M$ for all α , where $||a_{\alpha}||$ is the Banach space norm of a_{α} . In particular, if this series converges in a neighbourhood of c, then taking z with $z_1 - c_1 = \ldots = z_n - c_n$, we infer the existence of constants $M < \infty$ and r > 0, such that $||a_{\alpha}|| \leq Mr^{-|\alpha|}$ for all α , where $|\alpha| = \alpha_1 + \ldots + \alpha_n$. On the other hand, if we have $||a_{\alpha}|| \leq Mr^{-|\alpha|}$ for all α , then the series converges absolutely and uniformly in each compact subset of the polydisk $D_r(c) = \{z \in \mathbb{K}^n : |z_i - c_i| < r, \forall i\}$. Note that if $\mathbb{K} = \mathbb{R}$ then the polydisk $D_r(c)$ is simply a cube centred at c. Let us clarify what we mean by absolute and uniform convergence of the series (96).

Definition 39. Given a set $K \subset \mathbb{K}^n$, and a countable family $\{f_i : i \in I\}$ of functions $f_i : K \to X$, the series

$$\sum_{i \in I} f_i, \tag{97}$$

is said to converge absolutely and uniformly in K if there is a constant $M < \infty$ such that

$$\sum_{i\in J} \|f_i\|_K \le M,\tag{98}$$

for any finite subset $J \subset I$, where

$$\|f_i\|_K = \sup_{z \in K} \|f_i(z)\|.$$
(99)

For the power series (96), the family of functions is $\{f_{\alpha}\}$, where $f_{\alpha}(z) = a_{\alpha}(z-c)^{\alpha}$.

Remark 40. We use the term absolute convergence to remind ourselves the fact that we are not far from the basic cases $X = \mathbb{R}$ and $X = \mathbb{C}$. However, especially when X is infinite dimensional, the term normal convergence may be more appropriate.

Exercise 41. Show that if $||a_{\alpha}|| \leq Mr^{-|\alpha|}$ for all α , then the series (96) converges absolutely and uniformly in each compact subset of $D_r(c)$.

Exercise 42. Show that if the series (96) converges absolutely and uniformly in a compact set $K \subset \mathbb{K}^n$, then the partial sums

$$f_k(z) = \sum_{|\alpha| \le k} a_\alpha (z - c)^\alpha, \tag{100}$$

converge uniformly to some $f \in C(K, X)$, that is, $||f - f_k||_K \to 0$ as $k \to \infty$.

In what follows, we shall see that if a power series converges absolutely and uniformly then summing the series in any reasonable way leads to the same result.

Theorem 43. Let $f_{k,\ell} : K \to X$ for $k, l \in \mathbb{N}$ with some $K \subset \mathbb{K}^n$, and let $\sigma : \mathbb{N}^2 \to \mathbb{N}$ be a bijection. Define the sequence $\{g_m\}$ by $g_{\sigma(k,\ell)} = f_{k,\ell}$. Then the followings are equivalent:

- (a) The series $\sum_{m} g_{m}$ converges absolutely and uniformly in K.
- (b) The series $\sum_{k} \sum_{\ell} \|f_{k,\ell}\|_{K}$ converges. In particular, for each $k \in \mathbb{N}$, the series $\sum_{\ell} f_{k,\ell}$ is absolutely and uniformly convergent in K.
- (c) The series $\sum_{\ell} (\sum_{k} ||f_{k,\ell}||_K)$ converges. In particular, for each $\ell \in \mathbb{N}$, the series $\sum_{k} f_{k,\ell}$ is absolutely and uniformly convergent in K.

If any (so all) of the above conditions is satisfied, then we have

$$\sum_{\ell} \left(\sum_{k} f_{k,\ell} \right) = \sum_{k} \left(\sum_{\ell} f_{k,\ell} \right) = \sum_{m} g_{m}.$$

Proof. First we prove the implication (a) \Rightarrow (b). Let $N = \sum_n \|g_n\|_K < \infty$. This obviously implies that for each $k \in \mathbb{N}$, $M_k = \sum_{\ell} \|f_{k,\ell}\|_K < \infty$. Let $\varepsilon > 0$ and let m_k be such that $\sum_{\ell > m_k} \|f_{k,\ell}\|_K \le 2^{-k}\varepsilon$. So for any m we have

$$\sum_{k \le m} (\sum_{\ell} \|f_{k,\ell}\|_K) \le \sum_{k \le m} (\sum_{\ell \le m_k} \|f_{k,\ell}\|_K) + 2\varepsilon \le N + 2\varepsilon.$$

Now we shall prove that $g = \sum_{p} g_{p}$ is equal to $f = \sum_{k} (\sum_{\ell} f_{k,\ell})$. To this end, let m be such that $\sum_{k>m} (\sum_{\ell} ||f_{k,\ell}||_{K}) \leq \varepsilon$, and let $\tilde{f}_{\varepsilon} = \sum_{k \leq m} \sum_{\ell \leq m_{k}} f_{k,\ell}$. Then we have

$$\|f - \tilde{f}_{\varepsilon}\|_{K} \leq \sum_{k>m} (\sum_{\ell} \|f_{k,\ell}\|_{K}) + \sum_{k\leq m} (\sum_{\ell>m_{k}} \|f_{k,\ell}\|_{K}) \leq 3\varepsilon.$$

Similarly, for sufficiently large p, the partial sum $\tilde{g}_p = \sum_{q \leq p} g_q$ satisfies

$$\|\tilde{g}_p - \tilde{f}_{\varepsilon}\|_K \le \sum_{k>m} (\sum_{\ell} \|f_{k,\ell}\|_K) + \sum_{k\le m} (\sum_{\ell>m_k} \|f_{k,\ell}\|_K) \le 3\varepsilon,$$

and so we have

$$\|f-g\|_K \le \|g-\tilde{g}_p\|_K + 6\varepsilon.$$

Since $\tilde{g}_p \to g$ and ε is arbitrary, we conclude f = g.

For the other direction (b) \Rightarrow (a), we start with $M_k = \sum_{\ell} ||f_{k,\ell}||_K < \infty$ and the condition $M = \sum_k M_k < \infty$. Then for any p we have

$$\sum_{q \le p} \|g_q\|_K \le \sum_{k \le m} \sum_{\ell \le m} \|f_{k,\ell}\|_K \le M,$$

where *m* is such that $\{\ell \leq m\}^2 \supset \sigma^{-1}(\{q \leq p\})$. This completes the proof, since the equivalence of (a) and (c) follows by considering $h_{k,\ell} = f_{\ell,k}$.

By induction on the number of nested sums, we get the following result.

Corollary 44. Let $f_{\alpha} : K \to X$ for $\alpha \in \mathbb{N}_0^n$ with some $K \subset \mathbb{K}^n$, and let $\sigma : \mathbb{N}_0^n \to \mathbb{N}$ be a bijection. Define the sequence $\{g_m\}$ by $g_{\sigma(\alpha)} = f_{\alpha}$. Then the followings are equivalent: (a) The series $\sum_m g_m$ converges absolutely and uniformly in K. (b) The series $\sum_{\alpha_1} (\sum_{\alpha_2} \dots (\sum_{\alpha_n} ||f_{\alpha}||_K) \dots)$ converges. If any (so both) of the above conditions is satisfied, then we have

$$\sum_{\alpha_1} \left(\sum_{\alpha_2} \dots \left(\sum_{\alpha_n} f_\alpha \right) \dots \right) = \sum_m g_m.$$

Theorem 43 implies the following very strong result on rearrangements of absolutely uniformly convergent series.

Corollary 45. Suppose that the series $\sum_{p} g_p$ converges to g absolutely and uniformly. Let $\mathbb{N} = \bigcup_k M_k$ be a disjoint decomposition of \mathbb{N} , where k is from a countable set and each $M_k = \{m_{k,\ell}\}$ is countable. Then for each k, the series $\tilde{g}_k = \sum_{\ell} g_{m_{k,\ell}}$ converges absolutely and uniformly, and moreover the series $\sum_k \tilde{g}_k$ converges to g absolutely and uniformly.

Proof. Let us say k runs over $N \subset \mathbb{N}$, and for each $k \in N$, ℓ runs over $L_k \subset \mathbb{N}$. Define the set $M \subset \mathbb{N}^2$ by $M = \{(k,\ell) : k \in N \text{ and } \ell \in L_k\}$, and let $f_{k,\ell} = g_{m_{k,\ell}}$ if $(k,\ell) \in M$, and $f_{k,\ell} = 0$ otherwise. Then we get the proof by applying Theorem 43 with, e.g., $\sigma(k,\ell) = 2m_{k,\ell}$ for $(k,\ell) \in M$, and $\sigma(k,\ell) = 2\tau(k,\ell) - 1$ for $(k,\ell) \in \mathbb{N}^2 \setminus M$, where $\tau : \mathbb{N}^2 \setminus M \to \mathbb{N}$ is a bijection.

By choosing each M_k to have a single element, we get the following.

Corollary 46. Suppose that the series $\sum_{p} g_p$ converges absolutely and uniformly to g. Then given any bijection $\tau : \mathbb{N} \to \mathbb{N}$, the rearranged series $\sum_{p} g_{\tau(p)}$ also converges absolutely and uniformly to g.

APPENDIX B. ANALYTIC FUNCTIONS

We keep the setting of the previous section intact. In particular, we have $\mathbb{K} = \mathbb{R}$ or $\mathbb{K} = \mathbb{C}$, and X is a real or complex Banach space.

Definition 47. Let Ω be an open subset of \mathbb{K}^n . A function $f: \Omega \to X$ is called *analytic at* $c \in \Omega$ if it is developable into a power series around c, i.e., if there are coefficients $a_{\alpha} \in X$, $(\alpha \in \mathbb{N}^n_0)$, such that the power series (96) converges in a neighbourhood of c. Moreover, f is said to be *analytic in* Ω if it is analytic at each $c \in \Omega$. The set of analytic functions in Ω with values in X is denoted by $C^{\omega}(\Omega, X)$.

The following lemma shows that a convergent power series defines an analytic function in a neighbourhood of its centre.

Lemma 48. Suppose that the power series $f(z) = \sum a_{\alpha}(z-c)^{\alpha}$ converges absolutely uniformly in each compact subset of $D_r(c)$ for some r > 0, and let $d \in D_r(c)$. Then we have

$$f(z) = \sum_{\beta} \left(\sum_{\alpha \ge \beta} \binom{\alpha}{\beta} a_{\alpha} (d-c)^{\alpha-\beta} \right) (z-d)^{\beta},$$

which converges absolutely uniformly in each compact subset of $D_{r-|d-c|}(d)$. Here $\alpha \geq \beta$ means that $\alpha_i \geq \beta_i$ for each *i*, and $\binom{\alpha}{\beta} = \binom{\alpha_1}{\beta_1} \cdots \binom{\alpha_n}{\beta_n}$.

Proof. We have

$$(z-c)^{\alpha} = (z-d+d-c)^{\alpha} = \sum_{\beta \le \alpha} {\alpha \choose \beta} (z-d)^{\beta} (d-c)^{\alpha-\beta},$$

so that the proof is established upon justifying

$$\sum_{\alpha} a_{\alpha} \sum_{\beta \le \alpha} \binom{\alpha}{\beta} (z-d)^{\beta} (d-c)^{\alpha-\beta} = \sum_{\beta} \sum_{\alpha \ge \beta} \binom{\alpha}{\beta} a_{\alpha} (z-d)^{\beta} (d-c)^{\alpha-\beta},$$

for $z \in D_R(d)$ with R = r - |d - c|. This can be done by applying Corollary 45 if the left hand side converges absolutely uniformly in each compact subset of $D_R(d)$. To this end, let $|z - d| \leq \rho - |d - c|$ with $\rho < r$. Then we have

$$\sum_{\beta \le \alpha} \binom{\alpha}{\beta} |z_1 - d_1|^{\beta_1} \cdots |z_n - d_n|^{\beta_n} |d_1 - c_1|^{\alpha_1 - \beta_1} \cdots |d_n - c_n|^{\alpha_n - \beta_n} = (|z_1 - d_1| + |d_1 - c_1|)^{\alpha_1} \cdots (|z_n - d_n| + |d_n - c_n|)^{\alpha_n} \le \rho^{|\alpha|},$$

and since $a_{\alpha}\rho^{|\alpha|} = a_{\alpha}r^{|\alpha|}(\rho/r)^{|\alpha|}$ we obtain the desired absolute uniform convergence. \Box

Now we turn to the question of termwise differentiating power series.

Theorem 49. Suppose that the power series $f(z) = \sum a_{\alpha}(z-c)^{\alpha}$ converges absolutely uniformly in each compact subset of $D_r(c)$ for some r > 0. Then the series

$$g(z) = \sum_{\alpha} \alpha_k a_{\alpha} (z_1 - c_1)^{\alpha_1} \cdots (z_k - c_k)^{\alpha_k - 1} \cdots (z_n - c_n)^{\alpha_n},$$
(101)

converges absolutely uniformly in each compact subset of $D_r(c)$, and $g = \partial_k f$ in $D_r(c)$.

Proof. Let $\rho < r$. Then for any $\varepsilon > 0$ there is a constant $C_{\varepsilon} > 0$ such that

$$\alpha_k \|a_\alpha\|\rho^{|\alpha|-1} \le C_{\varepsilon} (1+\varepsilon)^{|\alpha|} \|a_\alpha\|\rho^{|\alpha|} \le C_{\varepsilon} (1+\varepsilon)^{|\alpha|} (\rho/r)^{|\alpha|} \|a_\alpha\|r^{|\alpha|}$$

and choosing ε small enough we see that (101) converges absolutely uniformly in $D_{\rho}(c)$.

Now we will show that $\partial_k f = g$ in $D_r(c)$, i.e., that for each $z \in D_r(c)$ one has

$$f(z + he_k) = f(z) + g(z)h + o(|h|).$$

Without loss of generality, assuming that k = n and c = 0, we write

$$f(z+he_n) - f(z) = \sum_{\alpha} a_{\alpha} \left((z+he_n)^{\alpha} - z^{\alpha} \right)$$
$$= h \sum_{\alpha} a_{\alpha} z_1^{\alpha_1} \cdots z_{n-1}^{\alpha_{n-1}} \sum_{j=0}^{\alpha_n - 1} (z_n+h)^j z_n^{\alpha_n - 1-j}$$
$$=: h\lambda_z(h).$$

Let $\rho < r$ be such that $z \in D_{\rho}(0)$, and consider all h satisfying $|z_n + h| \leq \rho$. Then

$$\sum_{\alpha} \|a_{\alpha}\| |z_{1}|^{\alpha_{1}} \cdots |z_{n-1}|^{\alpha_{n-1}} \sum_{j=0}^{\alpha_{n-1}} |z_{n}+h|^{j} |z_{n}|^{\alpha_{n-1}-j} \le \sum_{\alpha} \|a_{n}\| \alpha_{n} \rho^{|\alpha|-1} < \infty,$$

so the series for λ_z converges absolutely uniformly in a neighbourhood of the origin. Hence λ_z is continuous at 0, and moreover from $\lambda_z(0) = g(z)$, we infer

$$\lambda_z(h) = g(z) + o(1),$$

with $o(1) \to 0$ as $|h| \to 0$. Therefore

$$f(z + he_n) - f(z) = h(g(z) + o(1)) = hg(z) + o(|h|),$$

and the proof is established.

By repeatedly applying Theorem 49 we see that the coefficients of the power series of f about $c \in \Omega$ are given by

$$a_{\alpha} = \frac{\partial^{\alpha} f(c)}{\alpha!} \equiv \frac{\partial_{1}^{\alpha_{1}} \dots \partial_{n}^{\alpha_{n}} f(c)}{\alpha_{1}! \cdots \alpha_{n}!},$$
(102)

where we have introduced the convention $\alpha! = \alpha_1! \cdots \alpha_n!$. In other words, if $f \in C^{\omega}(\Omega)$ and $c \in \Omega$ then the following *Taylor series* converges in a neighbourhood of c.

$$f(z) = \sum_{\alpha} \frac{\partial^{\alpha} f(c)}{\alpha!} (z - c)^{\alpha}.$$
(103)

We have the identity theorem for multivariate analytic functions, which is necessarily a bit weaker than its single variable counterpart. Namely, the zeros of a multivariate analytic function can form a non-discrete set. For example, the zero set of $f(z) = z_1$ in \mathbb{R}^2 is $\{0\} \times \mathbb{R}$.

Theorem 50 (Identity theorem). Let $f \in C^{\omega}(\Omega, X)$ with Ω a connected open set in \mathbb{K}^n , and with some $b \in \Omega$, let $\partial^{\alpha} f(b) = 0$ for all α . Then $f \equiv 0$ in Ω . In particular, the same conclusion holds if f vanishes on some open subset of Ω .

Proof. Each $\Sigma_{\alpha} = \{z \in \Omega : \partial^{\alpha} f(z) = 0\}$ is relatively closed in Ω , so the intersection $\Sigma = \bigcap_{\alpha} \Sigma_{\alpha}$ is also closed. But Σ is also open, because $z \in \Sigma$ implies that $f \equiv 0$ in a neighbourhood of z by a Taylor series argument. Since $b \in \Sigma$, Σ is nonempty, implying that $\Sigma = \Omega$.

Exercise 51. Let X and Y be Banach spaces, and let

$$F(z) = \sum_{\alpha} A_{\alpha} z^{\alpha}, \quad \text{and} \quad g(z) = \sum_{\alpha} b_{\alpha} z^{\alpha}, \quad (104)$$

be two powers series convergent in $D_r(0)$, with values in L(X, Y) and X, respectively. Show that

$$F(z)g(z) = \sum_{\alpha} \left(\sum_{\beta \le \alpha} A_{\beta} b_{\alpha-\beta} \right) z^{\alpha}, \tag{105}$$

with the series converging absolutely and uniformly in each compact subset of $D_r(0)$.

Exercise 52. Assume the same setting as in the previous exercise, except that now g is a Y-valued power series. Assume also that the coefficient $A_0 \in L(X,Y)$ of F is invertible. Then show that

$$[F(z)]^{-1}g(z) = \sum_{\alpha} e_{\alpha} z^{\alpha}, \quad \text{with} \quad e_{\alpha} = A_0^{-1} \left(b_{\alpha} - \sum_{\{\beta \le \alpha, \beta \ne \alpha\}} A_{\alpha-\beta} e_{\beta} \right), \quad (106)$$

where the power series converges in a neighbourhood of 0.