# Analytic regularity and nonlinear approximation of a class of parametric semilinear elliptic PDEs 

## Report

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# Analytic regularity and nonlinear approximation of a class of parametric semilinear elliptic PDEs* 

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# Analytic Regularity and Nonlinear Approximation of a class of parametric semilinear elliptic PDEs * 

Markus Hansen and Christoph Schwab

May 17, 2011


#### Abstract

We investigate existence and regularity of a class of semilinear, parametric elliptic PDEs with affine dependence of the principal part of the differential operator on countably many parameters. We establish a-priori estimates and analyticity of the parametric solutions. We establish summability results of coefficient sequences of polynomial chaos type expansions of the parametric solutions in terms of tensorized Taylor-, Legendre- and Chebyshev polynomials on the infinite-dimensional parameter domain. We deduce rates of convergence for $N$ term truncated approximations of expansions of the parametric solution. We also deduce spatial regularity of the solution, and establish convergence rates of $N$-term discretizations of the parametric solutions with respect to these polynomials in parameter space and with respect to a multilevel hierarchy of Finite Element spaces in the spatial domain of the PDE.


Key Words: Semilinear Elliptic Partial Differential Equations, Infinite Dimensional Spaces, $N$-term approximation, Analyticity in Infinite Dimensional Spaces, Tensor Product Taylor-, Legendre- and Chebyshev polynomial Approximation.
AMS Mathematics Subject Classification: Primary: 35J61 Secondary: 35B30, 65N30, 41A58

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## 1 Introduction

The numerical solution of partial differential equations (PDEs) with random inputs has received considerable attention in recent years. This is in particular motivated by questions around uncertainty quantification, where PDEs with random inputs are to be solved efficiently for all values taken by the uncertain input parameters. Depending on the number of parameters, there arises the question of efficient numerical solution of such PDEs on possibly high-dimensional parameter spaces. In the context of PDEs with spatially inhomogeneous, random coefficients, parametric expansions of such coefficients, such as Karhunen-Loève expansions, give rise to parametric, deterministic PDEs on infinite dimensional parameter spaces (the parameters being the coefficients in the Karhunen-Loève expansions). For linear partial differential equations, it has been initiated in [14] and extended in $[1,2,7,8]$ to rather general classes of sparse Karhunen-Loève expansions that the solutions of these parametric, deterministic PDEs admit correspondingly sparse expansions in terms of tensorized polynomial systems on the infinite dimensional parameter space, under the provision of sparsity of the random inputs. To extend these results to a class of nonlinear, elliptic PDEs is the purpose of the present paper.
The analysis in $[2,7,8]$ involved in a crucial way the analytic dependence of the solution on the parameters; the extension to the present, nonlinear case involves therefore likewise analyticity of the nonlinearity in the PDE, in addition to the usual structural conditions such as convexity resp. monotonicity and polynomial growth conditions at infinity.
The main results of the present paper are as follows: we establish, under the assumption of affine dependence of the differential operator's principal part on the countably many parameters, that the unique, parametric solution of nonlinear, elliptic PDE depends holomorphically on the countably many parameters in the input data. We then prove that the unique solution of the parametric, nonlinear elliptic PDE admits representations in unconditionally convergent power series expansions about the origin in parameter space as well as unconditionally convergent "polynomial chaos" type expansions with respect to countable families of tensorized orthogonal polynomials on the parameter space. We consider in particular expansions into tensorized Legendre and in Chebyshev polynomial expansions; the latter representation generalizes, in the linear case, the results in $[2,7,8]$ and is useful in connection with recent sampling schemes on the parameter domain which are of "compressed sensing type" (see, e.g. [10]).
The outline of this paper is as follows: in Section 2, we introduce the class of problems of interest, in particular its parametric version depending on countably many parameters, and formulate structural assumptions on the problem data and prove basic well-posedness results for this class of PDEs. All parameters as well as data and solutions in this section are real-valued. Section 3 establishes then the existence, uniqueness and holomorphic dependence of solutions for complex-valued parameter vectors $z \in \mathbb{C}^{\mathbb{N}}$. This holomorphic dependence then is exploited in Section 4 for estimating the Taylor coefficients, and establishing their $p$-summability as well as absolute and uniform convergence of the respective Taylor series. In Section 5 these results are extended to systems of tensorized Legendre- and Chebyshev-polynomials. Under increased regularity assumptions on the inputs, we show in Section 6 that also the solution possesses higher spatial regularity, which carries over to the respective series expansions. In Section 7 we finally investigate convergence rates for $N$-term approximation of the solution by Taylor-, Legendre- and Chebyshev-partial sums.

## 2 A class of semilinear elliptic PDEs

In the following we shall consider the Dirichlet boundary value problem

$$
\begin{equation*}
-\nabla(a \nabla u)+G(u)=f \quad \text { in } \quad D, \quad u=0 \quad \text { on } \quad \partial D, \tag{2.1}
\end{equation*}
$$

where $a=a(x)$ is a given fixed function on $D$.
We will denote the differential operator by $T$ consisting of a linear part $L=-\nabla(a \nabla \cdot)$ and the (nonlinear) composition operator $T_{G}=G(\cdot)$. Here we shall not strive for utmost generality of the function $G$. Since we wish to study parametric problems, which (as an intermediate step) will be extended to complex parameters, our approach limits the possible choices for the function $G: \mathbb{C} \longrightarrow \mathbb{C}$. More specifically, since we aim at proving analyticity of the parameter dependence of the solutions, we shall only consider analytic functions $G$. Taking into account integrability requirements (see below), we are basically left with studying polynomial functions $G$.
This leads to studying the problems

$$
\begin{align*}
-\nabla(a \nabla u)+u^{m}=f & \text { in } \quad D \subset \mathbb{R}^{n},  \tag{2.2}\\
u=0 & \text { on } \quad \partial D,
\end{align*}
$$

where the necessary restrictions on the power $m \in \mathbb{N}$ and the dimension $n$ of the domain $D$ will be recalled below.

### 2.1 Basic properties of some composition operators

If we stick to the space $H_{0}^{1}(D)$ for solutions of (2.2), then this automatically yields restrictions on the one hand on the dimension $n$ of the domain $D \subset \mathbb{R}^{n}$ and on the other hand on the exponent $m$. The weak formulation of (2.2),

$$
b_{1}(u, v)+b_{2}(u, v)=l(v),
$$

where $u, v \in H_{0}^{1}(D)$ and

$$
b_{1}(u, v)=\int_{D} a(x) \nabla u(x) \nabla v(x) d x, \quad b_{2, m}(u, v)=\int_{D} u(x)^{m} v(x) d x, \quad l(v)=\int_{D} f(x) v(x) d x
$$

particularly requires $u \in L_{m+1}(D)$ for $b_{2, m}$ to be well-defined on $H_{0}^{1}(D) \times H_{0}^{1}(D)$, hence we obtain the following restrictions on the range of $m$ from the Sobolev embedding theorem:

$$
\begin{array}{cc}
n=1 \text { or } n=2: & m \in \mathbb{N} \\
n=3: H_{0}^{1}(D) \hookrightarrow L_{6}(D) & \text { hence } 1 \leq m \leq 5 \\
n=4: H_{0}^{1}(D) \hookrightarrow L_{4}(D) & \text { hence } 1 \leq m \leq 3 \\
n=5: H_{0}^{1}(D) \hookrightarrow L_{10 / 3}(D) & \text { hence } 1 \leq m \leq 2 \\
n=6: H_{0}^{1}(D) \hookrightarrow L_{3}(D) & \text { hence } 1 \leq m \leq 2 \\
n \geq 7: & m=1 .
\end{array}
$$

We will assume the parameter $m$ to be in this range thoughout this paper and denote the corresponding set of parameters $(n, m)$ by $\mathcal{M}$. Thus apart from the possibility of restricting ourselves to smoother functions, i.e. $u \in H^{s}(D)$ for $s>1$, there is only a very limited number of possible parameters.
In this subsection we shall be concerned with some basic properties of the composition operators $T_{G}$ associated with the functions $G_{m}(t)=t^{m}$. The following lemma contains some mapping properties needed in our further considerations.

Lemma 2.1. Let $(n, m) \in \mathcal{M}$. Further, let $D \subset \mathbb{R}^{n}$ be a Lipschitz domain, and put $G(t)=t^{m}$. Then $T_{G}: H^{1}(D) \longrightarrow H^{-1}(D)$ is bounded and locally Lipschitz continuous, i.e. Lipschitz continuous on every bounded subset of $H_{0}^{1}(D)$.

Proof. By the Sobolev embedding theorem we know $H_{0}^{1}(D) \hookrightarrow L_{m+1}(D)$ with constant $\iota$ (depending on $m$ and on $D$ ) and we find for $u, v \in H_{0}^{1}(D)$ by Hölder's inequality

$$
\begin{aligned}
\left\|G(u) \mid H^{-1}(D)\right\| & =\sup _{\left\|v \mid H_{0}^{1}(D)\right\| \leq 1}\left|\int_{D} u(x)^{m} v(x) d x\right| \\
& \leq \sup _{\left\|v \mid H_{0}^{1}(D)\right\| \leq 1}\left(\int_{D}|u(x)|^{m+1} d x\right)^{m /(m+1)}\left(\int_{D}|v(x)|^{m+1} d x\right)^{1 /(m+1)} \\
& \leq \iota\left\|u\left|L_{m+1}(D)\left\|^{m} \leq \iota^{m+1}\right\| u\right| H_{0}^{1}(D)\right\|^{m}
\end{aligned}
$$

Moreover, we find again by Hölder's inequality

$$
\begin{aligned}
\| G(u) & -G(v)\left|H^{-1}(D) \|=\sup _{\left\|w \mid H_{0}^{1}(D)\right\| \leq 1}\right| \int_{D}\left(u(x)^{m}-v(x)^{m}\right) w(x) d x \mid \\
& =\sup _{\left\|w \mid H_{0}^{1}(D)\right\| \leq 1}\left|\int_{D}(u-v)\left(u^{m-1}+u^{m-2} v+\cdots+u v^{m-2}+v^{m-1}\right) w(x) d x\right| \\
& \leq \sum_{k=0}^{m-1} \sup _{\left\|w \mid H_{0}^{1}(D)\right\| \leq 1}\left|\int_{D}(u-v)\left(u^{m-1-k} v^{k}\right) w(x) d x\right| \\
& \leq \sum_{k=0}^{m-1} \sup _{\left\|w \mid H_{0}^{1}(D)\right\| \leq 1}\left\|u-v\left|L_{m+1}(D)\|\cdot\| w\right| L_{m+1}(D)\right\| \cdot\left(\int_{D}|u|^{(m+1) \frac{m-1-k}{m-1}}|v|^{(m+1) \frac{k}{m-1}} d x\right)^{\frac{m-1}{m+1}} \\
& \leq \iota \sum_{k=0}^{m-1}\left\|u-v\left|L_{m+1}(D)\|\cdot\| u\right| L_{m+1}(D)\right\|^{m-1-k} \cdot\left\|v \mid L_{m+1}(D)\right\|^{k} .
\end{aligned}
$$

Thus if $M \subset H_{0}^{1}(D)$ is bounded with bound $C_{M}$, then $T_{G}$ is Lipschitz continuous on $M$ with Lipschitz constant at most $m \iota^{m+1} C_{M}^{m-1}$.

Remark 2.1. The elementary fact

$$
\begin{equation*}
a^{n}-b^{n}=(a-b)\left(a^{n-1}+a^{n-2} b+a^{n-3} b^{2}+\cdots+a b^{n-2}+b^{n-1}\right) \tag{2.3}
\end{equation*}
$$

has been essential in this argument. A more general version of this result can be found in [12], Theorem 5.3.2/1.

Under the ellipticity assumption

$$
\begin{equation*}
0<r \leq a(x) \leq R<\infty, \quad x \in D \tag{2.4}
\end{equation*}
$$

the problem (2.2), restricted to real-valued functions $a$ and $G$, can be seen to have real-valued solutions $u$ for every real-valued right-hand-side $f$. Without going into details, this follows from the invertibility of $L$ (by Lax-Milgram-theory) and (Lipschitz-)continuity, coercivity and monotonicity of $T$ by standard results on existence of solutions for nonlinear equations, see e.g. [12], Theorem 6.1/1.

We shall, for purposes of analysis, require solutions of (2.1) or (2.2) for complex-valued functions $a$. We first note that it is easily seen that the identity (2.3) and Lemma 2.1 remain valid for complex variables and complex-valued functions, respectively. Moreover, under the generalized ellipticity assumption

$$
\begin{equation*}
0<r \leq \mathfrak{R} a(x) \leq|a(x)| \leq R<\infty, \tag{2.5}
\end{equation*}
$$

the Lax-Milgram Theorem shows the bounded invertibility of $L$ also in this case. However, the argument above for the existence of solutions for (2.2) then would fail, since the mentioned theorem in [12] as well as modifications thereof are only valid for (reflexive) real Banach spaces. Thus in Section 3 we shall use a different approach based on fixed point iterations for spaces over $\mathbb{C}$.

### 2.2 The parametric problem

As explained in the introduction, our main interest is the analysis of a parametric version of the problems (2.1) or (2.2). More precisely, we shall investigate the parametric problem

$$
\begin{equation*}
-\nabla(a(y) \nabla u(y))+G(u(y))=f(y) \quad \text { in } \quad D, \quad u=0 \quad \text { on } \quad \partial D . \tag{2.6}
\end{equation*}
$$

Herein, the mappings $f: U \longrightarrow H^{-1}(D)$ and $a: U \longrightarrow L_{\infty}(D)$ are given, real-valued. As the parameter domain $U$ we choose $U=[-1,1]^{\mathbb{N}}$, i.e. this can be interpreted as countably many parameters each taking values in the unit interval $[-1,1]$. We are particularly interested in problems, where the dependence of the mapping $a$ on the parameter vector $y$ is affine, i.e. it is of the form

$$
a(y) \equiv a(\cdot, y), \quad a(x, y)=\bar{a}(x)+\sum_{j \geq 0} y_{j} \psi_{j}(x), \quad x \in D, y \in U
$$

The functions $\psi_{j} \in L_{\infty}(D)$ are assumed given and real-valued, and additional assumptions will be discussed later on.
We shall extend this problem to complex parameters, i.e. we shall assume that $f$ admits an extension to some $H^{-1}(D)$-valued mapping defined on the complex parameter domain

$$
\mathcal{U}=\left\{z \in \mathbb{C}^{\mathbb{N}}:\left|z_{j}\right| \leq 1, j \in \mathbb{N}\right\} \equiv\{z \in \mathbb{C}:|z| \leq 1\}^{\mathbb{N}}
$$

Clearly, an analytic extension of $a$ is given by

$$
a(x, z)=\bar{a}(x)+\sum_{j \geq 0} z_{j} \psi_{j}(x), \quad x \in D, z \in \mathcal{U}
$$

where we now permit $\bar{a}$ and $\psi_{j}, j \in \mathbb{N}$, to be complex-valued. Then our main assumption concerning the function $a$ will be the following counterpart of condition (2.5):

Assumption 1. There exist constants $0<r \leq R<\infty$, such that we have

$$
\begin{equation*}
0<r \leq \mathfrak{\Re} a(x, z) \leq|a(x, z)| \leq R<\infty \quad \text { for a.e. } x \in D \quad \text { and all } z \in \mathcal{U} \tag{2.7}
\end{equation*}
$$

In other words, every function $a(z)=a(\cdot, z)$ satisfies condition (2.5), and the respective upper and lower bounds may be chosen independent of $z \in \mathcal{U}$. In view of this last remark, the following assumption would seem more natural

Assumption 2. For every $z \in \mathbb{C}$ there exist constants $0<r(z) \leq R(z)<\infty$, such that we have

$$
\begin{equation*}
0<r(z) \leq \mathfrak{R} a(x, z) \leq|a(x, z)| \leq R(z)<\infty \quad \text { for a.e. } x \in D . \tag{2.8}
\end{equation*}
$$

The relation between these assumptions is the content of the following lemma.
Lemma 2.2. Assume

$$
\sup _{z \in \mathcal{U}}\left\|\sum_{j>J} z_{j} \psi_{j} \mid L_{\infty}(D)\right\| \longrightarrow 0 \quad(J \rightarrow \infty)
$$

Then the Assumptions 1 and 2 are equivalent.

Proof. Obviously, Assumption 1 implies Assumption 2, hence it remains the reverse implication. The crucial step to see this consists in proving that the mapping $z \mapsto a(z)$ is continuous as a mapping from $\mathcal{U}$, equipped with the product topology, into $L_{\infty}(D)$.
To prove continuity, let $\varepsilon>0$ and a convergent sequence $\left(z^{k}\right)_{k \in \mathbb{N}} \subset \mathcal{U}$ be given. In the product topology convergence towards $z \in \mathcal{U}$ means that for every $j \in \mathbb{N}$ we have $z_{j}^{k} \rightarrow z_{j}$ as $k \rightarrow \infty$. Thus for every $J \in \mathbb{N}$ we can find some $k_{0}(J, \varepsilon)$ such that

$$
\sum_{j=1}^{J}\left|z_{j}^{k}-z_{j}\right| \cdot\left\|\psi_{j} \mid L_{\infty}(D)\right\| \leq \frac{\varepsilon}{3} \quad \text { for all } \quad k \geq k_{0}
$$

Moreover, we find $J_{0}(\varepsilon)$ such that

$$
\sup _{w \in \mathcal{U}}\left\|\sum_{j>J} w_{j} \psi_{j} \mid L_{\infty}(D)\right\| \leq \frac{\varepsilon}{3} \quad \text { for all } \quad J \geq J_{0} .
$$

Then it holds for all $k \geq k_{1}(\varepsilon) \equiv k_{0}\left(J_{0}(\varepsilon), \varepsilon\right)$

$$
\begin{aligned}
&\left\|a\left(z^{k}\right)-a(z) \mid L_{\infty}(D)\right\| \leq\left\|a\left(z^{k}\right)-\bar{a}-\sum_{j=1}^{J} z_{j}^{k} \psi_{j}\left|L_{\infty}(D)\|+\| \sum_{j=1}^{J} z_{j}^{k} \psi_{j}-\sum_{j=1}^{J} z_{j} \psi_{j}\right| L_{\infty}(D)\right\| \\
&+\left\|a(z)-\bar{a}-\sum_{j=1}^{J} z_{j} \psi_{j} \mid L_{\infty}(D)\right\| \\
& \leq \frac{\varepsilon}{3}+\sum_{j=1}^{J}\left|z_{j}^{k}-z_{j}\right| \cdot\left\|\psi_{j} \mid L_{\infty}(D)\right\|+\frac{\varepsilon}{3} \leq \varepsilon
\end{aligned}
$$

The result itself now is a consequence of the compactness of $\mathcal{U}$ (which in turn follows from Tychonoff's Theorem). The upper estimate $\sup _{z \in \mathcal{U}} R(z)<\infty$ is immediate, and the lower estimate follows from the observation that due to (2.8) also $z \mapsto \frac{1}{a(z)}$ is a well-defined continuous mapping from $\mathcal{U}$ into $L_{\infty}(D)$ as well as

$$
\sup _{z \in \mathcal{U}} r(z)^{-1}=\left(\inf _{z \in \mathcal{U}} r(z)\right)^{-1}, \quad \text { i.e. } \quad \sup _{z \in \mathcal{U}} \frac{1}{r(z)}<\infty \Longleftrightarrow \inf _{z \in \mathcal{U}} r(z)>0 .
$$

This completes the proof.
The extra-condition in the above lemma (uniform approximation by finitely many parameters) will later on automatically be satisfied due to assumptions on the functions $\psi_{j}$, cf. Theorem 4.1.
In the linear case, i.e. without the term $G(u)$, in problem (2.6), the ellipticity Assumption 1 already suffices to prove existence and uniqueness of solutions for the respective parametric problem for every fixed parameter $z \in \mathcal{U}$ (Lax-Milgram theory). These arguments fail in the above semi-linear case. Instead in the next section we will use an approach via fixed point assertions.

## 3 Existence and regularity of solutions for complex parameters

In this section we are going to study existence and uniqueness of solutions $u \in H_{0}^{1}(D)$ for the semilinear problem (2.1), and thus also solutions $u(z) \in H_{0}^{1}(D)$ for the parametric problem (2.6) for every fixed parameter $z \in \mathcal{U}$. Afterward we shall investigate the regularity of the mapping $z \longmapsto u(z)$.

### 3.1 A fixed point iteration

The approach is based on the mentioned well-known fact the linear problem

$$
\begin{equation*}
-\nabla(a \nabla u)=f \quad \text { in } \quad D, \quad u=0 \quad \text { on } \quad \partial D, \tag{3.1}
\end{equation*}
$$

has a uniquely determined solution $u_{\text {lin }}$ whenever $a$ satisfies condition (2.5). This is an immediate consequence of the Lax-Milgram-Theorem. Upon denoting the corresponding linear differential operator by $T^{\text {lin }}: H_{0}^{1}(D) \longrightarrow H^{-1}(D)$, i.e. $T^{\operatorname{lin}}(u)=-\nabla(a \nabla u)$, and its bounded linear inverse by $S^{\text {lin }}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)$, we further find that every solution $u$ of the semi-linear problem (2.1) satisfies

$$
u=S^{\operatorname{lin}}(f-G(u))=u_{\operatorname{lin}}-S^{\operatorname{lin}} \circ G(u) .
$$

This gives rise to a fixed point problem for the (nonlinear) operator $T: H_{0}^{1}(D) \longrightarrow H_{0}^{1}(D)$, defined by $T u=u_{\text {lin }}-S^{\text {lin }} \circ G(u)$.
Since $S^{\text {lin }}$ is linear, it immediately follows $T u-T v=-S^{\operatorname{lin}}(G(u)-G(v))$, and

$$
\left\|T u-T v\left|H_{0}^{1}(D)\|\leq\| S^{\operatorname{lin}}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\|\cdot\| G(u)-G(v)\right| H^{-1}(D)\right\|, \quad u, v \in H_{0}^{1}(D)
$$

In order to show the operator to be a contraction we have to suppose that the composition operator $T_{G}$ is locally Lipschitz with appropriate bound on the Lipschitz constant.

Theorem 3.1. Let $a$ satisfy the condition (2.5) with lower bound $r$, and put

$$
\rho=\frac{\left\|f \mid H^{-1}(D)\right\|}{r}
$$

Let the function $G$ fulfill $G(0)=0$, and let the composition operator $T_{G}: H_{0}^{1}(D) \longrightarrow H^{-1}(D)$ be bounded and locally Lipschitz. More precisely, assume

$$
\left\|G(u)-G(v)\left|H^{-1}(D)\|\leq L(2 \rho)\| u-v\right| H_{0}^{1}(D)\right\|
$$

for all $u, v \in B_{2 \rho}(0)=\left\{w \in H_{0}^{1}(D):\left\|w \mid H_{0}^{1}(D)\right\| \leq 2 \rho\right\}$, where the constant $L(2 \rho)$ satisfies

$$
\begin{equation*}
L(2 \rho)\left\|S^{\operatorname{lin}}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \leq \frac{1}{2} \tag{3.2}
\end{equation*}
$$

Then the problem (2.1) admits a uniquely determined solution $u \in H_{0}^{1}(D)$ satisfying

$$
\left\|u-u_{\operatorname{lin}} \mid H_{0}^{1}(D)\right\| \leq \rho
$$

where $u_{\text {lin }}$ is the solution of the linear problem (3.1).

Proof. We want to apply Banach's fixed point theorem. Therefore, we consider the operator $T$ on the set $M=B_{\rho}\left(u_{\operatorname{lin}}\right)=\left\{v \in H_{0}^{1}(D):\left\|v-u_{\operatorname{lin}} \mid H_{0}^{1}(D)\right\| \leq \rho\right\}$. Due to the well-known a priori estimate $\left\|u_{\text {lin }} \mid H_{0}^{1}(D)\right\| \leq \rho$, we find $\left\|v \mid H_{0}^{1}(D)\right\| \leq 2 \rho$ for all $v \in M$. Moreover, $M$ is a closed subset of a Banach space, hence it can be considered as a complete metric space.
Furthermore, we find for $u \in M$

$$
\begin{aligned}
\left\|T u-u_{\operatorname{lin}} \mid H_{0}^{1}(D)\right\| & =\left\|S^{\operatorname{lin}} \circ G(u)\left|H_{0}^{1}(D)\|\leq\| S^{\operatorname{lin}}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\|\cdot\| G(u)\right| H^{-1}(D)\right\| \\
& \leq L(2 \rho)\left\|S^{\operatorname{lin}}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \cdot\left\|u \mid H_{0}^{1}(D)\right\| \leq \rho .
\end{aligned}
$$

Similarly, the assumed local Lipschitz continuity shows that $T: M \longrightarrow M$ is a contraction. Now an application of Banach's fixed point theorem proves the claim, where $u$ is the corresponding fixed point of $T$ in $B_{\rho}\left(u_{\text {lin }}\right)$.
The last assertion of the Theorem, $u \in B_{\rho}\left(u_{\mathrm{lin}}\right) \subset B_{2 \rho}(0)$, can be interpreted as an a-priori estimate, hence we obtain

$$
\begin{equation*}
\left\|u \mid H_{0}^{1}(D)\right\| \leq 2 \rho \tag{3.3}
\end{equation*}
$$

Moreover, when checking whether the constants $L(\varrho)$ satisfy the assumption (3.2), we can further apply the a-priori estimates for solutions of the linear problem (3.1), which can be formulated as $\left\|S^{\text {lin }}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \leq r^{-1}$. Together with Lemma 2.1 we then find

$$
L(2 \rho)\left\|S: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \leq m \iota(m+1)^{m+1}(2 \rho)^{m-1} \cdot r^{-1}
$$

Thus (3.2) turns out to be a smallness assumption for $\rho$ and therefore for $f \in H^{-1}(D)$.
Remark 3.1. Note that without the restriction to "solutions close to the linear ones", i.e. without the additional condition

$$
\left\|u-u_{\operatorname{lin}} \mid H_{0}^{1}(D)\right\| \leq \frac{\left\|f \mid H^{-1}(D)\right\|}{r}
$$

the uniqueness of the solution generally fails. For example the (real) problem for $G(t)=-t|t|^{p-2}$ and $a(x) \equiv 1$ (with certain restrictions on $p$ ) even has infinitely many solutions with arbitrarily large norms, see [13], Theorem 7.2 and Remark 7.3.

We now return to the parametric prolem (2.6). Then clearly the solution $u(z)$ can be obtained by applying Theorem 3.1 for every fixed parameter $z \in \mathcal{U}$. It only remains to check whether the condition (3.2) can be fulfilled.

Theorem 3.2. Consider the problem (2.1) with $G(\zeta)=\zeta^{m}$ for $\zeta \in \mathbb{C}$ and $(n, m) \in \mathcal{M}$. Let $a$ satisfy Assumption 1, and suppose

$$
\begin{equation*}
\left\|f(z) \mid H^{-1}(D)\right\|^{m-1}<\frac{r^{m}}{2^{m} m \iota(m+1)^{m+1}} \quad \text { for all } z \in \mathcal{U} \tag{3.4}
\end{equation*}
$$

Then problem (2.1) admits a uniquely determined solution $u(z) \in H_{0}^{1}(D)$ for every parameter $z \in \mathcal{U}$, satisfying

$$
\begin{equation*}
\left\|u(z)-u_{\operatorname{lin}}(z) \mid H_{0}^{1}(D)\right\| \leq \frac{M_{f}}{r} \quad \text { and } \quad\left\|u(z) \mid H_{0}^{1}(D)\right\| \leq \frac{2 M_{f}}{r} \tag{3.5}
\end{equation*}
$$

where $u_{\operatorname{lin}}(z)$ is the solution of the corresponding linear problem and $M_{f}=\sup _{z \in \mathcal{U}}\left\|f(z) \mid H^{-1}(D)\right\|$.

Note that if $z \mapsto f(z)$ is continuous on $\mathcal{U}$, then automatically $M_{f}<\infty$, and the smallness condition then reads as $M_{f}^{m-1}<\frac{r^{m}}{2^{m} m \iota(m+1)^{m+1}}$.

### 3.2 A fixed point iteration for Helmholtz-type problems.

Before we come to a stability assertion for the problem (2.1), we shall have a look at existence of solutions for Helmholtz-type equations

$$
\begin{equation*}
-\nabla(a \nabla u)+b u=f \quad \text { in } \quad D, \quad u=0 \quad \text { on } \quad \partial D, \tag{3.6}
\end{equation*}
$$

where the function $b$ may be complex-valued. These problems are particularly simple cases of the general problem (2.1), where the (nonlinear) composition operator $T_{G}$ is replaced by a (linear) multiplication operator $T^{b}$, where $T^{b} u=b u$.

Proposition 3.1. Consider the problem (3.6), where $a$ satisfies condition (2.5). Moreover, assume

$$
\begin{equation*}
\left\|b u\left|H^{-1}(D)\left\|\leq C_{b}\right\| u\right| H_{0}^{1}(D)\right\| \tag{3.7}
\end{equation*}
$$

for every $u \in H_{0}^{1}(D)$, where the constant $C_{b} \equiv\left\|T^{b}: H_{0}^{1}(D) \longrightarrow H^{-1}(D)\right\|$ satisfies

$$
\begin{equation*}
C_{b}\left\|S^{\operatorname{lin}}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \leq \frac{1}{2} \tag{3.8}
\end{equation*}
$$

and where the operator $S$ has the same meaning as in the previous section. Then there exists a uniquely determined solution $u \in H_{0}^{1}(D)$ of (3.6) satisfying

$$
\left\|u-u_{\text {lin }} \mid H_{0}^{1}(D)\right\| \leq \rho=\frac{\left\|f \mid H^{-1}(D)\right\|}{r} \quad \text { and } \quad\left\|u \mid H_{0}^{1}(D)\right\| \leq 2 \rho,
$$

where $u_{\operatorname{lin}}=S^{\operatorname{lin} f}$ is the unique solution of the problem (3.1).

Proof. As before, we shall use a fixed point iteration to obtain the desired result from Banach's fixed point theorem. To this end, consider for $u \in H_{0}^{1}(D)$ the operator $T$,

$$
T u=u_{\operatorname{lin}}-S^{\operatorname{lin}} \circ T^{b} u=u_{\operatorname{lin}}-S^{\operatorname{lin}}(b u) .
$$

Then we find by Assumption (3.7)

$$
\begin{aligned}
\left\|T u-u_{\text {lin }} \mid H_{0}^{1}(D)\right\| & \leq\left\|S^{\text {lin }}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \cdot\left\|b u \mid H^{-1}(D)\right\| \\
& \leq C_{b}\left\|S^{\operatorname{lin}}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \cdot\left\|u \mid H_{0}^{1}(D)\right\| .
\end{aligned}
$$

Hence $T$ is a mapping from $B_{\rho}\left(u_{\text {lin }}\right)$ into itself, since $\left\|u_{\text {lin }} \mid H_{0}^{1}(D)\right\| \leq \rho$ and $\left\|u \mid H_{0}^{1}(D)\right\| \leq 2 \rho$ for all $u \in B_{\rho}\left(u_{\text {lin }}\right)$. Moreover, we similarly obtain

$$
\begin{aligned}
\left\|T u-T v \mid H_{0}^{1}(D)\right\| & =\left\|S^{\operatorname{lin}} \circ T^{b}(u-v) \mid H_{0}^{1}(D)\right\| \\
& \leq\left\|S^{\operatorname{lin}}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \cdot\left\|b(u-v) \mid H^{-1}(D)\right\| \\
& \leq C_{b}\left\|S^{\operatorname{lin}}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \cdot\left\|u-v \mid H_{0}^{1}(D)\right\|,
\end{aligned}
$$

which shows that $T$ is a contraction. An application of Banach's fixed point theorem now shows the existence of precisely one fixed point of $T$ in $B_{\rho}\left(u_{\text {lin }}\right)$, which is the desired solution of (3.6).

As opposed to the semilinear problem (2.1) we this time have a smallness condition on the input function $b$. However, in the following application of this result, this in turn will again be reformulated as a smallness condition for $f$.

Proposition 3.2. Let $u$ and $\widetilde{u}$ be solutions of (2.1) for $G(\zeta)=\zeta^{m}$ and $(n, m) \in \mathcal{M}$ with data $a$ and $\widetilde{a}$, both satisfying (2.5) with lower bound $r$, and the right hand sides $f, \widetilde{f}$. If $f$ and $\widetilde{f}$ satisfy the condition

$$
\begin{equation*}
\max \left(\left\|f\left|H^{-1}(D)\|,\| \widetilde{f}\right| H^{-1}(D)\right\|\right)<r^{\frac{m}{m-1}}\left(2^{m} m \iota(m+1)^{m+1}\right)^{-\frac{1}{m-1}} \tag{3.9}
\end{equation*}
$$

then it holds

$$
\left\|u-\widetilde{u} \mid H_{0}^{1}(D)\right\| \leq \frac{2}{r}\left(\left\|f-\widetilde{f}\left|H^{-1}(D)\|+2\| a-\widetilde{a}\right| L_{\infty}(D)\right\| \cdot \frac{\left\|\widetilde{f} \mid H^{-1}(D)\right\|}{r}\right)
$$

Proof. Subtracting the corresponding problems for $u$ and $\widetilde{u}$ yields with $w=u-\widetilde{u}$

$$
\begin{aligned}
\int_{D}(f-\widetilde{f}) v d x & =\int_{D} a \nabla u \cdot \nabla v d x-\int_{D} \widetilde{a} \nabla \widetilde{u} \cdot \nabla v d x+\int_{D}(G(u)-G(\widetilde{u})) v d x \\
& =\int_{D} a \nabla w \cdot \nabla v d x+\int_{D}(a-\widetilde{a}) \nabla \widetilde{u} \cdot \nabla v d x+\int_{D} v w \sum_{j=0}^{m-1} u^{m-1-j} \widetilde{u}^{j} d x
\end{aligned}
$$

Defining

$$
b(x)=\sum_{j=0}^{m-1} u^{m-1-j}(x) \widetilde{u}^{j}(x),
$$

then $w=u-\widetilde{u}$ is a weak solution of the Helmholtz-type problem

$$
-\nabla(a \nabla w)+b w=\ell \quad \text { in } D,\left.\quad w\right|_{\partial D}=0,
$$

where

$$
\ell(v)=\int_{D}(f-\widetilde{f}) v d x+\int_{D}(\widetilde{a}-a) \nabla \widetilde{u} \cdot \nabla v d x .
$$

To apply the previous proposition we have to check condition (3.7). Since we have the a priori estimates $\left\|u \mid H_{0}^{1}(D)\right\| \leq 2 \rho$ and $\left\|\widetilde{u} \mid H_{0}^{1}(D)\right\| \leq 2 \rho$, where $\rho=r^{-1} \max \left(\left\|f\left|H^{-1}(D)\|,\| \widetilde{f}\right| H^{-1}(D)\right\|\right)$, we find

$$
\begin{aligned}
\left\|T^{b} w \mid H^{-1}(D)\right\| & \leq \iota(m+1)^{m+1} \sum_{j=0}^{m-1}\left\|u\left|H_{0}^{1}(D)\left\|^{m-1-j} \cdot\right\| \widetilde{u}\right| H_{0}^{1}(D)\right\|^{j} \cdot\left\|w \mid H_{0}^{1}(D)\right\| \\
& \leq \iota(m+1)^{m+1} m(2 \rho)^{m-1} \cdot\left\|w \mid H_{0}^{1}(D)\right\|
\end{aligned}
$$

In other words, $T^{b}$ is bounded (and hence globally Lipschitz continuous) with norm at most $C:=$ $\iota(m+1)^{m+1} m(2 \rho)^{m-1}$. Moreover, we obtain

$$
\begin{aligned}
\left\|\ell \mid H^{-1}(D)\right\| & \leq\left\|f-\widetilde{f}\left|H^{-1}(D)\|+\| a-\widetilde{a}\right| L_{\infty}(D)\right\| \cdot\left\|\widetilde{u} \mid H_{0}^{1}(D)\right\| \\
& \leq\left\|f-\widetilde{f}\left|H^{-1}(D)\|+\| a-\widetilde{a}\right| L_{\infty}(D)\right\| \cdot r^{-1}\left\|\widetilde{f} \mid H^{-1}(D)\right\| .
\end{aligned}
$$

In this case, the counterpart of condition (3.7) reads as

$$
C\left\|S^{\operatorname{lin}}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \leq \frac{1}{2}
$$

which can be reformulated as the imposed smallness-condition on $f$. Applying Proposition 3.1 now yields that there is a uniquely determined solution $w_{0}$ of (3.6) in $B_{\widetilde{\rho}}\left(S^{\operatorname{lin}} \ell\right)$, where $\widetilde{\rho}=r^{-1}\left\|\ell \mid H^{-1}(D)\right\|$.

Thus it remains to show that this solution coincides with $w$. To this end we note that also $w$ is a fixed point of $T$ (recall $T u=S^{\text {lin }} \ell-S^{\text {lin }} \circ T^{b} u$ ), and that $T$ is a contraction on the whole space $H_{0}^{1}(D)$. Then we obtain

$$
\left\|w-w_{0}\left|H_{0}^{1}(D)\|=\| T\left(w-w_{0}\right)\right| H_{0}^{1}(D)\right\| \leq C_{0}\left\|S^{\operatorname{lin}}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\| \cdot\left\|w-w_{0}\right\|,
$$

which for $w \neq w_{0}$ yields a contradiction. Finally, the a priori estimate from Proposition 3.1 yields $\left\|w \mid H_{0}^{1}(D)\right\| \leq 2 \widetilde{\rho}$, which completes the proof.
Note that a priori it is not clear, whether $w \in B_{\tilde{\rho}}(S \ell)$, hence we used this indirect approach to show that the solutions coincide.

### 3.3 Partial derivatives

For the arguments used in the later sections it will be necessary to further extend the considered parameter domains. We follow [2] and define for $\delta>0$ domains

$$
\mathcal{A}_{\delta}=\left\{z \in \mathbb{C}^{\mathbb{N}}: \delta \leq \mathfrak{R} a(x, z) \leq|a(x, z)| \leq R(z)<\infty \text { for every } x \in D\right\} .
$$

If $a$ fulfills Assumption 1 , then $\mathcal{A}_{\delta} \supset \mathcal{U}$ for every $0<\delta \leq r$. Moreover, we put $\mathcal{A}=\bigcup_{\delta>0} \mathcal{A}_{\delta}$, and assume $z \mapsto f(z)$ to admit an extension from $\mathcal{U}$ to $\mathcal{A}$.

Furthermore, we suppose the data $f(z) \in H^{-1}(D)$ to satisfy the smallness-condition

$$
\begin{equation*}
\left\|f(z) \mid H^{-1}(D)\right\|<\left(\frac{\delta^{m}}{2^{m} m \iota(m+1)^{m+1}}\right)^{\frac{1}{m-1}} \quad \text { for all } z \in \mathcal{A}_{\delta} \text { and all } 0<\delta \leq r \tag{3.10}
\end{equation*}
$$

Then the assumptions of both Theorem 3.2 and Proposition 3.2 are fulfilled. Note that there is no further assumption on the function $a$ or the parameter $z$, apart from the ellipticity condition (2.5). Hence for every $z \in \mathcal{A}_{\delta}$ we have a uniquely determined solution $u(z) \in H_{0}^{1}(D)$ satisfying the estimate

$$
\begin{equation*}
\left\|u(z)\left|H_{0}^{1}(D)\left\|\leq \frac{2\left\|f(z) \mid H^{-1}(D)\right\|}{\delta} \leq M(\delta)=2 \delta^{-1} \sup _{z \in \mathcal{A}_{\delta}}\right\| f(z)\right| H^{-1}(D)\right\| \tag{3.11}
\end{equation*}
$$

where the finiteness of the supremum follows from (3.10).
With this estimates in hand, we begin with a preparation.
Proposition 3.3. Let $(n, m) \in \mathcal{M}$, and consider $G(\zeta)=\zeta^{m}$ with $m \geq 2$. Then it holds

$$
\lim _{h \rightarrow 0} \frac{\left\|G\left(u\left(z+h e^{j}\right)\right)-G(u(z))-G^{\prime}(u(z))\left(u\left(z+h e^{j}\right)-u(z)\right) \mid H^{-1}(D)\right\|}{|h|}=0 .
$$

Note that although this expression appears reminiscent of Fréchet or Gâteaux differentiability, it does not imply either (neither is $G^{\prime}(u(z))\left(u\left(z+\cdot e^{j}\right)-u(z)\right)$ a (bounded) linear mapping, nor is $u\left(z+h e^{j}\right)-u(z)$ a (fixed) direction).
Proof. Fix $z \in \mathcal{A}$, and choose $\delta>0$ such that $z \in \mathcal{A}_{\delta}$. First we note that then for $h \in \mathbb{C} \backslash\{0\}$ we have $z+h e^{j} \in \mathcal{A}_{\delta / 2}$ : With $z \in \mathcal{A}_{\delta}$ and $|h| \cdot\left\|\psi_{j} \mid L_{\infty}(D)\right\|<\delta / 2$ we find

$$
\Re a\left(x, z+h e^{j}\right) \geq \Re a(x, z)-|h| \cdot\left\|\psi_{j} \mid L_{\infty}(D)\right\| \geq \delta-\delta / 2=\delta / 2
$$

as well as

$$
\left|a\left(x, z+h e^{j}\right)\right| \leq|a(x, z)|+|h| \cdot\left\|\psi_{j} \mid L_{\infty}(D)\right\| \leq R(z)+\delta / 2<\infty .
$$

In particular, this yields $u\left(z+h e^{j}\right) \in H_{0}^{1}(D)$ to be well-defined. Next we recall that by Taylor's Theorem it holds

$$
g(z)-g\left(z_{0}\right)-g^{\prime}\left(z_{0}\right)\left(z-z_{0}\right)=\left(z-z_{0}\right)^{2} \int_{0}^{1}(1-t) g^{\prime \prime}\left(z_{0}+t\left(z-z_{0}\right)\right) d t
$$

for every analytic function $g$ (this version easily follows from the $\mathbb{C}$-valued real-variable version by considering the Taylor polynomial for the mapping $t \mapsto g\left(z_{0}+t\left(z-z_{0}\right)\right)$ around 0 , evaluated in $t=1$ ). In particular, this can be applied to the polynomial $G$. Thus we can apply Taylor's formula pointwise (where $z$ corresponds to $u\left(z+h e^{j}\right)$ and $z_{0}$ corresponds to $u(z)$ ) to obtain

$$
\begin{aligned}
G\left(u\left(z+h e^{j}\right)\right) & -G(u(z))-G^{\prime}(u(z))\left(u\left(z+h e^{j}\right)-u(z)\right) \\
& =\left(u\left(z+h e^{j}\right)-u(z)\right)^{2} \int_{0}^{1}(1-t) G^{\prime \prime}\left(u(z)+t\left(u\left(z+h e^{j}\right)-u(z)\right)\right) d t .
\end{aligned}
$$

Multiplying with $v \in H_{0}^{1}(D)$ and integrating over $D$ yields

$$
\begin{aligned}
\int_{D} \mid\left(G\left(u\left(z+h e^{j}\right)\right)-\right. & \left.G(u(z))-G^{\prime}(u(z))\left(u\left(z+h e^{j}\right)-u(z)\right)\right) v(x) \mid d x \\
= & \int_{D}\left|u\left(z+h e^{j}\right)-u(z)\right|^{2}|v(x)| \int_{0}^{1}\left|(1-t) G^{\prime \prime}\left(u(z)+t\left(u\left(z+h e^{j}\right)-u(z)\right)\right)\right| d t d x \\
\leq & \left\|u\left(z+h e^{j}\right)-u(z)\left|L_{m+1}(D)\left\|^{2} \cdot\right\| v(x)\right| L_{m+1}(D)\right\| \\
& \quad \times\left\|\int_{0}^{1}\left|(1-t) G^{\prime \prime}\left(u(z)+t\left(u\left(z+h e^{j}\right)-u(y)\right)\right)\right| d t \left\lvert\, L_{\frac{m+1}{m-2}}(D)\right.\right\|,
\end{aligned}
$$

where in the last line we used Hölder's inequality (clearly the case $m=2$ is trivial). The last factor can be further estimated using the (generalized) Minkowski inequality. Then we obtain

$$
\begin{aligned}
& \left\|\int_{0}^{1}\left|(1-t) G^{\prime \prime}\left(u(z)+t\left(u\left(z+h e^{j}\right)-u(z)\right)\right)\right| d t \left\lvert\, L_{\frac{m+1}{m-2}}(D)\right.\right\| \\
& \quad \leq \int_{0}^{1}\left\|(1-t) G^{\prime \prime}\left(u(z)+t\left(u\left(z+h e^{j}\right)-u(z)\right)\right) \left\lvert\, L_{\frac{m+1}{m-2}}(D)\right.\right\| d t \\
& \quad=\int_{0}^{1}(1-t) m(m-1)\left\|\left|u(z)+t\left(u\left(z+h e^{j}\right)-u(z)\right)\right|^{m-2} \left\lvert\, L_{\frac{m+1}{m-2}}(D)\right.\right\| d t \\
& \quad \leq \int_{0}^{1}(1-t) m(m-1)\left((1-t)\left\|u(z)\left|L_{m+1}(D)\|+t\| u\left(z+h e^{j}\right)\right| L_{m+1}(D)\right\|\right)^{m-2} d t \\
& \quad \leq m(m-1) \int_{0}^{1}(1-t) d t(\iota(m+1) M(\delta))^{m-2}=\frac{1}{2} m(m-1)(\iota(m+1) M(\delta))^{m-2}
\end{aligned}
$$

where at the end we used the a priori estimate (Theorems 3.1 and 3.2), and $\iota(m+1)$ is the embedding constant for $H_{0}^{1}(D) \hookrightarrow L_{m+1}(D)$. Combining these two estimates yields together with the stability estimate (Proposition 3.2)

$$
\left\|G\left(u\left(z+h e^{j}\right)\right)-G(u(z))-G^{\prime}(u(z))\left(u\left(z+h e^{j}\right)-u(z)\right) \mid H^{-1}(D)\right\|
$$

$$
\begin{aligned}
& \leq \frac{1}{2} m(m-1) \iota(m+1)^{m+1}\left\|u\left(z+h e^{j}\right)-u(z) \mid H_{0}^{1}(D)\right\|^{2} \cdot M(\delta)^{m-2} \\
& \leq \frac{1}{2} m(m-1) \iota(m+1)^{m+1} M(\delta)^{m-2} 4 \delta^{-2} \\
& \quad \times\left(\left\|f\left(z+h e^{j}\right)-f(z)\left|H^{-1}(D)\|+\| a(\cdot, z+h)-a(\cdot, z)\right| L_{\infty}(D)\right\| M(\delta)\right)^{2} \\
& \leq 2 m(m-1) \iota(m+1)^{m+1} M(\delta)^{m-2} \delta^{-2}\left(\mathcal{O}(|h|)+|h| \cdot\left\|\psi_{j} \mid L_{\infty}(D)\right\| M(\delta)\right)^{2}=\mathcal{O}\left(|h|^{2}\right) .
\end{aligned}
$$

This proves the claim.
We now aim at proving existence of partial derivatives for the mapping $z \mapsto u(z)$. To do so we require an additional assumption on $f$.

Assumption 3. We assume the mapping $z \mapsto f(z)$ to admit partial derivatives $\partial_{z_{j}} f(z) \in H^{-1}(D)$ at every point $z \in \mathcal{A}$ with respect to every variable $z_{j}, j \in \mathbb{N}$.

Then the desired regularity result for $z \mapsto u(z)$ reads as follows.
Theorem 3.3. Let $z \mapsto f(z) \in H^{-1}(D)$ satisfy the condition (3.10) and Assumption 3. Then at any point $z \in \mathcal{A}$, the function $z \mapsto u(z)$ admits partial derivatives $\partial_{z_{j}} u(z) \in H_{0}^{1}(D)$ with respect to each complex variable $z_{j}$. For each $j \in \mathbb{N}, \partial_{z_{j}} u$ is the unique solution in $H_{0}^{1}(D)$ of the linear problem:

$$
\begin{aligned}
\int_{D} a(x, z) & \nabla w(x) \cdot \nabla v(x) d x+\int_{D} G^{\prime}(u(x, z)) w(x) v(x) d x \\
& =\int_{D} \partial_{z_{j}} f(x, z) v(x, z) d x-\int_{D} \psi_{j}(x) \nabla u(x, z) \cdot \nabla v(x) d x \quad \text { for all } v \in H_{0}^{1}(D)
\end{aligned}
$$

Proof. Fix again $z \in \mathcal{A}$, and choose $\delta>0$ such that $z \in \mathcal{A}_{\delta}$. As in the proof of Proposition 3.3 we have $z+h e^{j} \in \mathcal{A}_{\delta / 2}$ for $|h|$ sufficiently small. For such $h$, the difference quotient

$$
w_{h}(z)=\frac{u\left(z+h e^{j}\right)-u(z)}{h} \in H_{0}^{1}(D)
$$

is well-defined and we find for arbitrary $v \in H_{0}^{1}(D)$ and such $h$ that

$$
\begin{aligned}
& h \int_{D} \frac{f\left(x, z+h e^{j}\right)-f(x, z)-h \partial_{z_{j}} f(x, z)}{h} v(x) d x+h \int_{D} \partial_{z_{j}} f(x, z) v(x) d x \\
& =\int_{D} a\left(x, z+h e^{j}\right) \nabla u\left(x, z+h e^{j}\right) \cdot \nabla v(x) d x-\int_{D} a(x, z) \nabla u(x, z) \cdot \nabla v(x) d x \\
& \quad+\int_{D}\left(G\left(u\left(x, z+h e^{j}\right)\right)-G(u(x, z))\right) v(x) d x \\
& =h \int_{D} a(x, z) \nabla w_{h}(x, z) \cdot \nabla v(x) d x+\int_{D}\left(a\left(x, z+h e^{j}\right)-a(x, z)\right) \nabla u\left(x, z+h e^{j}\right) \cdot \nabla v(x) d x \\
& \quad+\int_{D}\left(G\left(u\left(x, z+h e^{j}\right)\right)-G(u(x, z))\right) v(x) d x \\
& =h \int_{D} a(x, z) \nabla w_{h}(x, z) \cdot \nabla v(x) d x+h \int_{D} \psi_{j}(x) \nabla u\left(x, z+h e^{j}\right) \cdot \nabla v(x) d x \\
& \quad+h \int_{D} \frac{G\left(u\left(x, z+h e^{j}\right)\right)-G(u(x, z))-G^{\prime}(u(x, z))\left(u\left(x, z+h e^{j}\right)-u(x, z)\right)}{h} v(x) d x
\end{aligned}
$$

$$
+h \int_{D} G^{\prime}(u(x, z)) w_{h}(x, z) v(x) d x .
$$

By Proposition 3.3 the third integral in the last line converges to 0 as $h \rightarrow 0$, and the first integral on the left hand side converges to 0 as well due to Assumption 3.
If we define $\ell_{h}(v)=\int_{D} \psi_{j}(x) \nabla u\left(x, z+h e^{j}\right) \cdot \nabla v(x) d x$ and $\ell_{0}(v)=\int_{D} \psi_{j}(x) \nabla u(x, z) \cdot \nabla v(x) d x$, then this yields bounded linear functionals on $H_{0}^{1}(D)$, and together with (3.11) we find

$$
\begin{aligned}
\left|\ell_{h}(v)-\ell_{0}(v)\right| & =\left|\int_{D} \psi_{j}\left(\nabla u\left(z+h e^{j}\right)-\nabla u(z)\right) \cdot \nabla v d x\right| \\
& \leq\left\|\psi_{j}\left|L_{\infty}(D)\|\cdot\| u\left(z+h e^{j}\right)-u(z)\right| H_{0}^{1}(D)\right\| \cdot\left\|v \mid H_{0}^{1}(D)\right\| \\
& \leq \frac{1}{\delta}\left\|\psi_{j}\left|L_{\infty}(D)\left\|\left(\left\|f\left(z+h e^{j}\right)-f(z)\left|H^{-1}(D)\|+|h| \cdot\| \psi_{j}\right| L_{\infty}(D)\right\| M(\delta)\right)\right\| v\right| H_{0}^{1}(D)\right\|
\end{aligned}
$$

In the last line we applied the stability estimate (Proposition 3.2). We conclude $\ell_{h} \rightarrow \ell_{0}$ in $H^{-1}(D)$ as $h \rightarrow 0$. Furthermore, defining

$$
\begin{aligned}
L_{h}(v)= & \int_{D} \frac{f\left(x, z+h e^{j}\right)-f(x, z)-h \partial_{z_{j}} f(x, z)}{h} v(x) d x-\int_{D} \psi_{j}(x) \nabla u\left(x, z+h e^{j}\right) \cdot \nabla v(x) d x \\
& \quad-\int_{D} \frac{G\left(u\left(x, z+h e^{j}\right)\right)-G(u(x, z))-G^{\prime}(u(x, z))\left(u\left(x, z+h e^{j}\right)-u(x, z)\right)}{h} v(x) d x
\end{aligned}
$$

this yields another bounded linear functional on $H_{0}^{1}(D)$, and we find $L_{h} \rightarrow \ell_{0}$ in $H^{-1}(D)$ as $h \rightarrow 0$. The next step consists in an application of Proposition 3.1 with $b(x, z)=G^{\prime}(u(x, z))$, hence we have to check condition (3.7). Similar to Lemma 2.1, applying Hölder's inquality yields

$$
\left\|G^{\prime}(u(z)) v\left|H^{-1}(D)\left\|\leq m \iota(m+1)^{m+1}\right\| u(z)\right| H_{0}^{1}(D)\right\|^{m-1} \cdot\left\|v \mid H_{0}^{1}(D)\right\|, \quad v \in H_{0}^{1}(D)
$$

hence (3.7) is fulfilled with $C_{b} \widehat{=}(2 \rho(z))^{m-1} m \iota(m+1)^{m+1}$, where $\rho(z)=\delta^{-1}\left\|f(z) \mid H^{-1}(D)\right\|$. Thus condition (3.8) reduces once more to a smallness-condition for $f(z)$,

$$
\left\|f(z) \mid H^{-1}(D)\right\| \leq \delta\left(2^{m} m \iota(m+1)^{m+1}\left\|S_{z}^{\text {lin }}: H^{-1}(D) \longrightarrow H_{0}^{1}(D)\right\|\right)^{-1 /(m-1)}
$$

which in turn is satisfied due to assumption (3.10). This implies that the function $w_{h}$ is the unique weak solution of the problem

$$
\int_{D} a(x, z) \nabla w_{h}(x, z) \cdot \nabla v(x) d x+\int_{D} G^{\prime}(u(x, z)) w_{h}(x, z) v(x) d x=\int_{D} \partial_{z_{j}} f(x, z) v(x) d x+L_{h}(v) .
$$

Similarly, $w_{0}$ is the unique weak solution of the problem

$$
\int_{D} a(x, z) \nabla w_{0}(x, z) \cdot \nabla v(x) d x+\int_{D} G^{\prime}(u(x, z)) w_{0}(x, z) v(x) d x=\int_{D} \partial_{z_{j}} f(x, z) v(x) d x+\ell_{0}(v) .
$$

Combining these two problems, we once more obtain from Proposition 3.1 together with Theorem 3.1 the estimate

$$
\left\|w_{h}(z)-w_{0}(z) \mid H_{0}^{1}(D)\right\| \leq \frac{\left\|L_{h}-\ell_{0} \mid H^{-1}(D)\right\|}{r} \rightarrow 0 \quad(h \rightarrow 0) .
$$

Thus $\partial_{z_{j}} u(z)=w_{0}(z)=\lim _{h \rightarrow 0} w_{h}(z)$ exists in $H_{0}^{1}(D)$.

## 4 Estimates of Taylor coefficents

From now on we assume the functions $a$ and $f$ to satisfy Assumption 1, the smallness-condition (3.10) and Assumption 3, respectively.

## $4.1 \quad \delta$-admissible sequences

We remind on the domains

$$
\mathcal{A}_{\delta}=\left\{z \in \mathbb{C}^{\mathbb{N}}: \delta \leq \mathfrak{R} a(x, z) \leq|a(x, z)| \leq R(z)<\infty \text { for every } x \in D\right\} .
$$

where $\delta>0$, and $\mathcal{U} \subset \mathcal{A}_{\delta}$ for $\delta \leq r$ as a consequence of Assumption 1. Following [2] we call a sequence $\rho=\left(\rho_{j}\right)_{j \in \mathbb{N}}$ of positive real numbers $\delta$-admissible, if

$$
\sum_{j=1}^{\infty} \rho_{j}\left|\psi_{j}(x)\right| \leq \Re \bar{a}(x)-\delta
$$

for almost every $x \in D$. This particularly implies that the poly-disc

$$
\mathcal{U}_{\rho}=\prod_{j \in \mathbb{N}}\left\{\zeta \in \mathbb{C}:|\zeta| \leq \rho_{j}\right\}
$$

is contained in $\mathcal{A}_{\delta}$. Indeed, we find for $z \in \mathcal{U}_{\rho}$ and almost every $x \in D$

$$
\Re a(x, z) \geq \Re \bar{a}(x)-\sum_{j=1}^{\infty}\left|z_{j}\right| \cdot\left|\psi_{j}(x)\right| \geq \Re \bar{a}(x)-\sum_{j=1}^{\infty} \rho_{j}\left|\psi_{j}(x)\right| \geq \delta
$$

as well as

$$
|a(x, z)| \leq|\bar{a}(x)|+\sum_{j=1}^{\infty}\left|z_{j}\right| \cdot\left|\psi_{j}(x)\right| \leq|\bar{a}(x)|+\Re \bar{a}(x)-\delta \leq 2|\bar{a}(x)| \leq 2 R .
$$

Similarly, we put

$$
\mathcal{U}_{\rho, E}=\prod_{j \in E}\left\{\zeta \in \mathbb{C}:|\zeta| \leq \rho_{j}\right\}
$$

for every (finite) set $E \subset \mathbb{N}$.

### 4.2 A preliminary estimate for partial derivatives of $u$

We first prove a preliminary result for estimates for partial derivatives of $u(y)$ (and hence for its Taylor coefficients) for general $\delta$-admissible sequences $\rho$.
Before we come to the next proposition, we shall need some more notation. We denote by

$$
\mathcal{F}=\left\{\nu=\left(\nu_{j}\right)_{j \in \mathbb{N}} \in \mathbb{N}_{0}^{\mathbb{N}}:|\operatorname{supp} \nu|<\infty\right\}
$$

the set of all sequences of non-negative integers with at most finitely many non-vanishing components. Here

$$
\operatorname{supp} \nu=\left\{j \in \mathbb{N}: \nu_{j} \neq 0\right\}
$$

denotes the support of $\nu$, and we put $|\nu|=\sum_{j \in \mathbb{N}} \nu_{j}$, which is finite if, and only if, $|\operatorname{supp} \nu|<\infty$.

Proposition 4.1. Let $f \in H^{-1}(D)$ satisfy the condition (3.10). Let a $\delta$-admissible sequence $\rho=$ $\left(\rho_{j}\right)_{j \in \mathbb{N}}$ be given. Then for every $\nu \in \mathcal{F}$ we have the estimate

$$
\left\|\partial^{\nu} u(y) \mid H_{0}^{1}(D)\right\| \leq \nu!M(\delta) \prod_{j \in \operatorname{supp} \nu} \rho_{j}^{-\nu_{j}} .
$$

Proof. The proof is based on the argument in [2, Lemma 2.4].
Put $E=\operatorname{supp} \nu$ and $J=|E|$. Writing $z=\left(z_{E}, z_{E^{c}}\right)$, i.e. $z_{E} \in \mathbb{C}^{J}$ contains the components corresponding to inidces $j \in E$, the $\delta$-admissibility of $\rho$ then mplies the estimate

$$
\begin{equation*}
\left\|u\left(z_{E}, 0\right) \mid H_{0}^{1}(D)\right\| \leq \frac{2\left\|f\left(z_{E}, 0\right) \mid H^{-1}(D)\right\|}{\delta} \leq M(\delta) \tag{4.1}
\end{equation*}
$$

for every $z_{E} \in \mathcal{U}_{\rho, E}$. W.l.o.g. we assume $E=\{1, \ldots, J\}$, which always may be achieved by renumerating the variables. If we further define the sequence $\widetilde{\rho}$ by

$$
\widetilde{\rho}_{j}=\rho_{j}+\varepsilon, \quad j \in E, \quad \varepsilon=\frac{\delta}{2 \sum_{j \in E}\left\|\psi_{j} \mid L_{\infty}(D)\right\|}, \quad \widetilde{\rho}_{j}=\rho_{j}, \quad j \notin E,
$$

it is easily checked that $\widetilde{\rho}$ is $\delta / 2$-admissible. In particular, $u_{E}$ is analytic in an open neighbourhood of $\mathcal{U}_{\rho, E}$, where writing $u_{E}\left(z_{1}, \ldots, z_{J}\right)=u_{E}\left(z_{E}\right) \equiv u\left(z_{E}, 0\right)$.
We may thus apply Cauchy's integral formula (see, e.g., [6]) in each variable $z_{j}, j \in E$, to obtain

$$
\begin{aligned}
u_{E}\left(z_{1}, \ldots, z_{J}\right) & =(2 \pi i)^{-J} \oint_{\Gamma_{1}\left(z_{1}\right)} \cdots \oint_{\Gamma_{J}\left(z_{J}\right)} \frac{u\left(z_{E}^{\prime}, 0\right)}{\left(z_{1}^{\prime}-z_{1}\right) \cdots\left(z_{J}^{\prime}-z_{J}\right)} d z_{1}^{\prime} \cdots d z_{J}^{\prime} \\
& =(2 \pi i)^{-J} \oint_{\Gamma_{1}} \cdots \oint_{\Gamma_{J}} \frac{u\left(z_{1}^{\prime}+z_{1}, \ldots, z_{J}^{\prime}+z_{J}\right)}{z_{1}^{\prime} \cdots z_{J}^{\prime}} d z_{1}^{\prime} \cdots d z_{J}^{\prime}
\end{aligned}
$$

where $\Gamma_{j}\left(z_{j}\right)$ denotes the circle with radius $\rho_{j}$ and center $z_{j}$, and $\Gamma_{j} \equiv \Gamma_{J}(0)$. Differentiation (or directly applying the formula for derivatives) then yields

$$
\partial^{\nu} u(0)=\frac{\partial^{|\nu|} u}{\partial z_{1}^{\nu_{1}} \cdots \partial z_{J}^{\nu_{J}}}(0)=\nu!(2 \pi i)^{-J} \oint_{\Gamma_{1}} \cdots \oint_{\Gamma_{J}} \frac{u\left(z_{1}, \ldots, z_{J}\right)}{z_{1}^{\nu_{1}+1} \cdots z_{J}^{\nu_{J}+1}} d z_{1} \cdots d z_{J}
$$

Eventually, together with (4.1) we conclude

$$
\left\|\partial^{\nu} u(0) \mid H_{0}^{1}(D)\right\| \leq \nu!M(\delta) \prod_{j \in E} \rho_{j}^{-\nu_{j}}
$$

This proves the claim.
We put emphasis on the observation that Proposition 4.1 is valid for arbitrary $\delta$-admissible sequences. Hence the resulting estimate may be improved to

$$
\left\|\partial^{\nu} u(0) \mid H_{0}^{1}(D)\right\| \leq \nu!M(\delta) \inf _{\rho \delta \text {-admissible }} \prod_{j \in \operatorname{supp} \nu} \rho_{j}^{-\nu_{j}} .
$$

Moreover, note that the optimal sequence $\rho$ (if it exists) may be a different one for every $\nu \in \mathcal{F}$. We eventually aim at proving $\ell_{p}$-summability results for the Taylor coefficients. For this purpose it will be sufficient to construct "suitable" $\delta$-admissible sequences $\rho \equiv \rho(\nu)$, upon which the respective summability result will be based.

### 4.3 Construction of $\delta / 2$-admissible sequences

The construction is essentially based on corresponding arguments in [2]. The starting point for our construction will be a given fixed $\delta$-admissible sequence $\varrho$. For given multiindex $\nu \in \mathcal{F}$ we will now construct a $\delta / 2$-admissible sequence $\rho$. We shall use the abbreviation $\gamma_{j}=\left\|\psi_{j} \mid L_{\infty}(D)\right\|$. The basic assumption then reads as

$$
\begin{equation*}
\left(\varrho_{j}\left\|\psi_{j} \mid L_{\infty}(D)\right\|\right)_{j \in \mathbb{N}} \in \ell_{1}(\mathbb{N}), \quad \text { i.e. } \quad \sum_{j \geq 1} \varrho_{j} \gamma_{j}<\infty \tag{4.2}
\end{equation*}
$$

To begin with choose $M \in \mathbb{N}$ such that

$$
\sum_{j>M} \varrho_{j} \gamma_{j} \leq \frac{\delta}{12}
$$

which exists due to the choice of $\varrho$. Without loss of generality, we assume that the indexing of the parameters $y_{j}$ is chosen such that the sequence $\left(\varrho_{j} \gamma_{j}\right)_{j \in \mathbb{N}}$ is non-increasing. Then we partition $\mathbb{N}$ into two sets $E=\{1, \ldots, M\}$ and $F=\mathbb{N} \backslash E$. We further choose $\kappa>1$ such that

$$
(\kappa-1) \sum_{j \leq M} \varrho_{j} \gamma_{j} \leq \frac{\delta}{4}
$$

Finally, we define our sequence $\rho$ by

$$
\rho_{j}=\kappa \varrho_{j}, \quad j \in E ; \quad \rho_{j}=\max \left(\varrho_{j}, \frac{\delta \nu_{j}}{4\left|\nu_{F}\right| \gamma_{j}}\right), \quad j \in F,
$$

where $\nu_{E}$ denotes the restriction of $\nu$ to the index set $E$, and $\left|\nu_{F}\right|$ denotes the $\ell_{1}$-norm of the multiindex, i.e. $\left|\nu_{F}\right|=\sum_{j>M} \nu_{j}$. (with the convention $\frac{\nu_{j}}{\left|\nu_{F}\right|}=0$ if $\left|\nu_{F}\right|=0$ ).

Now we first check that this sequence $\rho$ indeed is $\delta / 2$-admissible. To do so, we estimate

$$
\begin{aligned}
\sum_{j \geq 1} \rho_{j}\left|\psi_{j}(x)\right| & =\kappa \sum_{j \leq J} \varrho_{j}\left|\psi_{j}(x)\right|+\sum_{j>J} \max \left(\varrho_{j}, \frac{\delta \nu_{j}}{4\left|\nu_{F}\right| \gamma_{j}}\right)\left|\psi_{j}(x)\right| \\
& \leq(\kappa-1) \sum_{j \leq J} \varrho_{j} \gamma_{j}+\sum_{j \leq J} \varrho_{j}\left|\psi_{j}(x)\right|+\sum_{j>J}\left(\varrho_{j}+\frac{\delta \nu_{j}}{4\left|\nu_{F}\right| \gamma_{j}}\right)\left|\psi_{j}(x)\right| \\
& \leq \frac{\delta}{4}+\sum_{j \in \mathbb{N}} \varrho_{j}\left|\psi_{j}(x)\right|+\frac{\delta}{4},
\end{aligned}
$$

and hence by choice of $\varrho$

$$
\sum_{j \geq 1} \rho_{j}\left|\psi_{j}(x)\right| \leq \sum_{j \in \mathbb{N}} \varrho_{j}\left|\psi_{j}(x)\right|+\frac{\delta}{2} \leq \Re \bar{a}(x)-\frac{\delta}{2} .
$$

Then we obtain by Proposition 4.1 the estimate

$$
\left\|t_{\nu} \mid H_{0}^{1}(D)\right\| \leq M(\delta / 2) \rho^{-\nu} \leq M(\delta / 2)\left(\prod_{j \in E} \eta^{\nu_{j}} \varrho_{j}^{-\nu_{j}}\right)\left(\prod_{j \in F}\left(\frac{\left|\nu_{F}\right| d_{j}}{\nu_{j}}\right)^{\nu_{j}} \varrho_{j}^{-\nu_{j}}\right)
$$

where $\eta=\frac{1}{\kappa}<1$ and $d_{j}=\frac{4 \varrho_{j} \gamma_{j}}{\delta}$. Moreover, factors with exponent $\nu_{j}=0$ are understood to be 1 . This estimate can be equivalently written as

$$
\begin{equation*}
\varrho^{\nu}\left\|t_{\nu} \mid H_{0}^{1}(D)\right\| \leq M(\delta / 2)\left(\prod_{j \in E} \eta^{\nu_{j}}\right)\left(\prod_{j \in F}\left(\frac{\left|\nu_{F}\right| d_{j}}{\nu_{j}}\right)^{\nu_{j}}\right) . \tag{4.3}
\end{equation*}
$$

Finally, we note that the choice of $M$ implies

$$
\begin{equation*}
\left\|d \mid \ell_{1}(F)\right\|=\sum_{j>M} d_{j} \leq \frac{1}{3} \tag{4.4}
\end{equation*}
$$

### 4.4 Summability of Taylor coefficients

The construction in the previous section are the basis for the following theorem. We shall follow closely the argument given in [2] (Sections 2.2, 3.2-3.3 and 4.4). Before stating the main theorem, we mention the following basic result of [1].

Proposition 4.2. Given $0<p<1$, it holds $\left(\frac{|\nu|!}{\nu!} b^{\nu}\right)_{\nu \in \mathcal{F}} \in \ell_{p}(\mathcal{F})$ if, and only if, $\sum_{j \in \mathbb{N}} b_{j}<1$ and $\left(b_{j}\right)_{j \in \mathbb{N}} \in \ell_{p}(\mathbb{N})$.

In the sequel we shall use the following notion: We say that a sequence $\left(\Lambda_{N}\right)_{N \in \mathbb{N}}$ of finite subsets of $\mathcal{F}$ exhausts $\mathcal{F}$, if every finite set $\Lambda \subset \mathcal{F}$ is contained in every set $\Lambda_{N}$ for all $N \geq N_{0}$ (with $N_{0}=N_{0}(\Lambda)$ chosen sufficiently large).

Theorem 4.1. Suppose $f$ satisfies condition (3.10). Moreover, let $\rho=\left(\rho_{j}\right)_{j \in \mathbb{N}}$ be an arbitrary $\delta$-admissible sequence, and assume

$$
\begin{equation*}
\left(\left\|\psi_{j} \mid L_{\infty}(D)\right\|\right)_{j \in \mathbb{N}} \in \ell_{1}(\mathbb{N}) \quad \text { and } \quad\left(\rho_{j}\left\|\psi_{j} \mid L_{\infty}(D)\right\|\right)_{j \in \mathbb{N}} \in \ell_{p}(\mathbb{N}) \quad \text { for some } p<1 \tag{4.5}
\end{equation*}
$$

Then the Taylor coefficients $t_{\nu}$ of the solution $u$ of (2.1) satisfy $\left(\rho^{\nu}\left\|t_{\nu} \mid H_{0}^{1}(D)\right\|\right)_{\nu \in \mathcal{F}} \in \ell_{p}(\mathcal{F})$. If additionally

$$
\begin{equation*}
\sup _{z \in \mathcal{U}_{\rho}}\left\|f(z)-f\left(z_{1}, \ldots, z_{J}, 0, \ldots,\right) \mid H^{-1}(D)\right\| \longrightarrow 0 \quad(J \longrightarrow \infty) \tag{4.6}
\end{equation*}
$$

then it holds

$$
\sum_{\nu \in \mathcal{F}} t_{\nu} z^{\nu}=u(z), \quad z \in \mathcal{U}_{\rho}
$$

with uniform and unconditional convergence, which has to be understood in the following sense: If $\left(\Lambda_{N}\right)_{N \geq 1}$ exhausts $\mathcal{F}$, then the partial sums $S_{\Lambda_{N}} u(z)=\sum_{\nu \in \Lambda_{N}} t_{\nu} z^{\nu}$ satisfy

$$
\lim _{N \rightarrow \infty} \sup _{z \in \mathcal{U}_{\rho}}\left\|u(z)-S_{\Lambda_{N}} u(z) \mid H_{0}^{1}(D)\right\|=0
$$

## Proof. Step 1: Proof of the $\ell_{p}$-summability.

We first note that every $\delta$-admissible sequence may serve as the fixed reference sequence $\varrho$ in the constructions of the last section. Hence, we conclude from (4.3)

$$
\sum_{\nu \in \mathcal{F}}\left(\rho^{\nu}\left\|t_{\nu} \mid H_{0}^{1}(D)\right\|\right)^{p} \leq M(\delta / 2)^{p} \sum_{\nu \in \mathcal{F}}\left(\prod_{j \in E} \eta^{\nu_{j}}\right)^{p}\left(\prod_{j \in F}\left(\frac{\left|\nu_{F}\right| d_{j}}{\nu_{j}}\right)^{\nu_{j}}\right)^{p}
$$

$$
\equiv M(\delta / 2)^{p}\left(\sum_{\nu \in \mathcal{F}_{E}} \alpha(\nu)^{p}\right)\left(\sum_{\nu \in \mathcal{F}_{F}} \beta(\nu)^{p}\right) \equiv M(\delta / 2)^{p} A_{E} A_{F},
$$

where $\mathcal{F}_{E}=\{\nu \in \mathcal{F}: \operatorname{supp} \nu \subset E\}$ and $\mathcal{F}_{F}=\{\nu \in \mathcal{F}: \operatorname{supp} \nu \subset F\}$. Then it further follows

$$
A_{E}=\left(\sum_{\nu \in \mathcal{F}_{E}} \alpha(\nu)^{p}\right)=\sum_{\nu \in \mathcal{F}_{E}} \prod_{j \in E} \eta^{\nu_{j} p}=\prod_{j \in E} \sum_{n=0}^{\infty} \eta^{n p}=\left(\frac{1}{1-\eta^{p}}\right)^{M}
$$

recall $\eta<1$ (since $E$ is finite and $|E|=M$ the index set $\mathcal{F}_{E}$ may be identified with $\mathbb{N}_{0}^{M}$ ). Now we turn to showing $A_{F}<\infty$. Using as before the convention $0^{0}=1$ and $d^{\nu_{F}}=\prod_{j \in F} d_{j}^{\nu_{j}}$ we find

$$
\begin{equation*}
\beta(\nu)=\prod_{j \in F}\left(\frac{\left|\nu_{F}\right| d_{j}}{\nu_{j}}\right)^{\nu_{j}}=\frac{\left|\nu_{F}\right|^{\left|\nu_{F}\right|}}{\prod_{j \in F} \nu_{j}^{\nu_{j}}} d^{\nu_{F}}, \quad \nu \in \mathcal{F}_{F} \tag{4.7}
\end{equation*}
$$

Applying the Stirling inequalities

$$
\frac{n!e^{n}}{e \sqrt{n}} \leq n^{n} \leq \frac{n!e^{n}}{\sqrt{2 \pi n}}, \quad n \geq 1
$$

we can further estimate the numerator and denominator in (4.7),

$$
\left|\nu_{F}\right|^{\left|\nu_{F}\right|} \leq\left|\nu_{F}\right|!e^{\left|\nu_{F}\right|} \quad \text { and } \quad \prod_{j \in F} \nu_{j}^{\nu_{j}} \geq \frac{\nu_{F}!e^{\left|\nu_{F}\right|}}{\prod_{j \in F} \max \left(1, e \sqrt{\nu_{j}}\right)} \geq \frac{\nu_{F}!e^{\left|\nu_{F}\right|}}{\prod_{j \in F} e^{\nu_{j}}}
$$

where at the end we used the bound $\max (1, e \sqrt{n}) \leq e^{n}$. Altogether we then obtain from (4.7)

$$
\beta(\nu) \leq \frac{\left|\nu_{F}\right|!}{\nu_{F}!} e^{\left|\nu_{F}\right|} d^{\nu_{F}}=\frac{\left|\nu_{F}\right|!}{\nu_{F}!} \prod_{j \in F}\left(e d_{j}\right)^{\nu_{j}} .
$$

We next apply Proposition 4.2 to the sequence $\left(e d_{j}\right)_{j \in \mathbb{N}}$. The assumptions of Proposition 4.2 are satisfied due to (4.4), $e<3$, and condition (4.5). Eventually, this yields $\sum_{\nu \in \mathcal{F}_{F}} \beta(\nu)^{p}<\infty$, and thus the asserted summability of the Taylor coefficients.

## Step 2: Convergence of the Taylor series.

We have just shown $\left(\rho^{\nu}\left\|t_{\nu} \mid H_{0}^{1}(D)\right\|\right)_{\nu \in \mathcal{F}} \in \ell_{p}(\mathcal{F})$, which particularly implies $\left(\rho^{\nu}\left\|t_{\nu} \mid H_{0}^{1}(D)\right\|\right)_{\nu \in \mathcal{F}} \in$ $\ell_{1}(\mathcal{F})$ and thus the absolute convergence of the Taylor series $\sum_{\nu \in \mathcal{F}} t_{\nu} z^{\nu}$ on $\mathcal{U}_{\rho}$. We now show its convergence towards $u(z)$ in two substeps: First we reduce it to the case of a finite number of parameters $z_{j}$, and then prove the convergence in that special case.

## Substep 2.1: Reduction to a finite number of parameters.

If we put

$$
a_{J}\left(\cdot, z_{1}, \ldots, z_{J}\right)=\bar{a}(\cdot)+\sum_{1 \leq j \leq J} z_{j} \psi_{j}(\cdot),
$$

we obtain

$$
\left\|a(\cdot, z)-a_{J}\left(\cdot, z_{1}, \ldots, z_{J}\right)\left|L_{\infty}(D)\|=\| \sum_{j>J} z_{j} \psi_{j}\right| L_{\infty}(D)\right\| \leq \sum_{j>J}\left|z_{j}\right| \cdot\left\|\psi_{j} \mid L_{\infty}(D)\right\|
$$

and thus

$$
\begin{equation*}
\sup _{z \in \mathcal{U}_{\rho}}\left\|a(\cdot, z)-a_{J}\left(\cdot, z_{1}, \ldots, z_{J}\right)\left|L_{\infty}(D)\left\|\leq \sum_{j>J} \rho_{j}\right\| \psi_{j}\right| L_{\infty}(D)\right\| \longrightarrow 0 \quad(J \longrightarrow \infty) \tag{4.8}
\end{equation*}
$$

due to condition (4.5). If we denote by $u_{J}$ the restriction of $u$ to $J$ parameters, i.e.

$$
u_{J}\left(z_{1}, \ldots, z_{J}\right)=u\left(z_{1}, \ldots, z_{J}, 0, \ldots\right) \equiv u\left(z_{(J)}\right),
$$

then it follows from Proposition 3.2 and $\delta$-admissibility of $\rho$

$$
\begin{align*}
&\left\|u(z)-u_{J}\left(z_{1}, \ldots, z_{J}\right) \mid H_{0}^{1}(D)\right\| \leq \frac{2}{\delta}\left(\left\|f(z)-f\left(z_{1}, \ldots, z_{J}, 0, \ldots,\right) \mid H^{-1}(D)\right\|\right. \\
&\left.+M(\delta)\left\|a(\cdot, z)-a_{J}\left(\cdot, z_{1}, \ldots, z_{J}\right) \mid L_{\infty}(D)\right\|\right) \tag{4.9}
\end{align*}
$$

We finally conclude from (4.8) and (4.9) together with the assumption (4.6)

$$
\sup _{z \in \mathcal{U}_{\rho}}\left\|u(z)-u_{J}\left(z_{1}, \ldots, z_{J}\right) \mid H_{0}^{1}(D)\right\| \longrightarrow 0 \quad(J \longrightarrow \infty) .
$$

In particular, we have

$$
\sup _{z \in \mathcal{U}_{\rho}}\left\|u(z)-u_{J}\left(z_{1}, \ldots, z_{J}\right) \mid H_{0}^{1}(D)\right\| \leq \frac{1}{2 N}
$$

for every $J \geq J_{0}(N)$ and $N \in \mathbb{N}$.

## Substep 2.2: Convergence for finitely many parameters.

As mentioned before, the mapping $z_{(J)} \mapsto u\left(z_{(J)}\right)$ is holomorphic in an open neighbourhood of the polydisc $\mathcal{U}_{\rho, J}=\bigotimes_{1 \leq j \leq J}\left\{\zeta \in \mathbb{C}:|\zeta| \leq \rho_{j}\right\}$, see the proof of Proposition 4.1. By standard results on Banach space-valued holomorphic functions (cf. [4, Proposition 3.5] or [6, Theorem 2.1.2]) this implies its analyticity and therefore the uniform convergence of its Taylor series on this polydisc. In particular, for the index sets

$$
\Lambda_{K}^{*}=\left\{\mu \in \mathbb{N}_{0}^{J}:|\mu| \leq K\right\}
$$

we have

$$
\sup _{z=\left(z_{1}, \ldots, z_{J}\right) \in \mathcal{U}_{\rho, J}}\left\|u_{J}(z)-S_{\Lambda_{K}^{*}} u_{J}(z) \mid H_{0}^{1}(D)\right\| \leq \frac{1}{2 N}
$$

for all $K \geq K_{0}(N)$. W.l.o.g. we choose $K_{0}(N) \geq N$ and $J_{0}(N) \geq N$ to be increasing. Moreover, the index set $\Lambda_{K}^{*} \subset \mathbb{N}_{0}^{J}$ can be identified with an index set $\bar{\Lambda}_{K}^{*} \subset \mathcal{F}$ in an obvious way (the contained indices being supported in $\{1, \ldots, J\}$ ). In the same spirit we identify the partial sums $S_{\Lambda_{K}^{*}} u_{J}\left(z_{1}, \ldots, z_{J}\right)$ and $S_{\bar{\Lambda}_{K}^{*}} u\left(z_{(J)}\right) \equiv S_{\bar{\Lambda}_{K}^{*}} u(z)$ for $z \in \mathcal{U}_{\rho}$. Finally, putting $\Lambda_{N}^{*}=\Lambda_{K_{0}(N)} \subset \mathbb{N}_{0}^{J_{0}(N)}$, we end up with

$$
\sup _{z \in \mathcal{U}_{\rho}}\left\|u(z)-S_{\bar{\Lambda}_{N}^{*}} u(z) \mid H_{0}^{1}(D)\right\| \leq \frac{1}{N},
$$

which proves the (uniform) convergence of the Taylor series for this particular order of summation. Clearly the sequence of sets $\left(\bar{\Lambda}_{N}^{*}\right)_{N \in \mathbb{N}}$ exhausts $\mathcal{F}$.

## Substep 3: Unconditional summability.

Now let an arbitrary exhausting sequence $\left(\Lambda_{N}\right)_{N \in \mathbb{N}}$ of finite subsets of $\mathcal{F}$ be given. Then for any fixed $\varepsilon>0$ we can find some $M=M(\varepsilon)$, such that

$$
\sup _{z \in \mathcal{U}_{\rho}}\left\|u(z)-S_{\bar{\Lambda}_{M}^{*}} u(z) \mid H_{0}^{1}(D)\right\| \leq \frac{\varepsilon}{2} .
$$

Moreover, we may assume

$$
\sum_{\nu \notin \bar{\Lambda}_{M}^{*}} \rho^{\nu}\left\|t_{\nu} \mid H_{0}^{1}(D)\right\| \leq \frac{\varepsilon}{2}
$$

since $\left(\bar{\Lambda}_{N}^{*}\right)_{N \in \mathbb{N}}$ exhausts $\mathcal{F}$. Because $\left(\Lambda_{N}\right)_{N \in \mathbb{N}}$ exhausts $\mathcal{F}$ there exists some $N_{0}(\varepsilon)$ such that $\bar{\Lambda}_{M}^{*} \subset \Lambda_{N}$ for all $N \geq N_{0}$. This finally implies

$$
\begin{aligned}
\sup _{z \in \mathcal{U}_{\rho}}\left\|u(z)-S_{\Lambda_{N}} u(z) \mid H_{0}^{1}(D)\right\| & \leq \sup _{z \in \mathcal{U}_{\rho}}\left\|u(z)-S_{\bar{\Lambda}_{M}^{*}} u(z)\left|H_{0}^{1}(D)\left\|+\sup _{z \in \mathcal{U}_{\rho}}\right\| S_{\Lambda_{N} \backslash \bar{\Lambda}_{M}^{*}} u(z)\right| H_{0}^{1}(D)\right\| \\
& \leq \frac{\varepsilon}{2}+\sup _{z \in \mathcal{U}_{\rho}} \sum_{\nu \in \Lambda_{N} \backslash \bar{\Lambda}_{M}^{*}}\left|z^{\nu}\right| \cdot\left\|t_{\nu} \mid H_{0}^{1}(D)\right\| \\
& \leq \frac{\varepsilon}{2}+\sum_{\nu \notin \bar{\Lambda}_{M}^{*}} \rho^{\nu}\left\|t_{\nu} \mid H_{0}^{1}(D)\right\| \leq \varepsilon .
\end{aligned}
$$

This completes the proof.

## 5 Series expansions in tensorized polynomial bases

### 5.1 Legendre Series

In this section we shall consider expansions of the mapping $z \mapsto u(z)$ into series of tensorized Legendre polynomials. These expansions and related $N$-term approximations are better suited than Taylor polynomials, if the error is measured in the least squares sense, i.e. it results in better decay estimates.
When talking about Legendre polynomials there are different versions according to the chosen normalization. In the univariate case, we define the system $\left(P_{n}\right)_{n \geq 0}$ to be Legendre polynomials with the $L_{\infty}$-normalization $\left\|P_{n} \mid L_{\infty}([-1,1])\right\|=P_{n}(1)=1$, and we denote by $L_{n}=\sqrt{2 n+1} P_{n}$ their $L_{2}$-normalized version, i.e.

$$
\int_{-1}^{1}\left|L_{n}(t)\right|^{2} \frac{d t}{2}=1
$$

We further put $P_{0}=L_{0}=1$. For $\nu \in \mathcal{F}$, we define

$$
P_{\nu}(z)=\prod_{j \geq 1} P_{\nu_{j}}\left(z_{j}\right) \quad \text { and } \quad L_{\nu}(z)=\prod_{j \geq 1} L_{\nu_{j}}\left(z_{j}\right) .
$$

We note that the choice $P_{0}=L_{0}=1$ renders $P_{\nu}(z)$ well-defined, since due to the finite support of $\nu$, only finitely many factors are different from 1.
As a direct consequence we note that $\left(L_{\nu}\right)_{\nu \in \mathcal{F}}$ is an orthonormal basis in $L_{2}(U, d \mu)$, where $d \mu$ is the countable product of the probability measures $\frac{d y_{j}}{2}$ on $[-1,1]$. The spaces $L_{p}(U, d \mu), 0<p \leq \infty$, as well as the vector-valued spaces $L_{p}(U, d \mu ; X)$ are to be understood similarly.
From condition (3.4) and the a priori estimate (3.5) we immediately conclude $u \in L_{\infty}\left(U, d \mu ; H_{0}^{1}(D)\right) \hookrightarrow$ $L_{2}\left(U, d \mu ; H_{0}^{1}(D)\right)$, thus we obtain the unique expansions

$$
u(y)=\sum_{\nu \in \mathcal{F}} u_{\nu} P_{\nu}(y)=\sum_{\nu \in \mathcal{F}} v_{\nu} L_{\nu}(y)
$$

with convergence in $L_{2}\left(U, d \mu ; H_{0}^{1}(D)\right)$, where the $H_{0}^{1}(D)$-valued coefficients $u_{\nu}$ and $v_{\nu}$ are given by

$$
v_{\nu}=\int_{U} u(y) L_{\nu}(y) d \mu(y) \quad \text { and } \quad u_{\nu}=\left(\prod_{j \geq 1}\left(1+2 \nu_{j}\right)\right)^{1 / 2} v_{\nu}
$$

We then find the following analog of Theorem 4.1 for tensorized Legendre expansions.
Theorem 5.1. Let $a$ fulfill Assumption 1 for some $0<r \leq R<\infty$, and suppose $f$ satisfies condition (3.10) and Assumption 3. Moreover, let $\rho=\left(\rho_{j}\right)_{j \in \mathbb{N}}$ be a $\delta$-admissible sequence with $\rho_{j} \geq 1, j \in \mathbb{N}$, and assume

$$
\begin{equation*}
\left(\left\|\psi_{j} \mid L_{\infty}(D)\right\|\right)_{j \in \mathbb{N}} \in \ell_{1}(\mathbb{N}) \quad \text { and } \quad\left(\rho_{j}\left\|\psi_{j} \mid L_{\infty}(D)\right\|\right)_{j \in \mathbb{N}} \in \ell_{p}(\mathbb{N}) \quad \text { for some } p<1 \tag{5.1}
\end{equation*}
$$

Then the Legendre coefficients $u_{\nu}$ and $v_{\nu}$ of the solution $u$ of (2.1) satisfy $\left(\rho^{\nu}\left\|u_{\nu} \mid H_{0}^{1}(D)\right\|\right)_{\nu \in \mathcal{F}} \in \ell_{p}(\mathcal{F})$ and $\left(\rho^{\nu}\left\|v_{\nu} \mid H_{0}^{1}(D)\right\|\right)_{\nu \in \mathcal{F}} \in \ell_{p}(\mathcal{F})$. If additionally

$$
\begin{equation*}
\sup _{z \in \mathcal{U}}\left\|f(z)-f\left(z_{1}, \ldots, z_{J}, 0, \ldots,\right) \mid H^{-1}(D)\right\| \longrightarrow 0 \quad(J \longrightarrow \infty) \tag{5.2}
\end{equation*}
$$

then it holds

$$
u(z)=\sum_{\nu \in \mathcal{F}} u_{\nu} P_{\nu}(z)=\sum_{\nu \in \mathcal{F}} v_{\nu} L_{\nu}(z), \quad z \in \mathcal{U}
$$

with uniform and unconditional convergence, which has to be understood in the following sense: If $\left(\Lambda_{N}\right)_{N \geq 1}$ is a sequence of subsets of $\mathcal{F}$ which exhausts $\mathcal{F}$, then the partial sums $S_{\Lambda_{N}} u(z)=$ $\sum_{\nu \in \Lambda_{N}} u_{\nu} P_{\nu}(z)=\sum_{\nu \in \Lambda_{N}} v_{\nu} L_{\nu}(z)$ satisfy

$$
\lim _{N \rightarrow \infty} \sup _{z \in \mathcal{U}}\left\|u(z)-S_{\Lambda_{N}} u(z) \mid H_{0}^{1}(D)\right\|=0
$$

As for the Taylor series, we start with a preliminary estimate.
Proposition 5.1. Let the coefficient $a$ satisfy Assumption 1. For any $\nu \in \mathcal{F}$ let $\rho$ denote a $\nu$ dependent $\delta$-admissible sequence which is chosen in the following fashion: $\rho_{j} \geq 1$ for all $j \in \mathbb{N}$ and $\rho_{j}>1$ for all $j \in \operatorname{supp} \nu$. Moreover, suppose $f$ fulfills condition (3.10) and Assumption 3. Then we have the estimate

$$
\left\|u_{\nu} \mid H_{0}^{1}(D)\right\| \leq M(\delta) \prod_{j \in \operatorname{supp} \nu} \frac{\pi \rho_{j}}{2\left(\rho_{j}-1\right)}\left(2 \nu_{j}+1\right) \rho_{j}^{-\nu_{j}} .
$$

The proof is exactly the same as the corresponding one in [2, Lemma 4.2], based on the a priori estimate (3.11). Equipped with this estimate, the proof of Theorem 5.1 follows by exactly the same argument as in [2], Sections 4.3 and 4.4. Specifically, this means a modification of the argument for Taylor coefficients and Taylor series (i.e. construction of suitable $\delta / 2$-admissible sequences $\rho(\nu)$, reformulation of the estimate resulting from Proposition 5.1, and finally applying Proposition 4.2). We omit the details.

### 5.2 Chebyshev Series

As another example of orthogonal polynomials defined on the interval $[-1,1]$ we consider the Chebyshev polynomials $T_{n}(t)=\cos (n \arccos (t)), n \in \mathbb{N}$, and $T_{0}(t)=1$. They satisfy

$$
\left\|T_{n} \mid L_{\infty}([-1,1])\right\|=1, \quad n \geq 0, \quad \text { and } \quad \int_{[-1,1]}\left|T_{n}(t)\right|^{2} \frac{d t}{\sqrt{1-t^{2}}}=\frac{\pi}{2}, \quad n \in \mathbb{N}
$$

As before we consider the set $U$ equipped with the Borel $\sigma$-algebra $\mathcal{B}(U)$ and the probability measure

$$
d \eta=\bigotimes_{j \in \mathbb{N}} \frac{d t}{\pi \sqrt{1-t^{2}}}
$$

As in the Legendre-case condition (3.4) and the a priori estimate (3.5) imply $u \in L_{\infty}\left(U, d \eta ; H_{0}^{1}(D)\right) \hookrightarrow$ $L_{2}\left(U, d \eta ; H_{0}^{1}(D)\right)$. Moreover, the system of tensorized polynomials $\left(T_{\nu}\right)_{\nu \in \mathcal{F}}$, where for $\nu \in \mathcal{F}$ we put $T_{\nu}(y)=\prod_{j \in \mathbb{N}} T_{\nu_{j}}\left(y_{j}\right)$, constitutes an orthogonal basis. Note that they are not orthonormal with respect to the measure $\eta$, but it holds

$$
\int_{U}\left|T_{\nu}(y)\right|^{2} d \eta(y)=\prod_{j \in \operatorname{supp} \nu} \frac{1}{2}=2^{-|\operatorname{supp} \nu|}
$$

Nevertheless, this yields the unique expansions

$$
u(y)=\sum_{\nu \in \mathcal{F}} w_{\nu} T_{\nu}(y), \quad y \in U, \quad w_{\nu}=2^{|\operatorname{supp} \nu|} \int_{U} u(y) T_{\nu}(y) d \eta(y)
$$

with convergence in $L_{2}\left(U, d \eta ; H_{0}^{1}(D)\right)$.
We again aim at a summability result as in Theorems 4.1 and 5.1, and once more we start with an estimate in terms of arbitrary $\delta$-admissible sequences.

Proposition 5.2. Under the assumptions of Proposition 5.1 it holds the estimate

$$
\left\|w_{\nu} \mid H_{0}^{1}(D)\right\| \leq M(\delta) \prod_{j \in \operatorname{supp} \nu} 2 \rho_{j}^{-\nu_{j}} \equiv M(\delta) 2^{|\operatorname{supp} \nu|} \rho^{-\nu}
$$

Proof. We will only discuss the case $\nu=n e^{1}, n \in \mathbb{N}$, the general case then is a straightforward modification by applying the one-variable case to every variable $z_{j}$ with $j \in \operatorname{supp} \nu$.
Similar to the proof of Proposition 4.1 we write $z=\left(z_{1}, z^{\prime}\right) \in \mathcal{A}_{\delta}$ and put $u_{1}\left(z_{1}\right)=u\left(z_{1}, z^{\prime}\right)$ for some arbitrary $z^{\prime} \in \mathcal{U}^{\prime}=\prod_{j \geq 2}\{\zeta \in \mathbb{C}:|\zeta| \leq 1\}$ (note that due to the assumption on $\rho$ we have $\mathcal{U} \subset \mathcal{U}_{\rho}$ ). Then we have $\left\|u_{1}\left(z_{1}\right) \mid H_{0}^{1}(D)\right\| \leq M(\delta)$, and we further find

$$
\int_{U} T_{\nu}(y) u(y) d \eta(y)=\int_{U^{\prime}} \int_{[-1,1]} T_{n}(t) u\left(t, y^{\prime}\right) \frac{d t}{\pi \sqrt{1-t^{2}}} d \eta\left(y^{\prime}\right) \equiv \int_{U^{\prime}} \mathcal{I}_{n}\left(y^{\prime}\right) d \eta\left(y^{\prime}\right)
$$

It clearly suffices to bound $\left\|\mathcal{I}_{n}\left(y^{\prime}\right) \mid H_{0}^{1}(D)\right\|$ independently of $y^{\prime}$. At first we find

$$
\begin{aligned}
\pi \mathcal{I}_{n}\left(y^{\prime}\right) & =\int_{0}^{\pi} \cos (n \theta) u_{1}(\cos \theta) d \theta=\frac{1}{2} \int_{-\pi}^{\pi} u_{1}(\cos \theta) \cos (n \theta) d \theta \\
& =\frac{1}{2 i} \int_{|\zeta|=1} u_{1}\left(\frac{\zeta+\zeta^{-1}}{2}\right)\left(\frac{\zeta^{n}+\zeta^{-n}}{2}\right) \frac{d \zeta}{\zeta} \equiv \frac{1}{2 i} \int_{|\zeta|=1} u_{1}(\mathcal{J}(\zeta))\left(\mathcal{J}\left(\zeta^{n}\right)\right) \frac{d \zeta}{\zeta} .
\end{aligned}
$$

The last step is verified by substituting the standard parametrization $\zeta(\theta)=e^{i \theta}$ and by the Joukowskytransform $\mathcal{J}(\zeta)=\left(\zeta+\zeta^{-1}\right) / 2$. Vice versa it is well-known that $\mathcal{J}$ maps the unit circle onto the interval $[-1,1]$ (traversed twice), since with $|\zeta|=1$ it follows $\mathcal{J}(\zeta)=\Re(\zeta)$. Then we can estimate

$$
4 \pi\left\|\mathcal{I}_{n}\left(y^{\prime}\right)\left|H_{0}^{1}(D)\|\leq\| \int_{|\zeta|=1} u_{1}(\mathcal{J}(\zeta)) \zeta^{n} \frac{d \zeta}{\zeta}\right| H_{0}^{1}(D)\right\|+\left\|\left.\int_{|\zeta|=1} u_{1}(\mathcal{J}(\zeta)) \zeta^{-n} \frac{d \zeta}{\zeta} \right\rvert\, H_{0}^{1}(D)\right\|
$$

$$
\begin{aligned}
& =\left\|\int_{|\zeta|=\rho_{1}} u_{1}(\mathcal{J}(\zeta)) \zeta^{n-1} d \zeta\left|H_{0}^{1}(D)\|+\| \int_{|\zeta|=\rho_{1}^{-1}} u_{1}(\mathcal{J}(\zeta)) \zeta^{-n-1} d \zeta\right| H_{0}^{1}(D)\right\| \\
& \leq M(\delta) \rho_{1}^{n-1} \cdot 2 \pi \rho_{1}+M(\delta) \rho_{1}^{-n-1} \cdot 2 \pi \rho_{1}^{-1}=4 \pi M(\delta) \rho_{1}^{-n}
\end{aligned}
$$

Here we used Cauchy's Theorem, since $|\mathcal{J}(\zeta)| \leq \rho_{1}$ for $|\zeta|=\rho_{1}$ or $|\zeta|=\rho_{1}^{-1}$, and $u_{1}$ is analytic in an open neighbourhood of the disc $\left\{\zeta \in \mathbb{C}:|\zeta| \leq \rho_{1}\right\}$, see the proof of Proposition 4.1.

The above classical argument can essentially be found in Section 3 of [11].
By a modification of the construction and summation argument from Section 4 similar to the Legendre-case, we then obtain an analogous summability result for coefficients in the Chebyshev expansion.

Theorem 5.2. Theorem 5.1 remains true upon replacing the Legendre polynomials $L_{\nu}$ by the Chebyshev polynomials $T_{\nu}$ as well as the Legendre coefficients $v_{\nu}$ by the Chebyshev coefficients $w_{\nu}$.

## 6 Higher spatial regularity

In connection with numerical approximations of the problem (2.6) in the domain $D$, it is of interest to establish sufficient conditions on the data which imply additional regularity of the parametric solutions in $D$, while maintainig the summability of their series expansions. In this section, we assume $f \in L_{2}(D) \hookrightarrow H^{-1}(D)$. Then we have

$$
\left\|f\left|H^{-1}(D)\left\|\leq C_{P}\right\| f\right| L_{2}(D)\right\|
$$

where $C_{P}$ denotes the Poincaré constant of $D$. Here, the smoothness space $W(D)$ is defined as space of all solutions to the Dirichlet problem

$$
-\Delta u=f \quad \text { in } D,\left.\quad u\right|_{\partial D}=0
$$

with $f \in L_{2}(D)$, i.e.

$$
W(D)=\left\{u \in H_{0}^{1}(D): \Delta u \in L_{2}(D)\right\} .
$$

The space $W(D)$ can be semi-normed and normed, respectively, by

$$
|v|_{W}=\left\|\Delta v\left|L_{2}(D)\|, \quad\| v\right| W(D)\right\|=\left\|v \left|H_{0}^{1}(D) \|+|v|_{W} .\right.\right.
$$

We remark that for polygonal resp. polyhedral domains $D$, this space is contained in certain weighted $H^{2}(D)$ spaces, with weight functions which vanish at corners resp. corners and vertices of $D$ (see, e.g., $[5,9]$ and the references there).

## 6.1 $W$-analyticity of $u$

Before we return to our parametric problem we shall first prove an a priori estimate and a stability assertion for the problem (2.1), now with respect to the $W$-norm. If we additionally assume $a \in$ $W_{\infty}^{1}(D)$, i.e. $\nabla a \in L_{\infty}(D)$, then we may reformulate the problem (2.1) as

$$
\begin{equation*}
-\Delta u=\frac{1}{a}(f+\nabla a \cdot \nabla u-G(u)) . \tag{6.1}
\end{equation*}
$$

Assuming as before the validity of the ellipticity condition (2.5) we need to assure $G(u) \in L_{2}(D)$ in order to conclude $\Delta u \in L_{2}(D)$. Restricting ourselves once more to functions $G(t)=t^{m}$, this further restricts the choices of possible pairs $(n, m)$. More precisely, we define

$$
\mathcal{M}^{\prime}=\left\{(n, m) \in \mathcal{M}: H_{0}^{1}(D) \hookrightarrow L_{2 m}(D)\right\} .
$$

Under these conditions we can conclude: If $u \in H_{0}^{1}(D)$ solves the problem (2.2), then we already have $\Delta u \in L_{2}(D)$, and from (6.1) we obtain the estimate

$$
\begin{align*}
|u|_{W}=\left\|\Delta u \mid L_{2}(D)\right\| & \leq \frac{1}{r}\left(\left\|f\left|L_{2}(D)\|+\| \nabla a\right| L_{\infty}(D)\right\| \cdot\left\|u\left|H_{0}^{1}(D)\|+\| G(u)\right| L_{2}(D)\right\|\right) \\
& \leq \frac{1}{r}\left(\left\|f\left|L_{2}(D)\|+\| \nabla a\right| L_{\infty}(D)\right\| \cdot \frac{2\left\|f \mid H^{-1}(D)\right\|}{r}+\left\|u \mid L_{2 m}(D)\right\|\right), \tag{6.2}
\end{align*}
$$

and hence

$$
\begin{equation*}
\left\|u\left|W(D)\left\|\leq \frac{1}{r}\left(1+C_{P}\left(2+\frac{2\left\|\nabla a \mid L_{\infty}(D)\right\|}{r}+\frac{2 \iota(2 m)}{r}\right)\right)\right\| f\right| L_{2}(D)\right\| \tag{6.3}
\end{equation*}
$$

We continue with the announced stability assertion.
Lemma 6.1. Let $u$ and $\widetilde{u}$ be solutions of (2.2) with $f, \tilde{f} \in L_{2}(D)$ and coefficients $a, \widetilde{a} \in W_{\infty}^{1}(D)$, satisfying condition (2.5) for the same constants $r$ and $R$. Then there holds the estimate

$$
\begin{aligned}
|u-\widetilde{u}|_{W} \leq & \frac{1}{r} \\
( & \left(\left\|f-\widetilde{f}\left|L_{2}(D)\|+\| a-\widetilde{a}\right| L_{\infty}(D)\right\| \cdot|u|_{W}\right. \\
& +\left\|\nabla(a-\widetilde{a})\left|L_{\infty}(D)\|\cdot\| u\right| H_{0}^{1}(D)\right\|+\left\|\nabla \widetilde{a}\left|L_{\infty}(D)\|\cdot\| u-\widetilde{u}\right| H_{0}^{1}(D)\right\| \\
& \left.+\iota(2 m)^{2 m} C_{P}^{2(m-1)}\left\|u-\widetilde{u} \mid H_{0}^{1}(D)\right\|^{2}\left(\frac{2 \max \left(\left\|f\left|L_{2}(D)\|,\| \tilde{f}\right| L_{2}(D)\right\|\right)}{r}\right)^{2(m-1)}\right) .
\end{aligned}
$$

Proof. We already know from the above considerations $u, \widetilde{u} \in W(D)$. Moreover, from

$$
-a \Delta u=f+\nabla a \cdot \nabla u-G(u) \quad \text { and } \quad-\widetilde{a} \Delta \widetilde{u}=\widetilde{f}+\nabla \widetilde{a} \cdot \nabla \widetilde{u}-G(\widetilde{u})
$$

we conclude

$$
\widetilde{a} \Delta(\widetilde{u}-u)=(a-\widetilde{a}) \Delta u+f-\widetilde{f}+\nabla(a-\widetilde{a}) \cdot \nabla u+\nabla \widetilde{a} \cdot \nabla(u-\widetilde{u})+G(\widetilde{u})-G(u) .
$$

Taking into account condition (2.5) for $\widetilde{a}$, the claim now follows by applying the $L_{2}(D)$-norm. However, we need to have a closer look on the term $G(\widetilde{u})-G(u)$. Applying Hölder's inequality we obtain

$$
\begin{aligned}
\left\|G(u)-G(\widetilde{u}) \mid L_{2}(D)\right\| & =\int_{D}\left|u(x)^{m}-\widetilde{u}(x)^{m}\right|^{2} d x \\
& =\int_{D}\left|u(x)^{m}-\widetilde{u}(x)^{m}\right| \cdot|u(x)-\widetilde{u}(x)| \cdot\left|\sum_{j=0}^{m-1} u(x)^{m-1-j} \widetilde{u}(x)^{j}\right| d x \\
& \leq\left\|u-\widetilde{u}\left|L_{2 m}(D)\|\cdot\| u^{m}-\widetilde{u}^{m}\right| L_{2 m}(D)\right\| \sum_{j=0}^{m-1}\left\|u\left|L_{2 m}(D)\left\|^{m-1-j}\right\| \widetilde{u}\right| L_{2 m}(D)\right\|^{j}
\end{aligned}
$$

$$
\begin{aligned}
& \leq\left\|u-\widetilde{u} \mid L_{2 m}(D)\right\|^{2}\left(\sum_{j=0}^{m-1}\left\|u\left|L_{2 m}(D)\left\|^{m-1-j}\right\| \widetilde{u}\right| L_{2 m}(D)\right\|^{j}\right)^{2} \\
& \leq \iota(2 m)^{2 m}\left\|u-\widetilde{u} \mid H_{0}^{1}(D)\right\|^{2}\left(\sum_{j=0}^{m-1}\left\|u\left|H_{0}^{1}(D)\left\|^{m-1-j}\right\| \widetilde{u}\right| H_{0}^{1}(D)\right\|^{j}\right)^{2} \\
& \leq \iota(2 m)^{2 m} C_{P}^{2(m-1)}\left\|u-\widetilde{u} \mid H_{0}^{1}(D)\right\|^{2}\left(\frac{2 \max \left(\left\|f\left|L_{2}(D)\|,\| \widetilde{f}\right| L_{2}(D)\right\|\right)}{r}\right)^{2(m-1)} .
\end{aligned}
$$

This finally completes the proof.
We now return to the parametric problem. In order to derive a counterpart of Theorem 3.3, we shall henceforth assume $\psi_{j} \in W_{\infty}^{1}(D)$ for all $j$, and we define domains

$$
\mathcal{A}_{\delta, B}:=\left\{z \in \mathbb{C}^{\mathbb{N}}: 0<\delta \leq \Re(a(x, z)) \leq|a(x, z)| \leq 2 R \text { and }|\nabla a(x, z)| \leq B \text { for a.e. } x \in D\right\} \subset \mathcal{A}_{\delta}
$$

Clearly, choosing $B$ sufficiently large implies that $\mathcal{A}_{\delta, B}$ is nonempty. Moreover, the estimate (6.3) implies the uniform bound

$$
\begin{equation*}
\left\|u(z)\left|W(D)\left\|\leq M(\delta, B)=\frac{1}{\delta}\left(1+2 C_{P}\left(1+\frac{B}{\delta}+\frac{\iota(2 m)}{\delta}\right)\right) \sup _{z \in \mathcal{A}_{\delta, B}}\right\| f(z)\right| L_{2}(D)\right\| . \tag{6.4}
\end{equation*}
$$

Finally, we put $\widetilde{\mathcal{A}}=\bigcup_{\delta>0, B>0} \mathcal{A}_{\delta, B}$.
Proposition 6.1. Let $(n, m) \in \mathcal{M}^{\prime}$, and consider $G(\zeta)=\zeta^{m}$ with $m \geq 2$. Then it holds

$$
\lim _{h \rightarrow 0} \frac{\left\|G\left(u\left(z+h e^{j}\right)\right)-G(u(z))-G^{\prime}(u(z))\left(u\left(z+h e^{j}\right)-u(z)\right) \mid L_{2}(D)\right\|}{|h|}=0, \quad z \in \widetilde{\mathcal{A}} .
$$

Proof. Fix $z \in \mathcal{A}_{\delta, B}$. Then similar to the proof of Proposition 6.1 we find $z+h e^{j} \in \mathcal{A}_{\delta / 2,2 B}$ for $|h|$ sufficiently small, i.e. $u\left(z+h e^{j}\right) \in H_{0}^{1}(D)$ is well-defined. As in the proof of Proposition 3.3 we apply Taylor's Theorem to obtain

$$
\begin{aligned}
& \left\|G\left(u\left(z+h e^{j}\right)\right)-G(u(z))-G^{\prime}(u(z))\left(u\left(z+h e^{j}\right)-u(z)\right) \mid L_{2}(D)\right\|^{2} \\
& \quad=\int_{D}\left|u\left(z+h e^{j}\right)-u(z)\right|^{4}\left|\int_{0}^{1}(1-t) G^{\prime \prime}\left(u(z)+t\left(u\left(z+h e^{j}\right)-u(z)\right)\right) d t\right|^{2} d x \\
& \quad \leq\left\|u\left(z+h e^{j}\right)-u(z)\left|L_{2 m}(D)\left\|^{4} \cdot\right\| \int_{0}^{1}(1-t) G^{\prime \prime}\left(u(z)+t\left(u\left(z+h e^{j}\right)-u(y)\right)\right) d t\right| L_{\frac{2 m}{m-2}}(D)\right\|^{2},
\end{aligned}
$$

where in the last line we used Hölder's inequality (again the case $m=2$ is trivial). As before the last factor can be further estimated using the (generalized) Minkowski inequality to find

$$
\begin{aligned}
& \left\|\int_{0}^{1}(1-t) G^{\prime \prime}\left(u(z)+t\left(u\left(z+h e^{j}\right)-u(z)\right)\right) d t \left\lvert\, L_{\frac{2 m}{m-2}}(D)\right.\right\| \\
& \quad \leq \int_{0}^{1}(1-t) m(m-1)\left\|\left|u(z)+t\left(u\left(z+h e^{j}\right)-u(z)\right)\right|^{m-2} \left\lvert\, L_{\frac{2 m}{m-2}}(D)\right.\right\| d t
\end{aligned}
$$

$$
\leq \frac{m(m-1)}{2}(\iota(2 m) M(\delta, B))^{m-2}
$$

Combining these two estimates finally yields

$$
\begin{aligned}
& \left\|G\left(u\left(z+h e^{j}\right)\right)-G(u(z))-G^{\prime}(u(z))\left(u\left(z+h e^{j}\right)-u(z)\right) \mid L_{2}(D)\right\| \\
& \quad \leq \frac{1}{2} m(m-1) \iota(2 m)^{m}\left\|u\left(z+h e^{j}\right)-u(z) \mid H_{0}^{1}(D)\right\|^{2} \cdot M(\delta, B)^{m-2}=\mathcal{O}\left(|h|^{2}\right),
\end{aligned}
$$

cf. either Proposition 3.2 or Theorem 3.3. This proves the claim.
As before, we still need an assumption on the regularity of $z \mapsto f(z)$.
Assumption 4. We assume the mapping $z \mapsto f(z)$ to admit partial derivatives $\partial_{z_{j}} f(z) \in L_{2}(D)$ at every point $z \in \widetilde{\mathcal{A}}$ with respect to every variable $z_{j}, j \in \mathbb{N}$.

Theorem 6.1. For every $z \in \widetilde{\mathcal{A}}$ the mapping $z \mapsto u(z)$ admits complex partial derivatives $\partial_{z_{j}} u(z) \in$ $W(D)$ with respect to each variable $z_{j}$.

Proof. From the proof of Theorem 3.3 we know that the difference quotient $w_{h}(z)$ is the unique weak solution of

$$
-\nabla(a(\cdot, z) \nabla w)+G^{\prime}(u(z)) w=\partial_{z_{j}} f(z)+L_{h},
$$

where

$$
\begin{aligned}
L_{h}= & \frac{f\left(z+h e^{j}\right)-f(z)-h \partial_{z_{j}} f(z)}{h}+\nabla\left(\psi_{j} \nabla u\left(z+h e^{j}\right)\right) \\
& \quad-\frac{G\left(u\left(z+h e^{j}\right)\right)-G(u(z))-G^{\prime}(u(z))\left(u\left(z+h e^{j}\right)-u(z)\right)}{h} .
\end{aligned}
$$

The assumptions $\nabla \psi_{j} \in L_{\infty}(D)$ and $(n, m) \in \mathcal{M}^{\prime}$ then further yield

$$
\begin{aligned}
L_{h}= & \frac{f\left(z+h e^{j}\right)-f(z)-h \partial_{z_{j}} f(z)}{h}+\nabla \psi_{j} \cdot \nabla u\left(z+h e^{j}\right)+\psi_{j} \Delta u\left(z+h e^{j}\right) \\
& -\frac{G\left(u\left(z+h e^{j}\right)\right)-G(u(z))-G^{\prime}(u(z))\left(u\left(z+h e^{j}\right)-u(z)\right)}{h} \in L_{2}(D) .
\end{aligned}
$$

As above this further implies $w_{h} \in W(D)$. Similarly, we know from Theorem 3.3 that $w_{h}(z)$ converges in $H_{0}^{1}(D)$ towards $w_{0}(z)=\partial_{z_{j}} u(z)$, which is the solution of

$$
-\nabla(a(\cdot, z) \nabla w)+G^{\prime}(u(z)) w=\partial_{z_{j}} f(z)+\nabla \psi_{j} \cdot \nabla u(z)+\psi_{j} \Delta u(z) .
$$

Again we conclude $w_{0} \in W(D)$, and we find by subtracting these two problems

$$
\begin{aligned}
-a(\cdot, z) \Delta\left(w_{h}-w_{0}\right)=- & G^{\prime}(u(z))\left(w_{h}-w_{0}\right)+\nabla \psi_{j} \cdot \nabla\left(u\left(z+h e^{j}\right)-u(z)\right) \\
& +\psi_{j} \Delta\left(u\left(z+h e^{j}\right)-u(z)\right)+\frac{f\left(z+h e^{j}\right)-f(z)-h \partial_{z_{j}} f(z)}{h} \\
& -\frac{G\left(u\left(z+h e^{j}\right)\right)-G(u(z))-G^{\prime}(u(z))\left(u\left(z+h e^{j}\right)-u(z)\right)}{h} .
\end{aligned}
$$

By dividing by $a(\cdot, z)$ and applying the $L_{2}(D)$-norm we then obtain

$$
\left|w_{h}-w_{0}\right|_{W} \longrightarrow 0 \quad(h \longrightarrow 0),
$$

since all the terms on the right hand side converge to 0 separately (this follows from Theorem 3.3, Proposition 3.2, Lemma 6.1, Assumption 4, and Proposition 6.1, in this order). This completes the proof.

### 6.2 Summability of Taylor, Legendre and Chebyshev coefficient sequences

We begin with the counterpart of Propositions 4.1, 5.1 and 5.2.
Lemma 6.2. Under the assumptions of Theorem 6.1, we have the estimate

$$
\left\|t_{\nu} \mid W(D)\right\| \leq M(\delta, B) \rho^{-\nu}
$$

for every $\nu \in \mathcal{F}$ and every $(\delta, B)$-admissible sequence $\rho$. If additionally we have $\rho_{j} \geq 1$ for all $j \in \mathbb{N}$ and $\rho_{j}>1$ for $j \in \operatorname{supp} \nu$, we also find

$$
\left\|u_{\nu} \mid W(D)\right\| \leq M(\delta, B) \prod_{j \in \operatorname{supp} \nu} \frac{\pi \rho_{j}}{2\left(\rho_{j}-1\right)}\left(2 \nu_{j}+1\right) \rho_{j}^{-\nu_{j}} \quad \text { and } \quad\left\|w_{\nu} \mid W(D)\right\| \leq M(\delta, B) \prod_{j \in \operatorname{supp} \nu} 2 \rho_{j}^{-\nu_{j}} .
$$

Proof. The proof of the first part is identical to the one of Proposition 4.1, using now the regularity result from Theorem 6.1, and replacing the estimate (4.1) by its counterpart based upon (6.4). In a similar way, we obtain the second part by the same argument as in the proofs of Propositions 5.1 and 5.2.

Theorem 6.2. Let $a$ satisfy Assumption 1 for some $0<r \leq R<\infty$. Suppose $f$ fulfills

$$
\begin{equation*}
\left\|f(z) \mid H^{-1}(D)\right\|<\left(\frac{\delta^{m}}{2^{2 m} m \iota(m+1)^{m+1}}\right)^{\frac{1}{m-1}}, \quad z \in \mathcal{A}_{\delta / 2, B} \tag{6.5}
\end{equation*}
$$

and Assumption 4. Moreover, let $\rho=\left(\rho_{j}\right)_{j \in \mathbb{N}}$ be an arbitrary $(\delta, B)$-admissible sequence with $B$ sufficiently large and $\rho_{j} \geq 1, j \in \mathbb{N}$, and assume

$$
\begin{equation*}
\left(\left\|\psi_{j} \mid W_{\infty}^{1}(D)\right\|\right)_{j \in \mathbb{N}} \in \ell_{1}(\mathbb{N}) \quad \text { and } \quad\left(\rho_{j}\left\|\psi_{j} \mid W_{\infty}^{1}(D)\right\|\right)_{j \in \mathbb{N}} \in \ell_{p}(\mathbb{N}) \quad \text { for some } p<1 \tag{6.6}
\end{equation*}
$$

Then the Taylor coefficients $t_{\nu}$ and the Legendre coefficients $u_{\nu}$ and $v_{\nu}$ of the solution $u$ of (2.2) satisfy $\left(\rho^{\nu}\left\|t_{\nu} \mid W(D)\right\|\right)_{\nu \in \mathcal{F}} \in \ell_{p}(\mathcal{F}),\left(\rho^{\nu}\left\|u_{\nu} \mid W(D)\right\|\right)_{\nu \in \mathcal{F}} \in \ell_{p}(\mathcal{F})$ and $\left(\rho^{\nu}\left\|v_{\nu} \mid W(D)\right\|\right)_{\nu \in \mathcal{F}} \in \ell_{p}(\mathcal{F})$. If additionally

$$
\begin{equation*}
\sup _{z \in \mathcal{U}_{\rho}}\left\|f(z)-f\left(z_{1}, \ldots, z_{J}, 0, \ldots,\right) \mid L_{2}(D)\right\| \longrightarrow 0 \quad(J \longrightarrow \infty) \tag{6.7}
\end{equation*}
$$

then it holds

$$
u(z)=\sum_{\nu \in \mathcal{F}} t_{\nu} z^{\nu}, \quad z \in \mathcal{U}_{\rho}, \quad \text { and } \quad u(z)=\sum_{\nu \in \mathcal{F}} u_{\nu} P_{\nu}(z)=\sum_{\nu \in \mathcal{F}} v_{\nu} L_{\nu}(z)=\sum_{\nu \in \mathcal{F}} w_{\nu} T_{\nu}(z), \quad z \in \mathcal{U}
$$

with uniform and unconditional convergence in $W(D)$, which is to be understood as in Theorems 4.1, 5.1 and 5.2.

Proof. The proof of the summability assertion is exactly the same as the one for Theorems 4.1, 5.1 and 5.2 upon replacing $\gamma_{j}=\left\|\psi_{j} \mid L_{\infty}(D)\right\|$ by $\widetilde{\gamma}_{j}=\left\|\psi_{j}\left|W_{\infty}^{1}(D)\|=\| \psi_{j}\right| L_{\infty}(D)\right\|+\left\|\nabla \psi_{j} \mid L_{\infty}(D)\right\|$, see also [2, Theorem 5.1]. Moreover, also the convergence proof remains the same upon now using the stability assertion for the $W$-semi-norm (Lemma 6.1).

## 7 Best $N$-term approximation rates

The convergence results from Theorems 4.1 and 5.1 may be re-interpreted in the following way: when approximating the mapping $y \mapsto u(y)$, i.e. approximately solving the parametric problem (2.6), instead of solving single instances of the problem (i.e. solving the equation for single fixed choices $\left.y^{(j)} \in U, j=1, \ldots, N\right)$ and afterward interpolating these points $u\left(y^{(j)}\right)$, we can try to simultaneously approximate the solution for all $y \in U$ by approximating $u$ by appropriate partial sums of the Taylor, Legendre or Chebyshev expansions. The main task consists in bounding truncation errors of best $N$-term approximation by monomials $y^{\nu}$ or polynomials $P_{\nu}, L_{\nu}$ or $T_{\nu}$, respectively. Moreover, we can utilize the summability results from Theorems 4.1, 5.1 and 5.2 to estimate the rate of convergence. To do so, we note that by Assumption 1 the constant sequence $\rho_{j}=1, j \in \mathbb{N}$, is $\delta$-admissible for all $0<\delta \leq r$. Then

$$
\sup _{y \in U}\left\|u(y)-S_{\Lambda}^{(T)} u(y)\left|H_{0}^{1}(D)\left\|=\sup _{y \in U}\right\| u(y)-\sum_{\nu \in \Lambda} t_{\nu} y^{\nu}\right| H_{0}^{1}(D)\right\| \leq \sum_{\nu \notin \Lambda}\left\|t_{\nu} \mid H_{0}^{1}(D)\right\|
$$

Similarly, we obtain by Parseval's identity

$$
\left\|u-S_{\Lambda}^{(L)} u\left|L_{2}\left(U ; d \mu ; H_{0}^{1}(D)\right)\left\|^{2}=\right\| u-\sum_{\nu \in \Lambda} v_{\nu} L_{\nu}(y)\right| L_{2}\left(U ; d \mu ; H_{0}^{1}(D)\right)\right\|^{2}=\sum_{\nu \notin \Lambda}\left\|v_{\nu} \mid H_{0}^{1}(D)\right\|^{2} .
$$

These observations are to be combined with Stechkin's Lemma, which in our case can be formulated as follows: Given a sequence $\gamma=\left(\gamma_{\nu}\right)_{\nu \in \mathcal{F}} \subset H_{0}^{1}(D)$, and denoting by $\Lambda_{N} \subset \mathcal{F}$ an index set corresponding to the $N$ largest values of $\left\|\gamma_{\nu} \mid H_{0}^{1}(D)\right\|$, it holds for $0<p \leq q \leq \infty$

$$
\left(\sum_{\nu \notin \Lambda_{N}}\left\|\gamma_{\nu} \mid H_{0}^{1}(D)\right\|^{q}\right)^{1 / q} \leq N^{-1 / p+1 / q}\left(\sum_{\nu \in \mathcal{F}}\left\|\gamma_{\nu} \mid H_{0}^{1}(D)\right\|^{p}\right)^{1 / p}
$$

and this choice of $\Lambda_{N}$ yields a best $N$-term approximation of $\gamma$. We have shown
Theorem 7.1. Let $a$ and $f$ satisfy the assumptions of Theorems 4.1, 5.1 and 5.2. If we denote by $\Lambda_{N}^{(T)} \subset \mathcal{F}$ an index set corresponding to the $N$ largest values of $\left\|t_{\nu} \mid H_{0}^{1}(D)\right\|$, then it holds

$$
\sup _{y \in U}\left\|u(y)-S_{\Lambda_{N}^{(T)}}^{(T)} u(y)\left|H_{0}^{1}(D)\|\leq\|\left(\left\|t_{\nu} \mid H_{0}^{1}(D)\right\|\right)_{\nu \in \mathcal{F}}\right| \ell_{p}(\mathcal{F})\right\| N^{-s}, \quad s=\frac{1}{p}-1
$$

This remains true for the Legendre Polynomials $P_{\nu}$ and the Chebyshev polynomials $T_{\nu}$ with the respective coefficients, partial sums and index sets for largest coefficients.
Moreover, if $\Lambda_{N}^{(L)} \subset \mathcal{F}$ is a set corresponding to the $N$ largest values of $\left\|v_{\nu} \mid H_{0}^{1}(D)\right\|$, then it holds

$$
\left\|u-S_{\Lambda_{N}^{(L)}}^{(L)} u\left|L_{2}\left(U, d \mu ; H_{0}^{1}(D)\right)\|\leq\|\left(\left\|v_{\nu} \mid H_{0}^{1}(D)\right\|\right)_{\nu \in \mathcal{F}}\right| \ell_{p}(\mathcal{F})\right\| N^{-s}, \quad s=\frac{1}{p}-\frac{1}{2}
$$

and once more this remains true for Chebyshev polynomials $T_{\nu}$ with respect to the measure $\eta$. Finally, if $a$ and $f$ satisfy the assumptions of Theorem 6.2, then all the estimates remain true upon replacing the $H_{0}^{1}(D)$-norm by the $W(D)$-norm.

If it comes to an actual computation of these approximating partial sums, one has to take into consideration that also the coefficients itself must be approximated, e.g. for Taylor coefficients by approximately solving linear problems as in Theorem 3.3.

Here the increased spatial regularity from Section 6 comes into play. These results enable us to solve the corresponding linear problems for the Taylor coefficients using Finite element techniques. If we denote by $t$ the rate of convergence for the Finite element method (depending on the method itself as well as the shape of the domain $D$ ), then the error of computing $t_{\nu} \in W(D)$, measured in the $H_{0}^{1}(D)$-norm, can be estimated by $C M_{\nu}^{-t}$, where $M_{\nu}$ is the number of degrees of freedom used in the computation of $t_{\nu}$. The total number of degrees of freedom required for approximating the partial $\operatorname{sum} \widetilde{S}_{\Lambda} u$ is then given by

$$
N_{\mathrm{dof}}=\sum_{\nu \in \Lambda} M_{\nu} .
$$

With this notation the approximation error in both, the parametric space as well as in the physical domain $D$, satisfies the estimate

$$
\sup _{y \in U}\left\|u(y)-\widetilde{S}_{\Lambda_{N}^{(T)}}^{(T)} u(y) \mid H_{0}^{1}(D)\right\| \leq C N_{\mathrm{dof}}^{-\min \{s, t\}}
$$

where the occurring constants can be estimated in dependence of $\left\|\left(\left\|t_{\nu} \mid H_{0}^{1}(D)\right\|\right)_{\nu \in \mathcal{F}} \mid \ell_{p}(\mathcal{F})\right\|$ and $\left\|\left(\left\|t_{\nu} \mid W(D)\right\|\right)_{\nu \in \mathcal{F}} \mid \ell_{p}(\mathcal{F})\right\|$. Similar considerations are true for Legendre or Chebyshev partial sums (see [2], Section 5).

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