# Analytic regularity and polynomial approximation of stochastic, parametric elliptic multiscale PDEs* 

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# Analytic regularity and polynomial approximation of stochastic, parametric elliptic multiscale PDEs * 

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#### Abstract

A class of second order, elliptic PDEs in divergence form with stochastic and anisotropic conductivity coefficients and $n$ known, separated microscopic length scales $\varepsilon_{i}, i=1, \ldots, n$ in a bounded domain $D \subset \mathbb{R}^{d}$ is considered. Neither stationarity nor ergodicity of these coefficients is assumed. Sufficient conditions are given for the random solution to converge $\mathbb{P}$-a.s, as $\varepsilon_{i} \rightarrow 0$, to a stochastic, elliptic one-scale limit problem in a tensorized domain of dimension $(n+1) d$. It is shown that this stochastic limit problem admits best $N$-term "polynomial chaos" type approximations which converge at a rate $\sigma>0$ that is determined by the summability of the random inputs' Karhúnen-Loève expansion. The convergence of the polynomial chaos expansion is shown to hold $\mathbb{P}$-a.s. and uniformly with respect to the scale parameters $\varepsilon_{i}$. Regularity results for the stochastic, one-scale limiting problem are established. An error bound for the approximation of the random solution at finite, positive values of the scale parameters $\varepsilon_{i}$ is established in the case of two scales, and in the case of $n>2$ scales convergence is shown, albeit without giving a convergence rate in this case.


## 1 Problem formulation

### 1.1 A class of stochastic multiscale elliptic problems

In a bounded Lipschitz domain $D \subset \mathbb{R}^{d}$ (to which we shall refer as "physical domain"), we consider diffusion problems in $D$ where the diffusion coefficients resp. the permeability is uncertain and exhibits microstructure on one or several microscopic length scales. In what follows, we assume these length scales to be separated and a priori known. To describe the periodic microstructure, let $Y$ denote the unit cube in $\mathbb{R}^{d}$ and let $Y_{1}, Y_{2}, \ldots, Y_{n}$ be $n$ copies of $Y$ which we assume to be the ranges of the $n$ fast- or microscopic variables (all our results generalize to the case when the $Y_{j}$ are nonidentical). To describe the random permeabilities that are admissible in our analysis, we assume given a probability space $(\Omega, \Sigma, \mathbb{P})$, and a random field

$$
\begin{equation*}
\Omega \ni \omega \mapsto A\left(\omega ; x, y_{1}, \ldots, y_{n}\right) \in L^{\infty}\left(D ; C_{\#}\left(Y_{1} \times \ldots \times Y_{n}\right)_{\mathrm{sym}}^{d \times d}\right) \tag{1.1}
\end{equation*}
$$

such that

$$
\begin{equation*}
A \in L^{\infty}\left(\Omega, d \mathbb{P} ; L^{\infty}\left(D ; C_{\#}\left(Y_{1} \times \ldots \times Y_{n}\right)_{s y m}^{d \times d}\right)\right) \tag{1.2}
\end{equation*}
$$

Throughout, for $0<p \leq \infty$ and a Banach space $B$, we denote by $L^{p}(\Omega, d \mathbb{P} ; B)$ the Bochner space of strongly $\mathbb{P}$ measurable mappings from $(\Omega, \Sigma)$ to $B$ with the sigma-algebra of Borel sets which are $p$ summable (resp. $\mathbb{P}$-a.s. bounded in $B$ in case that $p=\infty$ ). In (1.1) and the following, the notation \# indicates that the functions admit $Y_{i}$ periodic extensions to all of $\mathbb{R}^{d}$ with respect to each of the variables $y_{i}$ for $i=1, \ldots, n$ which locally, i.e. on compact subsets of $\mathbb{R}^{d}$, belong to the same function spaces on these sets. For notational conciseness, we denote by $\mathbf{Y}=Y_{1} \times \ldots \times Y_{n}$ and by $\boldsymbol{y}=\left(y_{1}, \ldots, y_{n}\right) \in \mathbf{Y}$. We will write $C_{\#}(\mathbf{Y})$ in place of $C_{\#}\left(Y_{1} \times \ldots \times Y_{n}\right)$. Spaces of vector functions with each component function belonging to a Banach space $B$ will be denoted by $B^{d}$, and of $d \times d$ matrix functions by $B^{d \times d}$. Integrals over such functions will be understood as vector functions of integrals over all component functions. To ensure well-posedness of our problem, we impose

[^1]Assumption 1.1 The diffusion matrix A satisfies (1.2). In particular, it is uniformly bounded and coercive, i.e. there are positive constants $\alpha$ and $\beta$ such that for all $\omega \in \Omega, x \in D$ and every $\boldsymbol{y} \in \mathbf{Y}$ it holds

$$
\forall \xi \in \mathbb{R}^{d}: \quad \alpha|\xi|^{2} \leq \xi^{\top} A(\omega ; x, \boldsymbol{y}) \xi \leq \beta|\xi|^{2}
$$

We assume $\mathbb{P}$-a.s. scale separation. This means that for a nondimensional scale parameter $\varepsilon>0$, there are $n$ known, deterministic, positive functions $1>\varepsilon_{1}(\varepsilon) \geq \ldots \geq \varepsilon_{n}(\varepsilon)>0$ which depend continuously and monotonically on $\varepsilon$, and which describe the $n$ microscopic length scales which the random diffusion coefficient depends on. Without loss of generality, we set $\varepsilon_{1}=\varepsilon$. If the random coefficient (1.4) has $n>1$ fast scales, we say that the coefficient is $\mathbb{P}$-a.s. scale separated if, for all $i=1, \ldots, n-1$ for $\mathbb{P}$-a.s., there holds

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \frac{\varepsilon_{i+1}}{\varepsilon_{i}}=0 \tag{1.3}
\end{equation*}
$$

If $n=1$, i.e. there is only one fast scale, we denote $\epsilon_{1}=\epsilon$ and the condition (1.3) is understood to be void. For $A$ as in (1.1) and satisfying Assumption 1.1, for a given family of scale parameters $\varepsilon_{i}$ satisfying (1.3), we define a family of $n$-scale, random multiscale diffusion tensors $A^{\varepsilon}(\omega ; x) \in L^{\infty}\left(\Omega, d \mathbb{P} ; L^{\infty}(D)_{s y m}^{d \times d}\right)$ by

$$
\begin{equation*}
A^{\varepsilon}(\omega ; x):=A\left(\omega ; x, \frac{x}{\varepsilon_{1}}, \ldots, \frac{x}{\varepsilon_{n}}\right) \tag{1.4}
\end{equation*}
$$

With $A^{\varepsilon}(\omega ; x)$ defined in this way, for given $f \in H^{-1}(D), 0<\varepsilon<1$ and for $\omega \in \Omega$, we consider in $D$ the $n$-scale stochastic Dirichlet problem:

$$
\begin{equation*}
-\operatorname{div} A^{\varepsilon}(\omega ; x) \nabla u^{\varepsilon}=f(x) \quad \text { in } D,\left.\quad u^{\varepsilon}\right|_{\partial D}=0 \tag{1.5}
\end{equation*}
$$

For simplicity of exposition, we assume in what follows that the source term $f \in H^{-1}(D)$ is deterministic and independent of $\varepsilon$. At this point we remark that stochastic homogenization problems have been considered before; we mention only $[3,5,12]$ and the references there. However, usually only two scales were considered and an ergodic hypothesis was imposed. In this work, neither stationarity nor ergodicity of the random coefficient will be assumed. We begin our analysis by casting problem (1.5) in variational form:

$$
\begin{equation*}
\text { find } u^{\varepsilon} \in H_{0}^{1}(D) \text { such that } \int_{D} A^{\varepsilon}(\omega ; x) \nabla u^{\varepsilon} \cdot \nabla \phi d x=\int_{D} f \phi d x \quad \forall \phi \in H_{0}^{1}(D) . \tag{1.6}
\end{equation*}
$$

We equip the space $H_{0}^{1}(D)$ with the norm $\|v\|_{H_{0}^{1}(D)}=\|\nabla v\|_{L^{2}(D)}$. Then the random solution $u^{\varepsilon}$ of (1.6) satisfies, for $\mathbb{P}-$ a.e. $\omega \in \Omega$,

$$
\left\|u^{\varepsilon}(\omega ; \cdot)\right\|_{H_{0}^{1}(D)} \leq \frac{\|f\|_{H^{-1}(D)}}{\alpha}
$$

We assume that the random coefficient $A$ in Assumption 1.1 and in (1.4) is characterized by a sequence $\boldsymbol{z}(\omega)=\left(z_{k}(\omega)\right)_{k \geq 1}$ of random variables as follows:

$$
\begin{equation*}
A(\omega ; x, \boldsymbol{y})=\bar{A}(x, \boldsymbol{y})+\sum_{k=1}^{\infty} z_{k}(\omega) \Psi_{k}(x, \boldsymbol{y}), \quad(\omega, x, \boldsymbol{y}) \in \Omega \times D \times \mathbf{Y} \tag{1.7}
\end{equation*}
$$

where $\Psi_{k}(x, \boldsymbol{y}) \in L^{\infty}(D, \mathbf{Y})_{\text {sym }}^{d \times d}$. Without any assumptions on the normalization of the $z_{k}, \Psi_{k}$, the parametric representation (1.7) is nonunique. We therefore assume that the $z_{k}(\omega)$ are i.i.d. and $z_{k} \sim$ $\mathcal{U}(-1,1)$. We further denote by $\boldsymbol{z}$ the coefficient vector $\left(z_{1}, z_{2}, \ldots\right) \in U:=[-1,1]^{\mathbb{N}}$ of realizations. For a sequence $\beta=\left(\beta_{k}\right)_{k \geq 1} \in \ell^{1}(\mathbb{N})$, we assume the matrix functions $\Psi_{k}$ in (1.7) to satisfy

$$
\begin{equation*}
\forall k \in \mathbb{N}: \forall \xi \in \mathbb{R}^{d}, x \in D, \boldsymbol{y} \in \mathbf{Y}: \quad\left|\xi^{\top} \Psi_{k}(x, \boldsymbol{y}) \xi\right| \leq \beta_{k}|\xi|^{2} \tag{1.8}
\end{equation*}
$$

which implies that the series (1.7) converges unconditionally, $\mathbb{P}-a . s$. We also assume that the mean field in (1.7), i.e. the matrix function $\bar{A} \in L^{\infty}(D \times \mathbf{Y})_{s y m}^{d \times d}$, satisfies

$$
\begin{equation*}
\forall \xi \in \mathbb{R}^{d}, x \in D, \boldsymbol{y} \in \mathbf{Y}: \quad \alpha_{0}|\xi|^{2} \leq \xi^{\top} \bar{A}(x, \boldsymbol{y}) \xi \leq \beta_{0}|\xi|^{2} \tag{1.9}
\end{equation*}
$$

To ensure that the random coefficient $A(\omega ; x, \boldsymbol{y})$ in (1.7) is well defined and coercive, we assume that in (1.7) the fluctuation expansion of $A-\bar{A}$ is dominated by the mean field $\bar{A}$ in the following sense:

Assumption 1.2 We assume (1.8), (1.9) and that in (1.7), the random variables $z_{k}$ are i.i.d. in $[-1,1]$. Moreover, with the ellipticity constant $\alpha_{0}$ in (1.9), we assume that the matrix functions $\bar{A}$ and $\Psi_{k}$ satisfy for some $\kappa>0$

$$
\sum_{k \geq 1} \beta_{k} \leq \frac{\kappa}{1+\kappa} \alpha_{0}
$$

Assumption 1.2 implies Assumption 1.1: we may choose

$$
\alpha=\alpha_{0}-\frac{\kappa}{1+\kappa} \alpha_{0}=\frac{1}{1+\kappa} \alpha_{0}, \quad \beta=\alpha_{0}+\frac{\kappa}{1+\kappa} \alpha_{0} .
$$

From (1.8) and Assumption 1.2, we have
Proposition 1.3 The following estimate holds

$$
\sum_{k \geq 1}\left\|\Psi_{k}\right\|_{L^{\infty}(D)^{d \times d}} \leq 2 \frac{\kappa}{1+\kappa} \alpha_{0}
$$

Proof From Assumption 1.2, for each $i=1, \ldots, d$ and for every $k=1,2, \ldots$, we have

$$
\left\|\left(\Psi_{k}\right)_{i i}\right\|_{L^{\infty}(D)} \leq \beta_{k}
$$

Fix two indices $i, j=1, \ldots, d$, and choose in (1.8) $\xi_{i}=1$ and $\xi_{j}=1$, and $\xi_{l}=0$ for $l \neq i, j$. Then

$$
\forall x \in D, \boldsymbol{y} \in \mathbf{Y}, k \in \mathbb{N}: \quad\left|\left(\Psi_{k}(x, \boldsymbol{y})\right)_{i i}+\left(\Psi_{k}(x, \boldsymbol{y})\right)_{j j}+\left(\Psi_{k}(x, \boldsymbol{y})\right)_{i j}+\left(\Psi_{k}(x, \boldsymbol{y})\right)_{j i}\right| \leq 2 \beta_{k}
$$

From this, we deduce

$$
\forall x \in D, \forall \boldsymbol{y} \in \mathbf{Y}, \forall k \in \mathbb{N}: \quad\left|\left(\Psi_{k}(x, \boldsymbol{y})\right)_{i j}\right|+\left|\left(\Psi_{k}(x, \boldsymbol{y})\right)_{j i}\right| \leq 4 \beta_{k}
$$

This implies the assertion.

### 1.2 Karhúnen-Loève expansion

We give a particular example of a parametric expansion (1.7), the Karhúnen-Loève expansion of a random matrix function $A(\omega ; x, \boldsymbol{y})$. We give, in particular, sufficient conditions in order for Assumption 1.2 to hold. We formulate these conditions in terms of the smoothness of the covariance of the matrix function $A(\omega ; x, \boldsymbol{y})$, which is given by the fourth order tensor

$$
\operatorname{Cov}[A]_{i j i^{\prime} j^{\prime}}\left(x, \boldsymbol{y}, x^{\prime}, \boldsymbol{y}^{\prime}\right)=\int_{\Omega}\left(A_{i j}(\omega ; x, \boldsymbol{y})-\bar{A}_{i j}(x, \boldsymbol{y})\right)\left(A_{i^{\prime} j^{\prime}}\left(\omega ; x^{\prime}, \boldsymbol{y}^{\prime}\right)-\bar{A}_{i^{\prime} j^{\prime}}\left(x^{\prime}, \boldsymbol{y}^{\prime}\right)\right) d \mathbb{P}(\omega)
$$

for $i, j, i^{\prime}, j^{\prime}=1, \ldots, d$. Then $\operatorname{Cov}[A]_{i j i^{\prime} j^{\prime}} \in L^{\infty}((D \times \mathbf{Y}) \times(D \times \mathbf{Y}), \mathbb{R})$, for all $i, j, i^{\prime}, j^{\prime}$ is the kernel of the (compact and self-adjoint) covariance operator $\mathbf{Q}_{A}: L^{2}(D \times \mathbf{Y})_{\mathrm{sym}}^{d \times d} \rightarrow L^{2}(D \times \mathbf{Y})_{\mathrm{sym}}^{d \times d}$ defined by

$$
\left(\mathbf{Q}_{A} \Phi\right)_{i j}(x, \boldsymbol{y}):=\int_{D} \int_{\mathbf{Y}} \operatorname{Cov}[A]_{i j i^{\prime} j^{\prime}}\left(x, \boldsymbol{y}, x^{\prime}, \boldsymbol{y}^{\prime}\right) \Phi_{i^{\prime} j^{\prime}}\left(x^{\prime}, \boldsymbol{y}^{\prime}\right) d \boldsymbol{y}^{\prime} d x^{\prime}
$$

Let $\lambda_{1} \geq \lambda_{2} \geq \ldots \geq 0$ denote the eigenvalues of $\mathbf{Q}_{A}$ and let $\Phi_{k} \in\left(L^{2}(D \times \mathbf{Y})\right)^{d \times d}$ denote the corresponding eigenvectors. We assume that $\left\|\Phi_{k}\right\|_{L^{2}(D \times \mathbf{Y})^{d \times d}}=1$ for all $k$. Any random field $A \in$ $L^{2}\left(\Omega ; L^{2}(D \times \mathbf{Y})_{\text {sym }}^{d \times d}\right)$ can be represented by a Karhúnen-Loève (KL) expansion

$$
\begin{equation*}
A(\omega ; x, \boldsymbol{y})=\bar{A}(x, \boldsymbol{y})+\sum_{k=1}^{\infty} \sqrt{\lambda_{k}} \Phi_{k}(x, \boldsymbol{y}) Z_{k}(\omega) \tag{1.10}
\end{equation*}
$$

where $Z_{k}$ are pairwise uncorrelated random variables that satisfy

$$
Z_{k}(\omega)=\frac{1}{\sqrt{\lambda_{k}}} \int_{D \times \mathbf{Y}}\left(A_{i j}-\bar{A}_{i j}\right)\left(\Phi_{k}\right)_{i j} d \boldsymbol{y} d x
$$

By Assumption 1.2, the random coefficients $Z_{k}$ in (1.10) are uniformly bounded, $\mathbb{P}-$ a.s. for all $k$. Note also that, due to the normalization Assumption $\left\|\Phi_{k}\right\|_{L^{2}(D \times \mathbf{Y})^{d \times d}}=1$ the probability densities of the random variables $Z_{k}$ are not necessarily supported in $[-1,1]$. To estimate the eigenvalues $\lambda_{k}$, we will use the following classical result (see, e.g., [16] and the references therein).

Lemma 1.4 Let $(H,\langle\cdot, \cdot\rangle)$ be a Hilbert space and let $\mathcal{C}$ be a symmetric, nonegative and compact linear operator from $H$ to $H$ whose eigenpairs are $\left(\lambda_{m}, \phi_{m}\right)_{m \geq 1}$. If $m \in \mathbb{N}$ and $\mathcal{C}_{m}$ is an operator of rank at most $m$, then

$$
\lambda_{m+1} \leq\left\|\mathcal{C}-\mathcal{C}_{m}\right\|_{\mathcal{L}(H, H)}
$$

We then have the following bounds for the eigenvalues $\lambda_{k}$ in terms of the regularity of the covariance function of the random diffusion matrix $A$ in (1.1).
Proposition 1.5 Assume that the random diffusion matrix $A$ in (1.1) satisfies

$$
\begin{equation*}
A \in L^{2}\left(\Omega ; H_{m i x}^{t}(D \times \mathbf{Y})\right) \tag{1.11}
\end{equation*}
$$

where, for $t \geq 0$, the space $H_{m i x}^{t}(D \times \mathbf{Y})$ is defined by $H_{m i x}^{t}(D \times \mathbf{Y})=H^{t}(D) \otimes H_{\#}^{t}\left(Y_{1}\right) \otimes \ldots \otimes H_{\#}^{t}\left(Y_{n}\right)$ with $\otimes$ denoting the tensor product of separable Hilbert spaces and with $H^{t}$ denoting, for noninteger values of $t$, the fractional order Sobolev space (e.g [18]).

Then $\operatorname{Cov}[A] \in H_{m i x}^{t}(D \times \mathbf{Y}) \otimes H_{m i x}^{t}(D \times \mathbf{Y})$. Moreover, for all $\varepsilon>0$, there exists a constant $c=c(\varepsilon)>0$ such that for all $k \geq 1$ holds $\lambda_{k} \leq c(\varepsilon) k^{-t / d+\varepsilon}$.
Proof The proof is adapted from that of Proposition 2.18 in [16]. Let $\mathbf{P}^{L}$ be the orthogonal projection into $\widehat{\mathbf{V}}^{L}$ in the $L^{2}(D \times \mathbf{Y})$ norm. Here $\widehat{\mathbf{V}}^{L}$ denotes a sparse tensor product space in the sense of $[15,10]$ of multilevel spaces in $D$ and in $Y_{i}, i=1, \ldots, n$. The rank of the operator $\mathbf{P}^{L} \mathbf{Q}_{A}$ is at most $\operatorname{dim} \hat{\mathbf{V}}^{L}$. Using results on sparse grid interpolation (see, e.g. [17] and the references there) we find that

$$
\left\|\mathbf{Q}_{A}-\mathbf{P}^{L} \mathbf{Q}_{A}\right\|_{\mathcal{L}\left(L^{2}(D \times \mathbf{Y}) \otimes L^{2}(D \times \mathbf{Y})\right)} \leq c L^{n / 2} 2^{-L t}
$$

As $\operatorname{dim} \widehat{\mathbf{V}}^{L}=O\left(L^{n} 2^{d L}\right)$, we get with Lemma 1.4 the conclusion by choosing $\mathcal{C}=\mathbf{Q}_{A}$ and $\mathcal{C}_{m}=\mathbf{P}^{L} \mathbf{Q}_{A}$ by choosing in Lemma $1.4 m=k=O\left(L^{n} 2^{d L}\right)$.

For the eigenfunctions $\Phi_{k}$, we have
Proposition 1.6 Assume that the random coefficient $A$ in (1.1) satisfies (1.11) for some $t>d / 2$. Then for every $d / 2<t^{*}<t$ there is a constant $c>0$ independent of $k$ such that

$$
\sum_{i, j=1}^{d}\left\|\left(\Phi_{k}\right)_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}^{2} \leq c \lambda_{k}^{-2 t^{*} / t}
$$

Proof The proof of this proposition follows that for Proposition 2.3 of Bieri et al. [4]. We note that

$$
\frac{\partial}{\partial x_{\alpha}}\left(\Phi_{k}\right)_{i j}(x, \boldsymbol{y})=\frac{1}{\lambda_{k}} \int_{D \times \mathbf{Y}} \frac{\partial}{\partial x_{\alpha}} \operatorname{Cov}[A]_{i j i^{\prime} j^{\prime}}\left(x, \boldsymbol{y}, x^{\prime}, \boldsymbol{y}^{\prime}\right) \Phi_{i^{\prime} j^{\prime}}\left(x^{\prime}, \boldsymbol{y}^{\prime}\right) d x^{\prime} d \boldsymbol{y}^{\prime} \quad i, j=1, \ldots, d
$$

with summation over repeated indices. Therefore, with the normalization $\left\|\Phi_{k}\right\|_{L^{2}(D \times \mathbf{Y})^{d \times d}}=1$,

$$
\forall k \in \mathbb{N}: \quad\left\|\left(\Phi_{k}\right)\right\|_{H_{m i x}^{t}(D \times \mathbf{Y})^{d \times d}} \leq \frac{1}{\lambda_{k}}\|\operatorname{Cov}[A]\|_{H_{m i x}^{t}(D \times \mathbf{Y}) \otimes L^{2}(D \times \mathbf{Y})}
$$

For $0<t^{*}<t$ hold the inclusions $L^{2}(D \times \mathbf{Y})=L^{2}(D) \otimes L^{2}\left(Y_{1}\right) \otimes \ldots \otimes L^{2}\left(Y_{n}\right) \supset H_{m i x}^{t^{*}}(D \times \mathbf{Y}) \supset$ $H_{m i x}^{t}(D \times \mathbf{Y})$ which follow by interpolation between $L^{2}$ and $H^{t}$ on $D$ respectively on $Y_{i}$ and by the fact that the Sobolev norms of mixed highest derivative are cross norms on the tensor products of the respective Hilbert spaces. It follows from the corresponding interpolation inequality (see, e.g., [18, Chap. 1]) that there exists a constant $C\left(t^{*}\right)>0$ such that

$$
\left\|\left(\Phi_{k}\right)_{i j}\right\|_{H_{m i x}^{t^{*}}(D \times \mathbf{Y})} \leq C\left\|\left(\Phi_{k}\right)_{i j}\right\|_{H_{m i x}^{t}(D \times \mathbf{Y})}^{t^{*} / t}\left\|\left(\Phi_{k}\right)_{i j}\right\|_{L^{2}(D \times \mathbf{Y})}^{1-t^{*} / t}
$$

Applying Hölder's inequality we then get for all $k$

$$
\begin{aligned}
\sum_{i, j=1}^{d}\left\|\left(\Phi_{k}\right)_{i j}\right\|_{H_{m i x}^{t *}(D \times \mathbf{Y})}^{2} & \leq\left(C \sum_{i, j=1}^{d}\left\|\left(\Phi_{k}\right)_{i j}\right\|_{H_{m i x}^{t}(D \times \mathbf{Y})}^{2}\right)^{t^{*} / t}\left(\sum_{i, j=1}^{d}\left\|\left(\Phi_{k}\right)_{i j}\right\|_{L^{2}(D \times \mathbf{Y})}^{2}\right)^{1-t^{*} / t} \\
& \leq C\left(\frac{c}{\lambda_{k}^{2}}\right)^{t^{*} / t}
\end{aligned}
$$

As $t^{*}>d / 2$, we deduce that exists $c>0$ such that

$$
\forall k \in \mathbb{N}: \quad \sum_{i, j=1}^{d}\left\|\left(\Phi_{k}\right)_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}^{2} \leq c \lambda_{k}^{-2 t^{*} / t}
$$

The conclusion then follows.
In the Karhúnen-Loève expansion (1.10), let $\Psi_{k}=\sqrt{\lambda_{k}} \Phi_{k}$. We then find that there exists a constant $c>0$ (depending on $t, t^{*}$ and on $d$ ) such that for all $k$

$$
\sum_{i, j=1}^{d}\left\|\left(\Psi_{k}\right)_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}^{2} \leq c \lambda_{k}^{1-2 t^{*} / t}
$$

From Proposition 1.5, we find that

$$
\sum_{i, j=1}^{d}\left\|\left(\Psi_{k}\right)_{i j}\right\|_{L^{\infty}(D \times Y)}^{2} \leq c k^{(-t / d+\varepsilon)\left(1-2 t^{*} / t\right)} .
$$

For each vector $\xi \in \mathbb{R}^{d}$, we have

$$
\left|\left(\Psi_{k}\right)_{i j}(x, \boldsymbol{y}) \xi_{i} \xi_{j}\right|^{2} \leq\left(\sum_{i, j=1}^{d}\left\|\left(\Psi_{k}\right)_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}^{2}\right)\left(\sum_{i, j=1}^{d} \xi_{i}^{2} \xi_{j}^{2}\right) \leq c k^{(-t / d+\varepsilon)\left(1-2 t^{*} / t\right)}|\xi|^{4}
$$

Therefore we may choose

$$
\begin{equation*}
\beta_{k}=c k^{(-t / d+\varepsilon)\left(1 / 2-t^{*} / t\right)} . \tag{1.12}
\end{equation*}
$$

When $t$ is sufficiently large, e.g. $(t / d-\varepsilon)\left(1 / 2-t^{*} / t\right)>1$, this implies that $\beta=\left\{\beta_{k}\right\}_{k \geq 1} \in \ell^{1}(\mathbb{N})$. Assuming that the random variables $Z_{k}$ in the expansion (1.10) are uniformly bounded, we can and will in what follows assume that they are rescaled so that the support of their laws equals $[-1,1]$. Assumption 1.2 holds when the constant $\alpha_{0}$ is sufficiently large.

### 1.3 Probability space

A key tool in our analysis will be a parametric deterministic representation of the law of the random multiscale solution $u^{\varepsilon}$. We shall use this representation in order to prove various convergence results of $u^{\varepsilon}$ as $\varepsilon \rightarrow 0$. Below, we shall investigate the precise regularity of dependence of this representation of $u^{\varepsilon}$ on the parameter vector $\boldsymbol{z}$. This, in turn, also allows for the proof of sharp bounds on spectral approximations of the parametric solution $u^{\varepsilon}$. To this end, following [7, 8], we parametrize the law of $u^{\varepsilon}(\omega ; x)$ in terms of countably many "random coordinates" $z_{k}(\omega)$ in the representation (1.7). We collect the random coordinates $\left(z_{k}\right)_{k \geq 1}$ in a vector $\boldsymbol{z}$ and define the parametric, deterministic multiscale coefficient $A^{\varepsilon}(\boldsymbol{z} ; \boldsymbol{x})$ as follows:

$$
\begin{equation*}
A^{\varepsilon}(\boldsymbol{z}, x):=A\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon_{1}}, \ldots, \frac{x}{\varepsilon_{n}}\right) . \tag{1.13}
\end{equation*}
$$

We define a probability measure on the parameter space $U=[-1,1]^{\mathbb{N}}$. To this end we introduce the $\sigma$-algebra $\Theta=\left(\mathcal{B}^{1}([-1,1])\right)^{\mathbb{N}}$ where $\mathcal{B}^{1}([-1,1])$ denotes the Borel $\sigma$-algebra on the interval $[-1,1]$. On the measurable space $(U, \Theta)$ thus obtained, we define a probability measure by

$$
\begin{equation*}
\rho(d \boldsymbol{z}):=\bigotimes_{j \geq 1} \frac{d z_{j}}{2} \tag{1.14}
\end{equation*}
$$

For any set of the form $S=\prod_{j=1}^{\infty} S_{j}$ with $S_{j} \in \mathcal{B}^{1}([-1,1])$, it holds $S \in \Theta$ and

$$
\rho(S)=\prod_{j=1}^{\infty} \mathbb{P}\left\{\omega: z_{j}(\omega) \in S_{j}\right\}
$$

### 1.4 Parametric deterministic multiscale problem

For each $\boldsymbol{z} \in U$, we define the deterministic coefficient matrix $A(\boldsymbol{z} ; x, \boldsymbol{y})$ by

$$
\begin{equation*}
A(\boldsymbol{z} ; x, \boldsymbol{y})=\bar{A}(x, \boldsymbol{y})+\sum_{k \geq 1} z_{k} \Psi_{k}(x, \boldsymbol{y}), \tag{1.15}
\end{equation*}
$$

where the matrix functions $\bar{A}$ and $\Psi_{k}$ are those in (1.7). The convergence of the sum on the right hand side is ensured by Proposition 1.3. For the parametric, deterministic coefficient $A^{\varepsilon}(\boldsymbol{z}, x)$ defined in (1.13), and for given $0<\varepsilon<1, \boldsymbol{z} \in U$ and $f \in H^{-1}(D)$, we consider the deterministic multiscale problem: for given $f \in H^{-1}(D)$ and $\boldsymbol{z} \in U$, find $u^{\varepsilon}(\boldsymbol{z}, \cdot) \in H_{0}^{1}(D)$ which satisfies

$$
\begin{equation*}
-\operatorname{div} A^{\varepsilon}(\boldsymbol{z} ; \cdot) \nabla u^{\varepsilon}(\boldsymbol{z} ; \cdot)=f(x),\left.\quad u^{\varepsilon}(\boldsymbol{z} ; \cdot)\right|_{\partial D}=0 . \tag{1.16}
\end{equation*}
$$

Again, Assumption 1.1 holds with $\alpha=\alpha_{0} /(1+\kappa)$ and $\beta=\alpha_{0}+\kappa \alpha_{0} /(1+\kappa)$, so problem (1.16) admits a unique solution which satisfies

$$
\begin{equation*}
\sup _{0<\varepsilon<1} \sup _{\boldsymbol{z} \in U}\left\|u^{\varepsilon}(\boldsymbol{z}, \cdot)\right\|_{H_{0}^{1}(D)} \leq \frac{\|f\|_{H^{-1}(D)}}{\alpha} . \tag{1.17}
\end{equation*}
$$

We first prove that the solution $u^{\varepsilon}(\boldsymbol{z}, \cdot)$ depends on $\boldsymbol{z}$ continuously.
Proposition 1.7 Under Assumption 1.1, there exists a constant $c>0$ which is independent of $\varepsilon$ such that

$$
\forall \boldsymbol{z}, \boldsymbol{z}^{\prime} \in U: \quad\left\|u^{\varepsilon}(\boldsymbol{z} ; \cdot)-u^{\varepsilon}\left(\boldsymbol{z}^{\prime} ; \cdot\right)\right\|_{H_{0}^{1}(D)} \leq c\left\|A(\boldsymbol{z} ; \cdot, \cdot)-A\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right\|_{L^{\infty}(D, C(\overline{\mathbf{Y}}))} .
$$

Proof Define $w^{\varepsilon}:=u^{\varepsilon}(\boldsymbol{z} ; \cdot)-u^{\varepsilon}\left(\boldsymbol{z}^{\prime} ; \cdot\right)$. The function $w^{\varepsilon}$ is a weak solution of the Dirichlet problem

$$
-\operatorname{div} A^{\varepsilon}(\boldsymbol{z} ; \cdot) \nabla w^{\varepsilon}=-\operatorname{div}\left[A^{\varepsilon}\left(\boldsymbol{z}^{\prime} ; \cdot\right)-A^{\varepsilon}(\boldsymbol{z} ; \cdot)\right] \nabla u^{\varepsilon}\left(\boldsymbol{z}^{\prime} ; \cdot\right),\left.\quad w^{\varepsilon}\right|_{\partial D}=0 .
$$

Therefore, it holds for every $\boldsymbol{z}, \boldsymbol{z}^{\prime} \in U$

$$
\forall w \in H_{0}^{1}(D): \quad \int_{D} A^{\varepsilon}(\boldsymbol{z}, x) \nabla w^{\varepsilon}(x) \cdot \nabla w(x) d x=\int_{D}\left[A^{\varepsilon}\left(\boldsymbol{z}^{\prime} ; x\right)-A^{\varepsilon}(\boldsymbol{z} ; x)\right] \nabla u^{\varepsilon}\left(\boldsymbol{z}^{\prime} ; \cdot\right) \cdot \nabla w(x) d x .
$$

From (1.17) and Assumption 1.1, we obtain the conclusion.
To study the law of the solution $u^{\varepsilon}$ of (1.5), we need to prove its measurability.
Proposition 1.8 For every $0<\varepsilon<1$, the solution $U \ni \boldsymbol{z} \mapsto u^{\varepsilon}(\boldsymbol{z} ; \cdot)$ of (1.16) is measurable as a map from $U$ to $H_{0}^{1}(D)$.

Proof As $H_{0}^{1}(D)$ is separable, it is sufficient to show that $u$ is weakly measurable, i.e. for all $\phi \in H_{0}^{1}(D)$, the $H_{0}^{1}(D)$ innerproduct $\left\langle u^{\varepsilon}(\boldsymbol{z}, \cdot), \phi\right\rangle$ is measurable as a map from $U$ to $\mathbb{R}$. For $a \in \mathbb{R}$ we denote by $Y_{a}=\left\{\boldsymbol{z} \in U:\left\langle u^{\varepsilon}(\boldsymbol{z}, \cdot), \phi\right\rangle>a\right\}$. From Proposition 1.7 it follows that if $\left\langle u^{\varepsilon}(\boldsymbol{z}), \phi\right\rangle>a$, then there exists a positive constant $r$ such that if

$$
\sup _{x, \boldsymbol{y}}\left|A_{i j}(\boldsymbol{z} ; x, \boldsymbol{y})-A_{i j}\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right|<r \quad \text { for all } \quad i, j=1, \ldots, d
$$

then $\left\langle u^{\varepsilon}\left(\boldsymbol{z}^{\prime} ; \cdot\right), \phi\right\rangle>a$. Let $T_{k}$ be the set of $\boldsymbol{z} \in U$ such that for all $\overline{\boldsymbol{z}}=\left(z_{1}, z_{2}, \ldots, z_{k}, \bar{z}_{1}, \bar{z}_{2}, \ldots\right)$ $\left\langle u^{\varepsilon}(\overline{\boldsymbol{z}}, \cdot), \phi\right\rangle>a$ for all $\bar{z}_{j} \in[-1,1], j=1,2, \ldots$ For each $\boldsymbol{z} \in U$, from Proposition 1.3, we deduce that for all $i, j=1, \ldots, d$,

$$
\sup _{x, \boldsymbol{y}}\left|A_{i j}(\boldsymbol{z} ; x, \boldsymbol{y})-A_{i j}(\overline{\boldsymbol{z}} ; x, \boldsymbol{y})\right|<r,
$$

for all $\bar{z}_{1}, \bar{z}_{2}, \ldots \in[-1,1]$ when $k$ is sufficiently large. Therefore each vector $\boldsymbol{z} \in Y_{a}$ belongs to a set $T_{k}$ for some constant $k$. Let $R_{k} \in[-1,1]^{k}$ be the set of $t=\left(t_{1}, \ldots, t_{k}\right)$ such that $\left(t_{1}, t_{2}, \ldots, t_{k}, \bar{z}_{1}, \bar{z}_{2}, \ldots\right) \in T_{k}$ for all $\bar{z}_{i} \in[-1,1]$. From Proposition 1.7, $R_{k}$ is an open set, and therefore is the union of a countable set of open rectangles. Thus $T_{k}$ is a countable union of sets in $\Theta$ and is therefore measurable, so is $Y_{a}=\cup_{k} T_{k}$.

Remark 1.9 The random solution $u^{\varepsilon}(\omega ; \cdot)$ of problem (1.5) can be recovered from the parametric, deterministic solution of (1.16), for $\mathbb{P}$-a.e. $\omega \in \Omega$ via

$$
\Omega \ni \omega \mapsto u^{\varepsilon}(\omega ; \cdot)=\left.u^{\varepsilon}(\boldsymbol{z} ; \cdot)\right|_{\boldsymbol{z}=\boldsymbol{z}(\omega)} \in H_{0}^{1}(D)
$$

where, for $0<\varepsilon<1$, $u^{\varepsilon}(\boldsymbol{z} ; x)$ denotes the weak solution of the parametric, deterministic problem (1.16).
We shall use Remark 1.9 in what follows to homogenize (1.5). We do this by first passing to the $(n+1)$ scale limit in the parametric, deterministic problem (1.16) and then "reinsert" $\boldsymbol{z}=\boldsymbol{z}(\omega)$.

### 1.5 One-scale stochastic limiting problem

For each realization $\omega \in \Omega$, we study the limit when $\varepsilon \rightarrow 0$ of the solution $u^{\varepsilon}$ of the problem (1.5). Multiscale convergence is an appropriate tool for this purpose. It was first introduced for two-scale problems by Nguetseng [14] and elaborated further by Allaire [1]. The definition of $n+1$-scale convergence we give below is due to Allaire and Briane [2]; we use their notion of multiscale convergence to study solutions of the problem (1.5) as $\varepsilon \rightarrow 0$.
Definition 1.10 A bounded sequence $\left\{u^{\varepsilon}\right\}_{\varepsilon}$ in $L^{2}(D) n+1$-scale converges to a function $u_{0} \in L^{2}(D \times \mathbf{Y})$ if for all test functions $\phi \in L^{2}\left(D, C_{\#}(\mathbf{Y})\right)$ it holds

$$
\lim _{\varepsilon \rightarrow 0} \int_{D} u^{\varepsilon}(x) \phi\left(x, \frac{x}{\varepsilon_{1}}, \ldots, \frac{x}{\varepsilon_{n}}\right) d x=\int_{D} \int_{\mathbf{Y}} u_{0}(x, \boldsymbol{y}) \phi(x, \boldsymbol{y}) d \boldsymbol{y} d x
$$

Here and throughout, we denote

$$
\int_{\mathbf{Y}} \cdot d \boldsymbol{y}=\int_{Y_{1}} \ldots \int_{Y_{n}} \cdot d y_{n} \ldots d y_{1}
$$

The use of the preceding definition in homogenization is due to the following theorem from [2].
Theorem 1.11 Any bounded sequence $\left\{u^{\varepsilon}\right\}$ in $L^{2}(D)$ contains an $n+1$-scale convergent subsequence.
For the variational formulation of the limiting problem of (1.16) using $n+1$-scale convergence, we introduce the space

$$
\mathbf{V}=\left\{\boldsymbol{v}=\left(v_{0},\left\{v_{i}\right\}\right): v_{0} \in H_{0}^{1}(D), v_{i} \in V_{i}\right\}
$$

where

$$
\left.V_{1}=L^{2}\left(D ; H_{\#}^{1}\left(Y_{1}\right) / \mathbb{R}\right), \quad V_{i}=L^{2}\left(D \times Y_{1} \times \ldots \times Y_{i-1} ; H_{\#}^{1}\left(Y_{i}\right) / \mathbb{R}\right)\right\}, \quad i=2,3, \ldots, n
$$

We equip $V_{i}$ with the norm

$$
\|\boldsymbol{v}\|\|=\| \nabla v_{0}\left\|_{L^{2}(D)}+\sum_{i=1}^{n}\right\| \nabla_{y_{i}} v_{i} \|_{L^{2}\left(D \times \mathbf{Y}_{i}\right)} \text { where } \quad \mathbf{Y}_{i}:=Y_{1} \times \ldots \times Y_{i}
$$

For each $\boldsymbol{v} \in \mathbf{V}$, we denote by

$$
\begin{equation*}
\nabla \boldsymbol{v}=\nabla_{x} v_{0}+\sum_{i=1}^{n} \nabla_{y_{i}} v_{i} \tag{1.18}
\end{equation*}
$$

Theorem 1.12 For every fixed $\boldsymbol{z} \in U$, as $\varepsilon \rightarrow 0$ the solution $u^{\varepsilon}(\boldsymbol{z} ; \cdot)$ of the parametric, deterministic multiscale problem (1.16) converges weakly in $H_{0}^{1}(D)$ to a function $u_{0}(\boldsymbol{z} ; \cdot) ;$ moreover, $\nabla u^{\varepsilon}(\boldsymbol{z}, \cdot) n+1$ scale converges to $\boldsymbol{\nabla} \boldsymbol{u}$ where $\boldsymbol{u}(\boldsymbol{z})=\left(u_{0}, u_{1}, \ldots, u_{n}\right) \in \mathbf{V}$ is the unique solution of the parametric, deterministic elliptic one-scale limiting problem

$$
\begin{equation*}
b(\boldsymbol{z} ; \boldsymbol{u}, \boldsymbol{v})=\int_{D} \int_{\mathbf{Y}} A(\boldsymbol{z} ; x, \boldsymbol{y}) \boldsymbol{\nabla} \boldsymbol{u} \cdot \boldsymbol{\nabla} \boldsymbol{v} d \boldsymbol{y} d x=\int_{D} f v_{0} d x, \quad \forall \boldsymbol{v}=\left(v_{0},\left\{v_{i}\right\}\right) \in \mathbf{V} \tag{1.19}
\end{equation*}
$$

Here, the parametric bilinear form $b(\boldsymbol{z} ; \boldsymbol{u}, \boldsymbol{v}): \mathbf{V} \times \mathbf{V} \rightarrow \mathbb{R}$ is bounded and coercive uniformly for $\boldsymbol{z} \in U$ : there exist positive constants $c_{1}$ and $c_{2}$ which are independent of $\boldsymbol{z} \in U$ such that

$$
\begin{equation*}
\forall \boldsymbol{z} \in U \quad \forall \boldsymbol{u} \in \mathbf{V}: \quad b(\boldsymbol{z} ; \boldsymbol{u}, \boldsymbol{u}) \geq c_{1}\| \| \boldsymbol{u}\| \|^{2} \tag{1.20}
\end{equation*}
$$

and

$$
\begin{equation*}
\forall \boldsymbol{u}, \boldsymbol{v} \in \mathbf{V}: \sup _{\boldsymbol{z} \in U}|b(\boldsymbol{z} ; \boldsymbol{u}, \boldsymbol{v})| \leq c_{2}\| \| \boldsymbol{u}\| \|\|\boldsymbol{v}\| \| \tag{1.21}
\end{equation*}
$$

For each fixed $\boldsymbol{z} \in U$, this theorem is a consequence of [2]; the coefficients $c_{1}$ and $c_{2}$ only depend on $\alpha$ and $\beta$ in Assumption 1.1 and are, therefore, independent of $\boldsymbol{z}$. From this, we obtain

$$
\sup _{\boldsymbol{z} \in U} \left\lvert\,\|u(\boldsymbol{z})\| \leq \frac{\|f\|_{H^{-1}(D)}}{c_{1}}\right.
$$

Using this a-priori bound, one verifies that the passage to the $n+1$-scale limit can be achieved uniformly with respect to $\boldsymbol{z} \in U$. Theorem 1.12 establishes convergence of the parametric solutions $u^{\varepsilon}(\boldsymbol{z} ; \cdot)$ as $\varepsilon \rightarrow 0$ to a solution to the high dimensional, parametric and deterministic one-scale problem for each fixed parameter vector $\boldsymbol{z} \in U$. To establish the connection between the solution of this problem and the laws of the random multiscale solutions $u^{\varepsilon}$ of (1.5), we next verify measurability of the solution $\boldsymbol{u}(\boldsymbol{z})$ with respect to $\rho(d \boldsymbol{z})$.
Proposition 1.13 The solution $\boldsymbol{u}(\boldsymbol{z})$ of (1.19) as a map from $(U, \Theta, \rho(d \boldsymbol{z}))$ to $\mathbf{V}$ is measurable.
Proof For any two vectors $\boldsymbol{z}, \boldsymbol{z}^{\prime} \in U$, let $\boldsymbol{u}(\boldsymbol{z})$ and $\boldsymbol{u}\left(\boldsymbol{z}^{\prime}\right)$ be the solutions of the problems (1.19). Define $\boldsymbol{w}=\boldsymbol{u}(\boldsymbol{z})-\boldsymbol{u}\left(\boldsymbol{z}^{\prime}\right)$. From (1.19), we find

$$
\forall \boldsymbol{v} \in \mathbf{V}: \quad \int_{D} \int_{\mathbf{Y}} A(\boldsymbol{z} ; x, \boldsymbol{y}) \boldsymbol{\nabla} \boldsymbol{w} \cdot \boldsymbol{\nabla} \boldsymbol{v} d \boldsymbol{y} d x=\int_{D} \int_{\mathbf{Y}}\left(A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})\right) \boldsymbol{\nabla} u\left(\boldsymbol{z}^{\prime}\right) \cdot \boldsymbol{\nabla} \boldsymbol{v} d \boldsymbol{y} d x
$$

We choose $\boldsymbol{v}=\boldsymbol{w}$. From (1.20), $\|\boldsymbol{u}(\boldsymbol{z})\| \|$ is bounded uniformly for all $\boldsymbol{z} \in U$. Therefore, there exists a constant $c$ which does not depend on $\boldsymbol{z}, \boldsymbol{z}^{\prime} \in U$ such that

$$
\begin{equation*}
\left\|\boldsymbol{u}(\boldsymbol{z})-\boldsymbol{u}\left(\boldsymbol{z}^{\prime}\right)\right\| \leq c\left\|A(\boldsymbol{z} ; x, \boldsymbol{y})-A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right\|_{L^{\infty}(D \times \mathbf{Y})} . \tag{1.22}
\end{equation*}
$$

The proof then follows the argument used in the proof of Proposition 1.8.
We define

$$
\begin{equation*}
\underline{\mathbf{V}}=L^{2}(U, \rho ; \mathbf{V})=\left\{\boldsymbol{v}=\left\{\left(v_{0},\left\{v_{i}\right\}\right): v_{0} \in L^{2}\left(U, \rho ; H_{0}^{1}(D)\right), v_{i} \in L^{2}\left(U, \rho ; V_{i}\right)\right\}\right\} \tag{1.23}
\end{equation*}
$$

We note in passing that $L^{2}(U, \rho ; \mathbf{V}) \cong L^{2}(U, \rho) \otimes \mathbf{V}$ and consider the variational parametric, deterministic problem:

$$
\begin{equation*}
\text { find } \boldsymbol{u} \in \underline{\mathbf{V}} \text { such that } B(\boldsymbol{u}(\boldsymbol{z} ; \cdot), \boldsymbol{v})=F(\boldsymbol{v}), \quad \forall \boldsymbol{v}=\left(v_{0},\left\{v_{i}\right\}\right) \in \underline{\mathbf{V}} \tag{1.24}
\end{equation*}
$$

Here, the linear functional $F: \underline{\mathbf{V}} \rightarrow \mathbb{R}$ and the variational form $B(\cdot, \cdot): \underline{\mathbf{V}} \times \underline{\mathbf{V}} \rightarrow \mathbb{R}$ are given by

$$
F(\boldsymbol{v})=\int_{U} \int_{D} f(x) v_{0}(\boldsymbol{z} ; x) d x d \rho(\boldsymbol{z}), \quad B(\boldsymbol{u}, \boldsymbol{v})=\int_{U} b(\boldsymbol{z} ; \boldsymbol{u}, \boldsymbol{v}) d \rho(\boldsymbol{z}) .
$$

Proposition 1.14 Problem (1.24) admits a unique parametric, deterministic solution $\boldsymbol{u}(\boldsymbol{z} ; \cdot, \cdot) \in \mathbf{V}$ which belongs to $L^{2}(U, \rho ; \mathbf{V})$. For $\rho$-a.e. $\boldsymbol{z} \in U$, this solution coincides with the solution $\boldsymbol{u}(\boldsymbol{z} ; \cdot, \cdot)$ of the parametric problem (1.19).
Proof The existence and uniqueness of a solution to (1.24) follow from Lax-Milgram theorem.
For each $\boldsymbol{z} \in U$, the solution $\boldsymbol{u}(\boldsymbol{z} ; \cdot, \cdot) \in \mathbf{V}$ of the parametric, deterministic elliptic one-scale problem (1.19) exists, is unique and is uniformly bounded with respect to $\boldsymbol{z} \in U$. As a mapping $U \ni \boldsymbol{z} \mapsto \boldsymbol{u}(\boldsymbol{z} ; \cdot) \in$ $\mathbf{V}$, it is measurable. As $d \rho(\boldsymbol{z})$ is a probability measure on $U$, this implies that the parametric solution $\boldsymbol{u}(\boldsymbol{z} ; \cdot)$ of (1.19) coincides with the solution $\boldsymbol{u} \in \underline{\mathbf{V}}$ of (1.24).

Remark 1.15 The random solution $u^{\varepsilon}(\omega ; \cdot)$ of problem (1.5) $(n+1)$-scale converges, for $\mathbb{P}$-a.e. $\omega \in \Omega$, towards the weak solution $\boldsymbol{u}(\omega ; x, \boldsymbol{y})$ of the random one-scale limiting problem

$$
\begin{equation*}
b(\omega ; \boldsymbol{u}(\omega ; \cdot) ; \boldsymbol{v})=\int_{D} \int_{\mathbf{Y}} A(\omega ; x, \boldsymbol{y}) \boldsymbol{\nabla} \boldsymbol{u} \cdot \nabla \boldsymbol{v} d \boldsymbol{y} d x=\int_{D} f v_{0} d x, \quad \forall \boldsymbol{v}=\left(v_{0},\left\{v_{i}\right\}\right) \in \mathbf{V} \tag{1.25}
\end{equation*}
$$

where, for $\omega \in \Omega$, we define

$$
\Omega \ni \omega \mapsto b(\omega ; \boldsymbol{v}, \boldsymbol{w}):=\left.b(\boldsymbol{z} ; \boldsymbol{v}, \boldsymbol{w})\right|_{\boldsymbol{z}=\boldsymbol{z}(\omega)}, \quad \boldsymbol{u}(\omega ; x, \boldsymbol{y})=\left.\boldsymbol{u}(\boldsymbol{z} ; x, \boldsymbol{y})\right|_{\boldsymbol{z}=\boldsymbol{z}(\omega)} .
$$

Our aim is to construct approximations of $u^{\varepsilon}(\omega ; x)$ which are, on the one hand, robust with respect to $\varepsilon$, and, on the other hand, allow for discretization of the randomness with convergence rates superior to that of Monte Carlo Methods. To this end, we shall investigate next a spectral, "polynomial chaos" type approximation of the solution $\boldsymbol{u}(\boldsymbol{z} ; x, \boldsymbol{y})$ of the $(n+1)$-scale limiting problem with respect to the parameter vector $\boldsymbol{z} \in U$, and then investigate the rate of convergence as $\varepsilon \rightarrow 0$ of $u^{\varepsilon}(\omega ; x)$ to the solution of the limiting problem.

## 2 Galerkin Approximations in $U$

### 2.1 Orthonormal basis of $L^{2}(U, \rho)$

We start by defining a "generalized polynomial chaos" basis of $L^{2}(U, \rho(d \boldsymbol{z}))$. Let $\left(L_{n}\right)_{n \geq 0}$ be the univariate Legendre polynomials normalized so that

$$
\begin{equation*}
\frac{1}{2} \int_{-1}^{1}\left|L_{n}(t)\right|^{2} d t=1 \tag{2.1}
\end{equation*}
$$

Let $\mathcal{F}$ be the (countable) set of all sequences $\nu=\left(\nu_{j}\right)_{j \geq 1}$ of nonnegative integers such that only a finite number of $\nu_{j}$ are non zero, i.e. $\mathcal{F}=\left\{\nu \in \mathbb{N}_{0}^{\mathbb{N}}:\|\nu\|_{1}<\infty\right\}$. For $\nu \in \mathcal{F}$, we define the tensorized Legendre polynomials as

$$
L_{\nu}(\boldsymbol{z})=\prod_{j \geq 1} L_{\nu_{j}}\left(z_{j}\right), \quad \nu \in \mathcal{F} .
$$

By the completeness of the Legendre polynomials $\left(L_{n}(t)\right)_{n \geq 0}$ in $L^{2}(-1,1)$, the family $L_{\nu}$ forms an orthonormal basis of $L^{2}(U, \rho)$ : each function $\boldsymbol{u} \in \underline{\mathbf{V}}$ can be expanded in the Legendre series

$$
\begin{equation*}
\boldsymbol{u}=\sum_{\nu \in \mathcal{F}} \boldsymbol{u}_{\nu} L_{\nu}, \quad \boldsymbol{u}_{\nu} \in \mathbf{V} \tag{2.2}
\end{equation*}
$$

### 2.2 Semidiscretization with respect to $z$

For a subset $\Lambda \subset \mathcal{F}$ of finite cardinality, we define the space

$$
\underline{\mathbf{V}}_{\Lambda}=\left\{\boldsymbol{u}_{\Lambda}=\sum_{\nu \in \Lambda} \boldsymbol{u}_{\nu}(x, \boldsymbol{y}) L_{\nu}(\boldsymbol{z}): \boldsymbol{u}_{\nu} \in \mathbf{V}\right\} \subset \underline{\mathbf{V}} .
$$

We then consider the following Galerkin semidiscretization in $z$ :

$$
\begin{equation*}
\text { Find } \quad \boldsymbol{u}_{\Lambda} \in \underline{\mathbf{V}}_{\Lambda}: \quad B\left(\boldsymbol{u}_{\Lambda}, \boldsymbol{v}_{\Lambda}\right)=F\left(\boldsymbol{v}_{\Lambda}\right) \quad \forall \boldsymbol{v}_{\Lambda} \in \underline{\mathbf{V}}_{\Lambda} \tag{2.3}
\end{equation*}
$$

Then the following approximation result holds.
Theorem 2.1 For all $\Lambda \subset \mathcal{F}$, problem (2.3) admits a unique solution $\boldsymbol{u}_{\Lambda} \in \underline{\mathbf{V}}_{\Lambda}$ which satisfies the following error estimate:

$$
\begin{equation*}
\left\|\boldsymbol{u}-\boldsymbol{u}_{\Lambda}\right\|_{\underline{\mathbf{V}}} \leq\left(\sum_{\nu \in \mathcal{F} \backslash \Lambda}\left\|\boldsymbol{u}_{\nu}\right\|_{\mathrm{V}}^{2}\right)^{1 / 2} \tag{2.4}
\end{equation*}
$$

Proof $A s \underline{\mathbf{V}}_{\Lambda}$ is a Hilbert space, from (1.20) and (1.21) and from the Lax-Milgram lemma, (2.3) admits a unique solution $\boldsymbol{u}_{\Lambda} \in \underline{\mathbf{V}}_{\Lambda}$. From Cea's lemma and from the normalization (2.1) with Parseval's equality, we find that

$$
\left\|\boldsymbol{u}-\boldsymbol{u}_{\Lambda}\right\|_{\underline{\mathbf{V}}} \leq \inf _{\boldsymbol{v}_{\Lambda} \in \underline{\mathbf{V}}_{\Lambda}}\left\|\boldsymbol{u}-\boldsymbol{v}_{\Lambda}\right\|_{\underline{\mathbf{v}}} .
$$

Choosing $\boldsymbol{v}_{\Lambda}=\sum_{\nu \in \Lambda} \boldsymbol{u}_{\nu} L_{\nu}$, we arrive at the conclusion.

## 3 Best $N$-term approximations

From Theorem 2.1, with a fixed cardinality $N$, we infer that an optimal choice of the set $\Lambda$ is to select $\Lambda$ corresponding to $N$ terms $\boldsymbol{u}_{\nu}$ with largest $\mathbf{V}$ norms. Since these norms are not known a priori, we establish in this section an apriori bound for them, and choose the set $\Lambda$ according to these bounds. In this way, we obtain a constructive approach for the choices of index sets $\Lambda$ with the prescribed cardinality which might be, however, suboptimal. Nevertheless, we shall prove that the sets obtained in this way will allow for the best $N$-term convergence rates to be achieved. The key ingredient for obtaining the rate of convergence in terms of the cardinality of $\Lambda$ is the following observation, due to Stechkin.

Lemma 3.1 Let $\alpha=\left(\alpha_{\nu}\right)_{\nu \in \mathcal{F}}$ be a sequence in $\ell^{p}(\mathcal{F})$. Let $q \geq p \geq 0$. If $\Lambda_{N} \subset \mathcal{F}$ is the set of indices corresponding to a set of $N$ largest $\left|\alpha_{\nu}\right|$, then for every $N$ holds

$$
\|\alpha\|_{\ell\left(\mathcal{F} \backslash \Lambda_{N}\right)}=\left(\sum_{\nu \in \mathcal{F} \backslash \Lambda_{N}}\left|\alpha_{\nu}\right|^{q}\right)^{1 / q} \leq\|\alpha\|_{\ell^{p}(\mathcal{F})} N^{-\sigma}, \text { where } \sigma=\frac{1}{p}-\frac{1}{q} \geq 0
$$

The convergence rate of truncated gpc expansions therefore depends on the $p$-summability of the sequence $\left(\left\|\boldsymbol{u}_{\nu}\right\|_{\mathbf{V}}\right)_{\nu \in \mathcal{F}}$. We show that the summability of this sequence depends on the summability of the sequence $\beta_{k}$ in (1.8).
Assumption 3.2 There exists $0<p<1$ such that in (1.8) the sequence $\left(\beta_{k}\right)_{k \geq 1} \in \ell^{p}(\mathbb{N})$.
Remark 3.3 Assumption 3.2 holds if the constants $t$ and $t^{*}$ in (1.12) satisfy

$$
p\left(\frac{t}{d}-\varepsilon\right)\left(\frac{1}{2}-\frac{t^{*}}{t}\right)>1
$$

### 3.1 Complex extension of the parametric deterministic problem

To bound $\left\|\boldsymbol{u}_{\nu}\right\|_{\mathbf{V}}$, we follow [8] and extend the parametric, deterministic limit problem (1.19) to complex values of the parameters $\boldsymbol{z}$. Let $M$ be a positive constant. Let $K<1$ be a positive constant such that

$$
K \sum_{j=1}^{\infty} \beta_{j}<\frac{\alpha}{2 M}
$$

We choose a constant $J_{0}$ such that

$$
\sum_{j>J_{0}} \beta_{j}<\frac{\alpha K}{6 M(1+K)}
$$

Let $E=\left\{1,2, \ldots, J_{0}\right\}$ and $F=\mathbb{N} \backslash E$. We define

$$
\left|\nu_{F}\right|=\sum_{j>J_{0}}\left|\nu_{j}\right|
$$

For each $\nu \in \mathcal{F}$, we define

$$
\begin{equation*}
r_{m}=K \text { when } m \leq J_{0}, \text { and } r_{m}=\max \left\{1, \frac{\alpha \nu_{m}}{M\left|\nu_{F}\right| \beta_{m}}\right\} \text { when } m>J_{0} \tag{3.1}
\end{equation*}
$$

where we make the convention that $\frac{\left|\nu_{j}\right|}{\left|\nu_{F}\right|}=0$ if $\left|\nu_{F}\right|=0$. For $m \geq 1$, we let the set $\mathcal{U}_{m} \in \mathbb{C}$ be defined as

$$
[-1,1] \subset \mathcal{U}_{m}:=\left\{\zeta_{m} \in \mathbb{C}: \operatorname{dist}\left(\zeta_{m},[-1,1]\right) \leq r_{m}\right\} \subset \mathbb{C}
$$

We next extend problem (1.19) to the complex parametric domain

$$
\begin{equation*}
\mathcal{U}=\bigotimes_{m=1}^{\infty} \mathcal{U}_{m} \subset \mathbb{C}^{\mathbb{N}} \tag{3.2}
\end{equation*}
$$

We define the complex parametric coefficient $A(\boldsymbol{\zeta} ; x, \boldsymbol{y})$

$$
A(\boldsymbol{\zeta} ; x, \boldsymbol{y}):=\bar{A}(x, \boldsymbol{y})+\sum_{m=1}^{\infty} \zeta_{m} \Psi_{m}(x, \boldsymbol{y}), \quad \boldsymbol{\zeta} \in \mathcal{U}, x \in D, \boldsymbol{y} \in \mathbf{Y}
$$

The sum on the right hand side of this definition converges uniformly for $\boldsymbol{\zeta} \in \mathcal{U}, x \in D$ and for $\boldsymbol{y} \in \mathbf{Y}$, as we obtain from Proposition 1.3 for every $\zeta \in \mathcal{U}, x \in D$, and $\boldsymbol{y} \in \mathbf{Y}$

$$
\begin{aligned}
\left|A_{i j}(\boldsymbol{\zeta} ; x, \boldsymbol{y})\right| & \leq\left|\bar{A}_{i j}(x, \boldsymbol{y})\right|+\sum_{m=1}^{\infty}\left|\left(\Psi_{m}\right)_{i j}(x, \boldsymbol{y})\right|\left(1+r_{m}\right) \\
& \leq \operatorname{esssup}_{(x, \boldsymbol{y}) \in D \times \mathbf{Y}}\left|\bar{A}_{i j}(x, \boldsymbol{y})\right|+\sum_{m=1}^{J_{0}}\left\|\left(\Psi_{m}\right)_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}(1+K) \\
& +\sum_{j>J_{0}}\left(2+\frac{\alpha \nu_{m}}{M\left|\nu_{F}\right| \beta_{m}}\right)\left\|\left(\Psi_{m}\right)_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}
\end{aligned}
$$

From Proposition 1.3, we find that

$$
\begin{equation*}
\left|A_{i j}(\boldsymbol{\zeta} ; x, \boldsymbol{y})\right| \leq\left\|\bar{A}_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}+4 \frac{\kappa}{1+\kappa} \alpha_{0}+\frac{2 \alpha}{M} \tag{3.3}
\end{equation*}
$$

In what follows, for functions taking values in $\mathbb{C}$ we still denote (with slight abuse of notation) by $\mathbf{V}$

$$
H_{0}^{1}(D) \times \prod_{i=1}^{n} L^{2}\left(D \times Y_{1} \times \ldots \times Y_{i-1} ; H_{\#}^{1}\left(Y_{i}\right)\right)
$$

Consider the complex parametric, one-scale limiting problem: given $\boldsymbol{\zeta} \in \mathcal{U}$, find $\boldsymbol{u} \in \mathbf{V}$ such that

$$
\begin{equation*}
b(\boldsymbol{\zeta} ; \boldsymbol{u}, \boldsymbol{v})=\int_{D} \int_{\mathbf{Y}} A(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \boldsymbol{\nabla} \boldsymbol{u} \cdot \overline{\boldsymbol{\nabla} \boldsymbol{v}} d \boldsymbol{y} d x=\int_{D} f \overline{v_{0}} d x \quad \forall \boldsymbol{v}=\left(v_{0},\left\{v_{i}\right\}\right) \in \mathbf{V} \tag{3.4}
\end{equation*}
$$

Proposition 3.4 Problem (3.4) admits a unique solution which is uniformly bounded in $\mathbf{V}$ for all $\boldsymbol{\zeta} \in \mathcal{U}$.
Proof We first show that the matrix function $A(\boldsymbol{\zeta} ; x, \boldsymbol{y})$ is uniformly bounded and coercive for all $\boldsymbol{\zeta} \in U$, $x \in D$ and $\boldsymbol{y} \in \mathbf{Y}$. To this end, we observe that for every in $\xi \in \mathbb{C}^{d}$ and every $\boldsymbol{\zeta} \in U$, we have:

$$
\begin{aligned}
\left|\xi^{H} A(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \xi\right| & \leq\left|\xi^{H} \bar{A}(x, \boldsymbol{y}) \xi\right|+\sum_{m=1}^{\infty}\left|\zeta_{m}\right|\left|\left(\xi^{H} \Psi_{m} \xi\right)\right| \\
& \leq\left(\beta_{0}+\sum_{m=1}^{J_{0}}(1+K) \beta_{m}+\sum_{m>J_{0}}\left(2+\frac{\alpha \nu_{m}}{M\left|\nu_{F}\right| \beta_{m}}\right) \beta_{m}\right)|\xi|^{2} \\
& \leq\left(\beta_{0}+\frac{2 \kappa}{1+\kappa} \alpha_{0}+\frac{\alpha}{M}\right)|\xi|^{2}
\end{aligned}
$$

To prove uniform coercivity, we note that for every $\zeta \in \mathcal{U}$ and for every $\xi \in \mathbb{C}^{d}$

$$
\begin{align*}
\Re\left(\xi^{H} A(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \xi\right) & \geq \Re\left(\xi^{H} \bar{A}(x, y) \xi\right)-\sum_{m=1}^{\infty}\left|\zeta_{m} \| \xi^{H} \Psi_{m} \xi\right| \\
& \geq\left(\alpha_{0}-\sum_{m=1}^{J_{0}}(1+K) \beta_{m}-\sum_{m>J_{0}}\left(2+\frac{\alpha \nu_{m}}{M\left|\nu_{F}\right| \beta_{m}}\right) \beta_{m}\right)|\xi|^{2} \\
& \geq\left(\alpha_{0}-\frac{\kappa}{1+\kappa} \alpha_{0}-K \sum_{m=1}^{J_{0}} \beta_{m}-2 \sum_{m>J_{0}} \beta_{m}-\sum_{m>J_{0}} \frac{\alpha \nu_{m}}{M\left|\nu_{F}\right|}\right)|\xi|^{2} \\
& \geq\left(\alpha-\frac{\alpha}{2 M}-\frac{\alpha}{3 M}-\frac{\alpha}{M}\right)|\xi|^{2} \\
& >\frac{\alpha}{2}|\xi|^{2} \tag{3.5}
\end{align*}
$$

if $M \geq 4$ where $\alpha=\alpha_{0} /(1+\kappa)$. The proposition then follows from the Lax-Milgram Lemma.
For an index $\nu \in \mathcal{F}$, we denote the support of $\nu$ by $\operatorname{supp}(\nu)$, i.e. the set of $j$ such that $\nu_{j} \neq 0$. We define the domain

$$
\mathcal{U}_{\nu}=\otimes_{j \in \operatorname{supp}(\nu)} \mathcal{U}_{j} .
$$

The following analyticity properties of $\boldsymbol{u}(\boldsymbol{z} ; \cdot, \cdot)$ hold.
Proposition 3.5 For $\nu \in \mathcal{F}$ and $\zeta \in \mathcal{U}$ with fixed $\zeta_{k}$ for all the indices $k \notin \operatorname{supp}(\nu)$, the map $\boldsymbol{u}: \mathcal{U}_{\nu} \rightarrow \mathbf{V}$ is analytic as a $\mathbf{V}$-valued function.
Proof For $m \in \mathbb{N}$, we fix all coordinates $\zeta_{k}$ for $k \neq m$ and partition each vector $\boldsymbol{\zeta} \in \mathbb{C}^{\mathbb{N}}$ as $\boldsymbol{\zeta}=\left(\zeta_{m}^{*}, \zeta_{m}\right)$. It is sufficient to show that there exists a function $\boldsymbol{v} \in \mathbf{V}$ such that for all $\boldsymbol{\zeta} \in \mathcal{U}$ holds

$$
\lim _{\delta \rightarrow 0}\left\|\frac{\boldsymbol{u}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; \cdot, \cdot\right)-\boldsymbol{u}(\boldsymbol{\zeta} ; \cdot, \cdot)}{\delta}-\boldsymbol{v}(\boldsymbol{\zeta} ; \cdot, \cdot)\right\|_{\mathbf{V}}=0
$$

For $\delta>0$, define the difference quotient $\boldsymbol{v}^{\delta}:=\delta^{-1}\left(\boldsymbol{u}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; \cdot, \cdot\right)-\boldsymbol{u}(\boldsymbol{\zeta} ; \cdot, \cdot)\right)$.

The function $\boldsymbol{v}^{\delta}$ is a weak solution of the parametric variational problem

$$
\int_{D} \int_{\mathbf{Y}} A(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \boldsymbol{\nabla} \boldsymbol{v}^{\delta} \cdot \overline{\boldsymbol{\nabla} \boldsymbol{w}} d \boldsymbol{y} d x=-\int_{D} \int_{\mathbf{Y}} \Psi_{m}(x, y) \boldsymbol{\nabla} \boldsymbol{u}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}\right) \cdot \overline{\boldsymbol{\nabla} \boldsymbol{w}} d \boldsymbol{y} d x, \quad \forall \boldsymbol{w} \in \mathbf{V}
$$

Let $\boldsymbol{v}$ denote the solution of the problem

$$
\int_{D} \int_{\mathbf{Y}} A(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \nabla \boldsymbol{v} \cdot \overline{\boldsymbol{\nabla} \boldsymbol{w}} d \boldsymbol{y} d x=-\int_{D} \int_{\mathbf{Y}} \Psi_{m}(x, y) \boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \cdot \overline{\boldsymbol{\nabla} \boldsymbol{w}} d \boldsymbol{y} d x, \quad \forall \boldsymbol{w} \in \mathbf{V}
$$

We deduce that for every $\boldsymbol{w} \in \mathbf{V}$ and every $\boldsymbol{\zeta} \in \mathcal{U}$ holds
$\int_{D} \int_{\mathbf{Y}} A(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \boldsymbol{\nabla}\left(\boldsymbol{v}^{\delta}-\boldsymbol{v}\right) \cdot \overline{\boldsymbol{\nabla} \boldsymbol{w}} d \boldsymbol{y} d x=-\int_{D} \int_{\mathbf{Y}} \Psi_{m}(x, y) \boldsymbol{\nabla}\left(\boldsymbol{u}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}\right)-\boldsymbol{u}(\boldsymbol{\zeta}, x, \boldsymbol{y})\right) \cdot \overline{\boldsymbol{\nabla} \boldsymbol{w}} d \boldsymbol{y} d x$.
From this we obtain

$$
\left\|\left\|\boldsymbol{v}^{\delta}-\boldsymbol{v}\right\|\right\| \leq c(m)\left\|\boldsymbol{u}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}\right)-\boldsymbol{u}(\boldsymbol{\zeta} ; x, \boldsymbol{y})\right\| \|
$$

From (1.22),

$$
\lim _{\delta \rightarrow 0}\| \| \boldsymbol{v}^{\delta}-\boldsymbol{v}\| \|=0
$$

Hartogs' theorem implies that for every $\nu \in \mathcal{F}, \boldsymbol{u}(\zeta)$ is analytic as a mapping from $\mathcal{U}_{\nu}$ to $\mathbf{V}$. This completes the proof.

We next investigate summability of the Legendre coefficients.

### 3.2 Coefficient estimates

Proposition 3.6 For every $\nu \in \mathcal{F}$, there holds

$$
\begin{equation*}
\left\|\boldsymbol{u}_{\nu}\right\|_{\mathbf{v}} \leq C\left(\prod_{m \in \operatorname{supp}(\nu)} \frac{2(1+K)}{K} \eta_{m}^{-\nu_{m}}\right) \tag{3.6}
\end{equation*}
$$

where $\eta_{m}:=r_{m}+\sqrt{1+r_{m}^{2}}$ with $r_{m}$ as in (3.1).
Proof We proceed as in the proof of Lemma A. 3 in Bieri, Andreev and Schwab [4] and Hoang and Schwab [11]. For $\nu \in \mathcal{F}$, the function $\boldsymbol{u}_{\nu} \in \mathbf{V}$ in (1.23) can be represented as

$$
\begin{equation*}
\boldsymbol{u}_{\nu}=\int_{U} \boldsymbol{u}(\boldsymbol{z}) L_{\nu}(\boldsymbol{z}) \rho(d \boldsymbol{z}) \tag{3.7}
\end{equation*}
$$

where the integral is understood as a Bochner integral of $\mathbf{V}$-valued functions. Let $S=\operatorname{supp}(\nu) \subset \mathbb{N}$ and define $\bar{S}:=\mathbb{N} \backslash S$. We then denote by $\mathcal{U}_{S}=\otimes_{m \in S} \mathcal{U}_{m}$ and $\mathcal{U}_{\bar{S}}=\otimes_{m \in \bar{S}} \mathcal{U}_{m}$, and by $\boldsymbol{z}_{S}=\left\{z_{i}, i \in S\right\}$, $\boldsymbol{z}_{\bar{S}}=\left\{z_{i}, i \in \bar{S}\right\}$ the extraction from $\boldsymbol{z}$, and analogously $\boldsymbol{\zeta}_{S}$ and $\boldsymbol{\zeta}_{\bar{S}}$.

Let $\mathcal{E}_{m}$ be the ellipse in $\mathcal{U}_{m}$ with foci at $\pm 1$ and the sum of the semiaxes being $\eta_{m}$; and $\mathcal{E}_{S}=$ $\prod_{m \in \operatorname{supp}(\nu)} \mathcal{E}_{m}$. We can then write (3.7) as

$$
\boldsymbol{u}_{\nu}=\frac{1}{(2 \pi i)^{|\nu|_{0}}} \int_{U} L_{\nu}(\boldsymbol{z}) \oint_{\mathcal{E}_{S}} \frac{\boldsymbol{u}\left(\boldsymbol{\zeta}_{S}, \boldsymbol{z}_{\bar{S}}\right)}{\left(\boldsymbol{\zeta}_{S}-\boldsymbol{z}_{S}\right)^{\mathbf{1}}} d \boldsymbol{\zeta}_{S} d \rho(\boldsymbol{z})
$$

For each $m \in \mathbb{N}$, let $\Gamma_{m}$ be a copy of $[-1,1]$ and $z_{m} \in \Gamma_{m}$. We denote by $U_{S}=\prod_{m \in S} \Gamma_{m}$ and $U_{\bar{S}}=\prod_{m \in \bar{S}} \Gamma_{m}$. We then have

$$
\boldsymbol{u}_{\nu}=\frac{1}{(2 \pi i)^{|\nu|_{0}}} \int_{U_{\bar{S}}} \oint_{\mathcal{E}_{S}} \boldsymbol{u}\left(\boldsymbol{\zeta}_{S}, \boldsymbol{z}_{\bar{S}}\right) \int_{U_{S}} \frac{L_{\nu}(\boldsymbol{z})}{\left(\boldsymbol{\zeta}_{S}-\boldsymbol{z}_{S}\right)^{\mathbf{1}}} d \rho_{S}\left(\boldsymbol{z}_{S}\right) d \boldsymbol{\zeta}_{S} d \rho_{\bar{S}}\left(\boldsymbol{z}_{\bar{S}}\right)
$$

We recall the definitions of the Legendre functions of the second kind:

$$
Q_{n}(\xi)=\frac{1}{2} \int_{[-1,1]} \frac{L_{n}(z)}{(\xi-z)} d z, \quad \xi \in \mathbb{C} \backslash[-1,1], n \in \mathbb{N}_{0}
$$

For $\nu \in \mathcal{F}$, we denote by $\nu_{S}$ the restriction of $\nu$ to $S$. We define for $\zeta \in \mathbb{C}^{\mathbb{N}}$

$$
\mathcal{Q}_{\nu_{S}}\left(\boldsymbol{\zeta}_{S}\right)=\prod_{m \in \operatorname{supp}(\nu)} Q_{\nu_{m}}\left(\zeta_{m}\right)
$$

Under the Joukovski transformation $\zeta_{m}=\frac{1}{2}\left(w_{m}+w_{m}^{-1}\right)$, the Legendre polynomials of the second kind are written as

$$
Q_{\nu_{m}}\left(\frac{1}{2}\left(w_{m}+w_{m}^{-1}\right)\right)=\sum_{k=\nu_{m}+1}^{\infty} \frac{q_{\nu_{m} k}}{w_{m}^{k}}
$$

with $\left|q_{\nu_{m} k}\right| \leq \pi$. Therefore

$$
\left|\mathcal{Q}_{\nu_{S}}\left(\boldsymbol{\zeta}_{S}\right)\right| \leq \prod_{m \in S} \sum_{k=\nu_{m}+1}^{\infty} \frac{\pi}{\eta_{m}^{k}}=\prod_{m \in S} \pi \frac{\eta_{m}^{-\nu_{m}-1}}{1-\eta_{m}^{-1}}
$$

We then have

$$
\begin{aligned}
\left\|\boldsymbol{u}_{\nu}\right\|_{\mathbf{V}} & =\left\|\frac{1}{(2 \pi i)^{|\nu|_{0}}} \int_{U_{\bar{S}}} \oint_{\mathcal{E}_{S}} \boldsymbol{u}\left(\boldsymbol{\zeta}_{S}, \boldsymbol{z}_{\bar{S}}\right) \mathcal{Q}_{\nu_{S}}\left(\boldsymbol{\zeta}_{S}\right) d \boldsymbol{\zeta}_{S} d \rho_{\bar{S}}\left(\overline{\boldsymbol{z}}_{S}\right)\right\|_{\mathbf{V}} \\
& \leq \frac{1}{(2 \pi)^{|\nu|_{0}}} \int_{U_{\bar{S}}} \oint_{\mathcal{E}_{S}}\left\|\boldsymbol{u}\left(\boldsymbol{\zeta}_{S}, \boldsymbol{z}_{\bar{S}}\right)\right\|_{\mathbf{V}} \mathcal{Q}_{\nu_{S}}\left(\boldsymbol{\zeta}_{S}\right) d \boldsymbol{\zeta}_{S} d \rho_{\bar{S}}\left(\overline{\boldsymbol{z}}_{S}\right) \\
& \leq \frac{1}{(2 \pi)^{|\nu|_{0}}\|\boldsymbol{u}(\boldsymbol{\zeta})\|_{L^{\infty}\left(\mathcal{E}_{S} \times U_{\bar{S}}, \mathbf{V}\right)} \max _{\mathcal{E}_{S}}\left|\mathcal{Q}_{\nu_{S}}\right| \prod_{m \in S} \operatorname{Len}\left(\mathcal{E}_{m}\right)} \\
& \leq \frac{1}{(2 \pi)^{|\nu|_{0}}}\|\boldsymbol{u}(\boldsymbol{\zeta})\|_{L^{\infty}\left(\mathcal{E}_{S} \times \mathcal{U}_{\bar{S}}, \mathbf{V}\right)} \prod_{m \in S} \pi \frac{\eta_{m}^{-\nu_{m}-1}}{1-\eta_{m}^{-1}} \operatorname{Len}\left(\mathcal{E}_{m}\right) \\
& \leq C \prod_{m \in S} \frac{2(1+K)}{K} \eta_{m}^{-\nu_{m}},
\end{aligned}
$$

as $\operatorname{Len}\left(\mathcal{E}_{m}\right) \leq 4 \eta_{m}, \eta_{m} \geq 1+K$ and $\boldsymbol{u}(\boldsymbol{\zeta})$ is uniformly bounded in $\mathbf{V}$.
To show the $\ell^{p}(\mathcal{F})$ summability of $\left\|\boldsymbol{u}_{\nu}\right\|_{\mathbf{V}}$, we use the following proposition, whose proof can be found in [7].
Proposition 3.7 For $0<p<1$, $\left(\frac{|\nu|!}{\nu!} b^{\nu}\right)_{\nu \in \mathcal{F}} \in \ell^{p}(\mathcal{F})$ iff (i) $\sum_{m \geq 1} b_{m}<1$ and (ii) $\left(b_{m}\right) \in \ell^{p}(\mathbb{N})$.
Proposition 3.8 For $0<p<1$ as in Assumption 3.2, $\left(\left\|u_{\nu}\right\|_{\mathbf{V}}\right)_{\nu \in \mathcal{F}} \in \ell^{p}(\mathcal{F})$.
Proof We have from the previous proposition that

$$
\begin{aligned}
\left\|\boldsymbol{u}_{\nu}\right\|_{\mathbf{v}} & \leq C \prod_{m \in S} \frac{2(1+K)}{K}\left(1+r_{m}\right)^{-\nu_{m}} \\
& \leq C\left(\prod_{m \in E, \nu_{m} \neq 0} \frac{2(1+K)}{K} \eta^{\nu_{m}}\right)\left(\prod_{m \in F, \nu_{m} \neq 0} \frac{2(1+K)}{K}\left(\frac{M\left|\nu_{F}\right| \beta_{m}}{\alpha \nu_{m}}\right)^{\nu_{m}}\right)
\end{aligned}
$$

where $\eta=1 /(1+K)$. Let $\mathcal{F}_{E}=\{\nu \in \mathcal{F}: \quad \operatorname{supp}(\nu) \subset E\}$ and $\mathcal{F}_{F}=\mathcal{F} \backslash E$. From this, we have

$$
\sum_{\nu \in \mathcal{F}}\left\|\boldsymbol{u}_{\nu}\right\|_{\mathrm{V}}^{p} \leq C A_{E} A_{F}
$$

where

$$
A_{E}=\sum_{\nu \in \mathcal{F}_{E}} \prod_{m \in E, \nu_{m} \neq 0}\left(\frac{2(1+K)}{K}\right)^{p} \eta^{p \nu_{m}}
$$

and

$$
A_{F}=\sum_{\nu \in \mathcal{F}_{F}} \prod_{m \in F, \nu_{m} \neq 0}\left(\frac{2(1+K)}{K}\right)^{p}\left(\frac{M|\nu| \beta_{m}}{\alpha \nu_{m}}\right)^{p \nu_{m}}
$$

We now show that both $A_{E}$ and $A_{F}$ are finite. For $A_{E}$, we have

$$
A_{E}=\left(1+\left(\frac{2(1+K)}{K}\right)^{p} \sum_{m \geq 1} \eta^{p m}\right)^{J_{0}}
$$

which is finite because $\eta<1$. For $A_{F}$, we note that for $\nu_{m} \neq 0$,

$$
\frac{2(1+K)}{K} \leq\left(\frac{2(1+K)}{K}\right)^{\nu_{m}}
$$

Therefore

$$
A_{F} \leq \sum_{\nu \in \mathcal{F}_{F}} \prod_{m \in F}\left(\frac{|\nu| d_{m}}{\nu_{m}}\right)^{p \nu_{m}}, \text { where } d_{m}=\frac{2 M(1+K) \beta_{m}}{K \alpha}
$$

and where we made the convention that $0^{0}=1$. We now proceed as in [8]: from the Stirling estimate $n!e^{n} /(e \sqrt{n}) \leq n^{n} \leq n!e^{n} / \sqrt{2 \pi n}$, we infer $|\nu|^{|\nu|} \leq|\nu|!e^{|\nu|}$ and obtain

$$
\prod_{m \in F} \nu_{m}^{\nu_{m}} \geq \frac{\nu!e^{|\nu|}}{\prod_{m \in F} \max \left\{1, e \sqrt{\nu_{m}}\right\}}
$$

Hence

$$
A_{F} \leq \sum_{\nu \in \mathcal{F}_{F}}\left(\frac{|\nu|!}{\nu!} d^{\nu}\right)^{p}\left(\prod_{m \in F} \max \left\{1, e \sqrt{\nu_{m}}\right\}\right)^{p} \leq \sum_{\nu \in \mathcal{F}_{F}}\left(\frac{|\nu|!}{\nu!} \bar{d}^{\nu}\right)^{p}
$$

where $\bar{d}_{m}=e d_{m}$ and where we have used the estimate $e \sqrt{n} \leq e^{n}$. From this, we have

$$
\sum_{m \geq 1} \bar{d}_{m} \leq \sum_{m \in F} \frac{6 M(1+K) \beta_{m}}{K \alpha}<1
$$

It is also obvious that

$$
\|\bar{d}\|_{\ell^{p}(\mathbb{N})}<\infty
$$

From these estimates and from Proposition 3.7 we obtain the conclusion.

### 3.3 Best $N$-term Approximation Rates

With Lemma 3.1, we have from Proposition 3.8 and Theorem 2.1 the following result:
Theorem 3.9 If Assumptions 1.1, 1.2 and 3.2 hold for some $0<p<1$, there exists a sequence $\left(\Lambda_{N}\right)_{N \in \mathbb{N}} \subset \mathcal{F}$ of index sets with cardinality not exceeding $N$ such that the solutions $\boldsymbol{u}_{\Lambda_{N}}$ of the Galerkin semidiscretized problems (2.3) satisfy

$$
\left\|\boldsymbol{u}-\boldsymbol{u}_{\Lambda_{N}}\right\|_{\underline{\mathbf{v}}} \leq C N^{-\sigma}, \quad \sigma=\frac{1}{p}-\frac{1}{2}
$$

## 4 Regularity

To obtain convergence rates of sparse tensor finite element discretizations for the fully discretized problem of (2.3), we introduce, following [15, 10], regularity spaces $\mathcal{H}_{i}(i=1, \ldots, n)$. The space $\mathcal{H}_{i}$ consists of all the functions $w\left(x, y_{1}, \ldots, y_{i}\right)$ that are $Y_{j}$-periodic in $y_{j}(j=1, \ldots, i)$ such that for any vectors $\left(\alpha_{0}, \alpha_{1}, \ldots, \alpha_{i}\right) \in\left(\mathbb{N}_{0}^{d}\right)^{i+1}$ such that $\left|\alpha_{j}\right| \leq 1$ for $0 \leq j \leq i-1$ and $\left|\alpha_{i}\right| \leq 2$ where $\left|\alpha_{i}\right|$ denotes the sum of all the components of $\alpha_{i}$,

$$
\frac{\partial^{\left|\alpha_{0}\right|+\ldots+\left|\alpha_{i}\right|} w}{\partial^{\alpha_{0}} x \partial^{\alpha_{1}} y_{1} \ldots \partial^{\alpha_{i}} y_{i}} \in L^{2}\left(D \times Y_{1} \times \ldots \times Y_{i}\right)
$$

The space $\mathcal{H}_{i}$ is equipped with the norm

$$
\|w\|_{\mathcal{H}_{i}}=\sum_{\substack{\left|\alpha_{i}\right| \leq 2 \\\left|\alpha_{j}\right| \leq 1, j=0, \ldots, i-1}}\left\|\frac{\partial^{\left|\alpha_{0}\right|+\ldots+\left|\alpha_{i}\right|} w}{\partial^{\alpha_{0}} x \partial^{\alpha_{1}} y_{1} \ldots \partial^{\alpha_{i}} y_{i}}\right\|_{L^{2}\left(D \times Y_{1} \times \ldots \times Y_{i}\right)} .
$$

We then define the subspace $\mathcal{H}$ of $\mathbf{V}$ as

$$
\mathcal{H}=\left\{\left(v_{0},\left\{v_{i}\right\}\right): v_{0} \in H^{2}(D), v_{i} \in \mathcal{H}_{i}, i=1, \ldots, n\right\}
$$

### 4.1 Regularity of the parametric, deterministic problem (1.19)

For each index $i=1, \ldots, n-1$, we denote

$$
\begin{equation*}
\boldsymbol{y}_{i}=\left(y_{1}, \ldots, y_{i}\right) \quad \text { and } \quad \mathbf{Y}_{i}=Y_{1} \times \ldots \times Y_{i} \tag{4.1}
\end{equation*}
$$

We define by $C_{i}^{1}, i=1, \ldots, n$ the space of functions $w\left(x, y_{1}, \ldots, y_{i}\right)$ that are continuous in each variables $x, y_{1}, \ldots, y_{i}$ and that are $Y_{j}$-periodic with respect to $y_{j}, j=1, \ldots, i$. For a vector $\left(\gamma_{0}, \ldots, \gamma_{i}\right) \in\{0,1\}^{i+1}$ and the index vector $\left(j_{0}, j_{1}, \ldots, j_{i}\right) \in\{1, \ldots, d\}^{i+1}$, the strong derivative

$$
\frac{\partial^{\gamma_{0}+\ldots+\gamma_{i}} w}{\partial x_{j_{0}}^{\gamma_{0}} \partial y_{1 j_{1}}^{\gamma_{1}} \ldots \partial y_{i j_{i}}^{\gamma_{i}}}
$$

exists for all $\left(x, y_{1}, \ldots, y_{i}\right) \in \bar{D} \times \bar{Y}_{1} \times \ldots \times \bar{Y}_{i}$ and is continuous. We define the seminorm

$$
\begin{equation*}
\|w\|_{C_{i}^{1}}=\sum_{\substack{\left(\gamma_{0}, \ldots, \gamma_{i}\right) \in\{0,1\}^{i+1} \\\left(j_{0}, j_{1}, \ldots, j_{i}\right) \in\{1, \ldots, d\}^{i+1}}}\left\|\frac{\partial^{\gamma_{0}+\ldots+\gamma_{i}} w}{\partial x_{j_{0}}^{\gamma_{0}} \partial y_{1 j_{1}}^{\gamma_{1}} \ldots \partial y_{i j_{i}}^{\gamma_{i}}}\right\|_{L^{\infty}\left(\Omega \times \mathbf{Y}_{i}\right)} \tag{4.2}
\end{equation*}
$$

The following homogenization result is, in principle, well known (see, e.g., [2]). As we require its parametric version, and also use its derivation later, we present its proof.

Proposition 4.1 There exists a symmetric matrix function $A_{0}(\boldsymbol{z} ; \cdot) \in L^{\infty}(D)_{\text {sym }}^{d \times d}$ that is uniformly bouunded and coercive for all $\boldsymbol{z} \in U$ such that the limit function $u_{0}(\boldsymbol{z}, \cdot) \in H_{0}^{1}(D)$ in Theorem 1.12 is the solution of the problem:

$$
\int_{D} A_{0}(\boldsymbol{z} ; x) \nabla u_{0}(\boldsymbol{z} ; x) \cdot \nabla \phi(x) d x=\int_{D} f(x) \phi(x) d x, \quad \forall \phi \in H_{0}^{1}(D)
$$

Proof With $\boldsymbol{u}(\boldsymbol{z})=\left(u_{0}, u_{1}, \ldots, u_{n}\right) \in \mathbf{V}$ as in (1.19), we have (with implied summation over the repeated index $l=1, \ldots, d)$

$$
u_{n}=w_{n l}\left(\frac{\partial u_{0}}{\partial x_{l}}+\frac{\partial u_{1}}{\partial y_{1 l}}+\ldots+\frac{\partial u_{n-1}}{\partial y_{(n-1) l}}\right)
$$

where the functions $w_{n l} \in L^{2}\left(D \times Y_{1} \times \ldots \times Y_{n-1} ; H_{\#}^{1}\left(Y_{n}\right) / \mathbb{R}\right)$ are the unique solutions of the parametric unit-cell problems

$$
\begin{equation*}
\int_{Y_{n}} A\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}, y_{n}\right)\left(e_{l}+\nabla_{y_{n}} w_{n l}\right) \cdot \nabla_{y_{n}} \phi_{n} d y_{n}=0, \quad \forall \phi_{n} \in H_{\#}^{1}\left(Y_{n}\right) \tag{4.3}
\end{equation*}
$$

(here $e_{l}$ denotes the $l$ th unit vector in $\mathbb{R}^{d}$ ). From (1.19), we have

$$
\begin{equation*}
\int_{D} \int_{\mathbf{Y}} A\left(I+\nabla_{y_{n}} w_{n}\right) \cdot\left(\nabla_{x} u_{0}+\sum_{k=1}^{n-1} \nabla_{y_{k}} u_{k}\right) \cdot \nabla_{y_{n-1}} \phi_{n-1} d \boldsymbol{y} d x=0 \tag{4.4}
\end{equation*}
$$

for all $\phi_{n-1} \in L^{2}\left(D \times Y_{1} \times \ldots \times Y_{n-2}, H_{\#}^{1}\left(Y_{n-1}\right)\right)$, where $w_{n}$ denotes the vector $\left(w_{n 1}, \ldots, w_{n d}\right)$ and $I$ is the identity matrix. By recursion, we define the "upscaled" conductivity matrices $A_{n-1}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}\right)$ as

$$
\begin{equation*}
A_{n-1}=\int_{Y_{n}} A\left(I+\nabla_{y_{n}} w_{n}\right) d y_{n}=\int_{Y_{n}} A\left(I+\nabla_{y_{n}} w_{n}\right) \cdot\left(I+\nabla_{y_{n}} w_{n}\right) d y_{n} \tag{4.5}
\end{equation*}
$$

We then consider the parametric unit cell problem on scale $n-1$ : find $w_{(n-1) l}$ such that

$$
\int_{Y_{n-1}} A_{n-1}\left(e_{l}+\nabla_{y_{n-1}} w_{(n-1) l}\right) \cdot \nabla_{y_{n-1}} \phi_{n-1} d y_{n-1}=0, \quad \forall \phi_{n-1} \in H_{\#}^{1}\left(Y_{n-1}\right) .
$$

We then have

$$
u_{n-1}=w_{(n-1) l}\left(\frac{\partial u_{0}}{\partial x_{l}}+\frac{\partial u_{1}}{\partial y_{1 l}}+\ldots+\frac{\partial u_{n-2}}{\partial y_{(n-2) l}}\right)
$$

With the convention that $A_{n}=A$, we define recursively for $i=n-2, n-3, \ldots$ the functions $w_{i l} \in$ $L^{2}\left(D \times \mathbf{Y}_{i-1} ; H_{\#}^{1}\left(Y_{i}\right) / \mathbb{R}\right)$ as (unique) solutions of the problems

$$
\int_{Y_{i}} A_{i}\left(e_{l}+\nabla_{y_{i}} w_{i l}\right) \cdot \nabla_{y_{i}} \phi_{i} d y_{i}=0, \quad \forall \phi_{i} \in H_{\#}^{1}\left(Y_{i}\right) .
$$

For $i=1,2, \ldots, n$, the scale interaction function $u_{i}$ is then determined as

$$
u_{i}=w_{i l}\left(\frac{\partial u_{0}}{\partial x_{l}}+\frac{\partial u_{1}}{\partial y_{1 l}}+\ldots+\frac{\partial u_{i-1}}{\partial y_{(i-1) l}}\right)
$$

and the "upscaled" matrix $A_{i-1}$ is defined in terms of $A_{i}$ as

$$
\begin{equation*}
A_{i-1}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{i-1}\right)=\int_{Y_{i}} A_{i}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{i-1}, y_{i}\right)\left(I+\nabla_{y_{i}} w_{i}\right) \cdot\left(I+\nabla_{y_{i}} w_{i}\right) d y_{i}, \quad \boldsymbol{z} \in U, x \in D, \boldsymbol{y}_{i-1} \in \mathbf{Y}_{i-1} \tag{4.6}
\end{equation*}
$$

where $w_{i}$ denotes the vector $\left(w_{i 1}, \ldots, w_{i d}\right)$. Upon completing the upscaling recursion at $i=1$ the effective diffusivity matrix $A_{0}(x)$ is obtained as

$$
A_{0}(\boldsymbol{z} ; x)=\int_{Y_{1}} A_{1}\left(\boldsymbol{z} ; x, y_{1}\right)\left(I+\nabla_{y_{1}} w_{1}\right) \cdot\left(I+\nabla_{y_{1}} w_{1}\right) d y_{1}
$$

and the function $u_{0}(\boldsymbol{z} ; \cdot) \in H_{0}^{1}(D)$ satisfies the homogenized, parametric limiting problem

$$
\begin{equation*}
\int_{D} A_{0}(\boldsymbol{z} ; x) \nabla u_{0}(\boldsymbol{z} ; x) \cdot \nabla \phi(x) d x=\int_{D} f(x) \phi(x) d x \tag{4.7}
\end{equation*}
$$

for all $\phi \in H_{0}^{1}(D)$.
As the matrix $A$ is symmetric, all matrices $A_{i}(i=0, \ldots, n-1)$ are symmetric. Fix $\xi \in \mathbb{R}^{d}$. Then (with summation over repeated indices)

$$
A_{(n-1) k l} \xi_{k} \xi_{l}=\int_{Y_{n}} A_{r s}\left(\xi_{r}+\frac{\partial\left(w_{n k} \xi_{k}\right)}{\partial y_{n r}}\right)\left(\xi_{s}+\frac{\partial\left(w_{n l} \xi_{l}\right)}{\partial y_{n s}}\right) d y_{n}
$$

For the constant $\alpha$ as in Assumption 1.1, and for every $\boldsymbol{z} \in U, x \in D, \boldsymbol{y}_{n-1} \in \mathbf{Y}_{n-1}$ and every $\xi \in \mathbb{R}^{d}$

$$
A_{(n-1) k l}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}\right) \xi_{k} \xi_{l} \geq \alpha \int_{Y_{n}}\left(\xi_{r}+\frac{\partial\left(w_{n k} \xi_{k}\right)}{\partial y_{n r}}\right)\left(\xi_{r}+\frac{\partial\left(w_{n l} \xi_{l}\right)}{\partial y_{n r}}\right) d y_{n} \geq \alpha|\xi|^{2}
$$

Furthermore with summation over repeated indices,

$$
\begin{aligned}
A_{(n-1) k l}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}\right) \xi_{k} \xi_{l} & \leq \beta\left(\sum_{r=1}^{d} \xi_{r}^{2}+\sum_{r=1}^{d} \int_{Y_{n}} \frac{\partial\left(w_{n k} \xi_{k}\right)}{\partial y_{n r}} \frac{\partial\left(w_{n l} \xi_{l}\right)}{\partial y_{n r}} d y_{n}\right) \\
& \leq \beta\left(\sum_{r=1}^{d} \xi_{r}^{2}+\sum_{r=1}^{d}\left(\sum_{k=1}^{d} \xi_{k}^{2}\right)\left(\sum_{k=1}^{d} \int_{Y_{n}}\left(\frac{\partial w_{n k}}{\partial y_{n r}}\right)^{2} d y_{n}\right)\right)
\end{aligned}
$$

From (4.3), we deduce that there is a constant $c=c(d)$ which depends only on the dimension $d$ such that

$$
\left\|\nabla_{y_{n}} w_{n l}\right\|_{L^{2}\left(Y_{n}\right)} \leq \frac{c(d)}{\alpha} \sup _{k, l}\left\|A_{k l}\right\|_{L^{\infty}(D \times \mathbf{Y})}
$$

Therefore, there is a constant $c=c(\alpha, d)$ such that

$$
A_{(n-1) k l}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}\right) \xi_{k} \xi_{l} \leq \beta c(\alpha, d)\left(1+\sup _{k, l}\left\|A_{k l}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right)^{2}|\xi|^{2}
$$

Repeating this argument for $A_{i}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{i}\right), i=n-1, \ldots, 1$, we deduce that for all $z \in U$ and $x \in D$,

$$
A_{0 k l}(z ; x) \xi_{k} \xi_{l} \geq \alpha|\xi|^{2}
$$

and

$$
\begin{array}{r}
A_{0 k l}(z ; x) \xi_{k} \xi_{l} \leq \beta c(\alpha, d)^{n}\left(1+\sup _{k, l}\left\|A_{k l}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right)^{2}\left(1+\sup _{k, l}\left\|A_{(n-1) k l}\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}\right)}\right)^{2} \cdots \\
.\left(1+\sup _{k, l}\left\|A_{1 k l}\right\|_{L^{\infty}(D \times Y)}\right)^{2}|\xi|^{2} .
\end{array}
$$

From (4.5), we deduce that

$$
\begin{aligned}
\sup _{k, l}\left\|A_{(n-1) k l}\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}\right)} & \leq c(d) \sup _{k, l}\left\|A_{k l}\right\|_{L^{\infty}(D \times \mathbf{Y})}\left(1+\left\|\nabla_{y_{n}} w_{n}\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}, L^{2}\left(Y_{n}\right)\right)}\right) \\
& \leq c(\alpha, d)\left(1+\sup _{k, l}\left\|A_{k l}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right)^{2}
\end{aligned}
$$

so

$$
1+\sup _{k, l}\left\|A_{(n-1) k l}\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}\right)} \leq c(\alpha, d)\left(1+\sup _{k, l}\left\|A_{k l}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right)^{2}
$$

Repeating this argument for $i=2, \ldots, n$, we get

$$
1+\sup _{k, l}\left\|A_{(n-i) k l}\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-i}\right)} \leq c(\alpha, d)^{2^{i}-1}\left(1+\sup _{k, l}\left\|A_{k l}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right)^{2^{i}} .
$$

Therefore

$$
\begin{equation*}
A_{0 k l}(z ; x) \xi_{k} \xi_{l} \leq \beta\left[c(\alpha, d)\left(1+\sup _{k, l}\left\|A_{k l}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right]^{2^{n+1}-2}|\xi|^{2} .\right. \tag{4.8}
\end{equation*}
$$

Proposition 4.2 Assume that the domain $D$ is convex and $f \in L^{2}(D)$. Assume further that $A(\boldsymbol{z}) \in \boldsymbol{C}_{n}^{1}$ and $\|A(\boldsymbol{z})\|_{\boldsymbol{C}_{n}^{1}}$ is uniformly bounded for all $\boldsymbol{z} \in U$. Then $\boldsymbol{u}(\boldsymbol{z}) \in \mathcal{H}$ and $\|\boldsymbol{u}(\boldsymbol{z})\|_{\mathcal{H}}$ is uniformly bounded for all $\boldsymbol{z} \in U$.

Proof The functions $u_{i}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{i}\right)$ can be expressed in terms of the functions $w_{i}=\left(w_{i 1}, \ldots, w_{i d}\right)$ as

$$
\begin{equation*}
u_{i}=w_{i} \cdot\left(I+\nabla_{y_{i-1}} w_{i-1}\right) \cdots\left(I+\nabla_{y_{1}} w_{1}\right) \cdot \nabla u_{0} \tag{4.9}
\end{equation*}
$$

From (4.3), for almost all $\left(x, \boldsymbol{y}_{n-1}\right) \in D \times \mathbf{Y}_{n-1}$

$$
\begin{equation*}
\int_{Y_{n}} A(\boldsymbol{z} ; x, \boldsymbol{y}) \nabla_{y_{n}} w_{n l}(\boldsymbol{z} ; x, \boldsymbol{y}) \cdot \nabla_{y_{n}} \phi_{n} d y_{n}=\int_{Y_{n}}\left(\nabla_{y_{n}} \cdot\left(A(\boldsymbol{z} ; x, \boldsymbol{y}) e_{l}\right) \phi_{n} d y_{n}, \quad \forall \phi_{n} \in H_{\#}^{1}\left(Y_{n}\right)\right. \tag{4.10}
\end{equation*}
$$

As any function in $\mathcal{D}\left(\mathbb{R}^{d}\right)$ with a sufficiently small support can be extended to a $Y_{n}$-periodic function of the same regularity, we see using a partition of unity, that

$$
\int_{\mathbb{R}^{d}} A(\boldsymbol{z} ; x, \boldsymbol{y}) \nabla_{y_{n}} w_{n l}(\boldsymbol{z} ; x, \boldsymbol{y}) \cdot \nabla_{y_{n}} \phi_{n} d y_{n}=\int_{\mathbb{R}^{d}}\left(\nabla_{y_{n}} \cdot\left(A(\boldsymbol{z} ; x, \boldsymbol{y}) e_{l}\right) \phi_{n} d y_{n}, \quad \forall \phi_{n} \in \mathcal{D}\left(\mathbb{R}^{d}\right)\right.
$$

We choose a smooth domain $D^{\prime}$ such that $Y_{n} \subset D^{\prime}$ and $\tau \in \mathcal{D}\left(D^{\prime}\right)$ such that $\tau\left(y_{n}\right)=1$ when $y_{n} \in Y_{n}$. For $\tau\left(y_{n}\right) w_{n l}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}, y_{n}\right)$ in $D^{\prime}$, we deduce that

$$
\left\|w_{n l}(\boldsymbol{z} ; x, \boldsymbol{y})\right\|_{H^{2}\left(Y_{n}\right)} \leq C\left(\left\|\nabla_{y_{n}} \cdot\left(A(\boldsymbol{z} ; x, \boldsymbol{y}) e_{l}\right)\right\|_{L^{2}\left(Y_{n}\right)}+\left\|w_{n l}\right\|_{L^{2}\left(Y_{n}\right)}\right)
$$

where the constant $C$ depends on the $\boldsymbol{C}^{1}$ norm of $A\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}, \cdot\right), \alpha, \beta$ and $\tau$, and is in particular independent of $\boldsymbol{z} \in U$ (see, e.g., Wloka [19] page 330).

Now, we freeze all the coordinates $\left(x, \boldsymbol{y}_{n-1}\right)$ except the $j$ th coordinate of the variable $y_{k}$ for an index $k=1, \ldots, n-1$, and denote by $\left(y_{k j}^{*}, y_{k j}\right)$ the vector $\boldsymbol{y}_{n-1}$. For $\delta>0$, let

$$
\chi^{\delta}(\boldsymbol{z} ; x, \boldsymbol{y})=\frac{w_{n l}\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-w_{n l}(\boldsymbol{z} ; x, \boldsymbol{y})}{\delta}
$$

For all $\phi_{n} \in H_{\#}^{1}\left(Y_{n}\right)$ we have

$$
\begin{array}{r}
\int_{Y_{n}} A(\boldsymbol{z} ; x, \boldsymbol{y}) \nabla_{y_{n}} \chi^{\delta} \cdot \nabla_{y_{n}} \phi_{n} d y_{n}=-\int_{Y_{n}} \frac{A\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})}{\delta} e_{l} \cdot \nabla_{y_{n}} \phi_{n} d y_{n} \\
\quad-\int_{Y_{n}} \frac{A\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})}{\delta} \nabla_{y_{n}} w_{n l}\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta, \cdot\right) \cdot \nabla_{y_{n}} \phi_{n} d y_{n} . \tag{4.11}
\end{array}
$$

Let $\chi\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}, \cdot\right) \in H_{\#}^{1}\left(Y_{n}\right) / \mathbb{R}$ denote the solution of the problem

$$
\begin{align*}
\int_{Y_{n}} A(\boldsymbol{z} ; x, \boldsymbol{y}) \nabla_{y_{n}} \chi & \cdot \nabla_{y_{n}} \phi_{n} d y_{n}=-\int_{Y_{n}} \frac{\partial A(\boldsymbol{z} ; x, \boldsymbol{y})}{\partial y_{k j}} e_{l} \cdot \nabla_{y_{n}} \phi_{n} d y_{n} \\
& -\int_{Y_{n}} \frac{\partial A(\boldsymbol{z} ; x, \boldsymbol{y})}{\partial y_{k j}} \nabla_{y_{n}} w_{n l}(\boldsymbol{z} ; x, \boldsymbol{y}) \cdot \nabla_{y_{n}} \phi_{n} d y_{n} \tag{4.12}
\end{align*}
$$

From these equations, we deduce

$$
\begin{aligned}
& \int_{Y_{n}} A(\boldsymbol{z} ; x, \boldsymbol{y}) \nabla_{y_{n}}\left(\chi^{\delta}-\chi\right) \cdot \nabla_{y_{n}} \phi_{n} d y_{n} \\
& =-\int_{Y_{n}}\left(\frac{A\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})}{\delta}-\frac{\partial A(\boldsymbol{z} ; x, \boldsymbol{y})}{\partial y_{k j}}\right) e_{l} \cdot \nabla_{y_{n}} \phi_{n} d y_{n} \\
& -\int_{Y_{n}}\left(\frac{A\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})}{\delta}-\frac{\partial A(\boldsymbol{z} ; x, \boldsymbol{y})}{\partial y_{k j}}\right) \nabla_{y_{n}} w_{n l}\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta, \cdot\right) \cdot \nabla_{y_{n}} \phi_{n} d y_{n} \\
& -\int_{Y_{n}} \frac{\partial A(\boldsymbol{z} ; x, \boldsymbol{y})}{\partial y_{k j}} \nabla_{y_{n}}\left(w_{n l}\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-w_{n l}(\boldsymbol{z} ; x, \boldsymbol{y})\right) \cdot \nabla_{y_{n}} \phi_{n} d y_{n} .
\end{aligned}
$$

From (4.11) (ignoring the constant $\delta$ ), we have for every $\boldsymbol{z} \in U$ and $x \in D$

$$
\begin{align*}
& \left\|w_{n l}\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-w_{n l}(\boldsymbol{z} ; x, \boldsymbol{y})\right\|_{H^{2}\left(Y_{n}\right) / \mathbb{R}} \\
& \leq C\left(\left\|A\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})\right\|_{L^{\infty}\left(Y_{n}\right)}+\left\|\nabla_{y_{n}}\left(A\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})\right)\right\|_{L^{\infty}\left(Y_{n}\right)}\right) . \tag{4.13}
\end{align*}
$$

Therefore,

$$
\begin{aligned}
\left\|\chi^{\delta}-\chi\right\|_{H^{2}\left(Y_{n}\right) / \mathbb{R}} \leq & C\left(\left\|\nabla_{y_{n}}\left(\frac{A\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})}{\delta}-\frac{\partial A}{\partial y_{k j}}\right)\right\|_{L^{\infty}\left(Y_{n}\right)}\right. \\
& +\left\|\frac{A\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})}{\delta}-\frac{\partial A}{\partial y_{k j}}\right\|_{L^{\infty}\left(Y_{n}\right)} \\
& +\left\|A\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})\right\|_{L^{\infty}\left(Y_{n}\right)} \\
& \left.+\left\|\nabla_{y_{n}}\left(A\left(\boldsymbol{z} ; x, y_{k j}^{*}, y_{k j}+\delta\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})\right)\right\|_{L^{\infty}\left(Y_{n}\right)}\right)
\end{aligned}
$$

which converges to 0 when $\delta$ tends to 0 as $A \in \boldsymbol{C}_{n}^{1}$. Therefore

$$
\chi=\frac{\partial w_{n l}}{\partial y_{k j}} \text { in } H_{\#}^{2}\left(Y_{n}\right) / \mathbb{R} .
$$

As $\chi$ satisfies (4.12), for each $\boldsymbol{z} \in U, \chi$ as a map from $D \times \mathbf{Y}_{n-1}$ to $H^{2}\left(Y_{n}\right) / \mathbb{R}$ is continuous, due to the continuity of the coefficient $A(\boldsymbol{z} ; x, \boldsymbol{y})$ and due to the continuity of $w_{n l}$ as a map from $D \times \mathbf{Y}_{n-1}$ (from (4.13)).

Performing a similar procedure for the remaining functions $\partial w_{n l} / \partial y_{k j}$ and their derivatives, we find $w_{n l} \in C_{n-1}^{1}\left(H_{\#}^{2}\left(Y_{n}\right) / \mathbb{R}\right)$. Therefore, from (4.5), $A_{n-1} \in C_{n-1}^{1}$. In the same fashion, we deduce that $w_{i l} \in C_{i-1}^{1}\left(H_{\#}^{2}\left(Y_{i}\right) / \mathbb{R}\right)$ for all $i=1, \ldots, n$ and $l=1, \ldots, d$, and $\left\|w_{i l}(\boldsymbol{z})\right\|_{\boldsymbol{C}_{i-1}^{1}\left(H^{2}\left(Y_{i}\right) / \mathbb{R}\right)}$ is uniformly bounded for all $\boldsymbol{z}$. Therefore for every $\boldsymbol{z} \in U, A_{0}(\boldsymbol{z}) \in C^{1}(\bar{D})^{d \times d}$ and $\left\|A_{0}(\boldsymbol{z})\right\|_{\left(C^{1}(\bar{D})\right)^{d \times d}}$ is uniformly bounded for all $\boldsymbol{z} \in U$.

Next we claim $u_{0}(\boldsymbol{z}, \cdot) \in H^{2}(D)$ and that its $H^{2}(D)$ norm is uniformly bounded for all $\boldsymbol{z} \in U$. We have shown that for all vectors $\xi \in \mathbb{R}^{d}$, and for every $\boldsymbol{z} \in U$

$$
\alpha|\xi|^{2} \leq \xi^{\top} A_{0}(\boldsymbol{z}, x) \xi \leq \beta^{\prime}|\xi|^{2}
$$

where $\alpha>0$ is the constant in Assumption 1.1 and $\beta^{\prime}$ is a positive constant that depends only on $\alpha, \beta$, $n$ and $d$. The entries of $A_{0}(\boldsymbol{z}, x)$ are therefore uniformly bounded by a constant depending on $\alpha$ and $\beta^{\prime}$. As $D$ is convex, Theorem 3.2.1.2 of Grisvard [9] shows that for each $\boldsymbol{z} \in U, u_{0}(\boldsymbol{z}) \in H^{2}(D)$. The proofs of Lemma 3.1.3.2 and of Theorem 3.2.1.2 in [9] show that

$$
\left\|u_{0}(\boldsymbol{z}, \cdot)\right\|_{H^{2}(D)} \leq c\|f\|_{L^{2}(D)},
$$

where the constant $c$ depends on the $C^{1}(\bar{D})$ norms of $A_{0}$, and the $L^{\infty}(D)$ norms of the entries of the matrix $A_{0}^{-1 / 2}(\boldsymbol{z}, x)$ which can be bounded by $\alpha$ and $\beta^{\prime}$. Therefore $\left\|u_{0}(\boldsymbol{z}, \cdot)\right\|_{H^{2}(D)}$ is uniformly bounded for all $\boldsymbol{z} \in U$. As $\left\|w_{i l}\left(\boldsymbol{z}, \boldsymbol{y}_{i}\right)\right\|_{\boldsymbol{C}_{i-1}^{1}\left(H^{2}\left(Y_{i}\right) / \mathbb{R}\right)}$ is uniformly bounded for $\boldsymbol{z} \in U$, we get from (4.9) that $u_{i} \in \mathcal{H}_{i}$ and $\left\|u_{i}(\boldsymbol{z})\right\|_{\mathcal{H}_{i}}$ is uniformly bounded for all $\boldsymbol{z} \in U$. Hence $\|\boldsymbol{u}(\boldsymbol{z})\|_{\mathcal{H}}$ is uniformly bounded for all $z \in U$.

To establish the measurability of $\boldsymbol{u}$, we make the following assumptions.
Assumption 4.3 We assume that the matrices $\Psi_{k}$ in (1.7) are in $\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}$ such that for all $i, j=1, \ldots, d$

$$
\sum_{k=1}^{\infty}\left\|\Psi_{k}\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}<\infty
$$

Remark 4.4 When $\operatorname{Cov}[A]_{i j i^{\prime} j^{\prime}} \in H^{t+1}(D) \otimes H_{\#}^{t+1}\left(Y_{1}\right) \otimes \ldots \otimes H^{t+1}\left(Y_{n}\right)$ for a sufficiently large constant $t$, for any vectors $\left(\gamma_{0}, \gamma_{1}, \ldots, \gamma_{n}\right) \in\{0,1\}^{n+1}$ and any $\left(j_{0}, j_{1}, \ldots, j_{n}\right) \in\{1, \ldots, d\}^{n+1}$,

$$
\frac{\partial^{\gamma_{0}+\ldots+\gamma_{n}} \mathbf{C o v}[A]_{i j i^{\prime} j^{\prime}}\left(x, \boldsymbol{y}, x^{\prime}, \boldsymbol{y}^{\prime}\right)}{\partial x_{j_{0}}^{\gamma_{0}} \partial y_{1 j_{1}}^{\gamma_{1}} \ldots \partial y_{n j_{n}}^{\gamma_{n}}} \in H^{t}(D) \otimes H_{\#}^{t}\left(Y_{1}\right) \otimes \ldots \otimes H_{\#}^{t}\left(Y_{n}\right)
$$

We then deduce that

$$
\begin{equation*}
\left\|\Psi_{k}\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}} \leq c(\varepsilon) k^{(-t / d+\varepsilon)\left(1-2 t^{*} / t\right)} \tag{4.14}
\end{equation*}
$$

Assumption 4.3 holds when $t$ is sufficiently large.
Proposition 4.5 With Assumption 4.3, the function $\boldsymbol{u}$ as a map from $U$ to $\mathcal{H}$ is measurable.
Proof We first prove that there exists a constant $c$ such that for all $\boldsymbol{z}, \boldsymbol{z}^{\prime} \in U$,

$$
\begin{equation*}
\left\|\boldsymbol{u}(\boldsymbol{z})-\boldsymbol{u}\left(\boldsymbol{z}^{\prime}\right)\right\|_{\mathcal{H}} \leq c\left\|A(\boldsymbol{z} ; \cdot, \cdot)-A\left(\boldsymbol{z}^{\prime}, \cdot, \cdot\right)\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}} . \tag{4.15}
\end{equation*}
$$

From (4.3), we have for every fixed $\boldsymbol{z} \in U, x \in D$ and $\boldsymbol{y}_{n-1} \in \mathbf{Y}_{n-1}$

$$
\begin{aligned}
& \int_{Y_{n}} A(\boldsymbol{z} ; x, \boldsymbol{y}) \nabla_{y_{n}}\left(w_{n l}(\boldsymbol{z} ; x, \boldsymbol{y})-w_{n l}\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right) \cdot \nabla_{y_{n}} \phi_{n} d y_{n} \\
& =\int_{Y_{n}}\left[\nabla_{y_{n}} \cdot\left(\left(A(\boldsymbol{z} ; x, \boldsymbol{y})-A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right) e_{l}\right)\right. \\
& \left.\quad \quad-\nabla_{y_{n}}\left(\left(A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})\right) \nabla_{y_{n}} w_{n l}\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right)\right] \phi_{n} d y_{n}, \quad \forall \phi_{n} \in H_{\#}^{1}\left(Y_{n}\right) .
\end{aligned}
$$

As $\left\|w_{n l}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}, \cdot\right)\right\|_{H^{2}\left(Y_{n}\right) / \mathbb{R}}$ is uniformly bounded for all $\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}\right) \in U \times D \times \mathbf{Y}_{n-1}$, we obtain

$$
\begin{equation*}
\|\left(w_{n l}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}, \cdot\right)-w_{n l}\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}_{n-1}, \cdot\right)\left\|_{H^{2}\left(Y_{n}\right) / \mathbb{R}} \leq c\right\| A(\boldsymbol{z} ; \cdot, \cdot)-A\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right) \|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}\right. \tag{4.16}
\end{equation*}
$$

Similarly, from (4.12)

$$
\begin{aligned}
& \int_{Y_{n}} A(\boldsymbol{z} ; x, \boldsymbol{y}) \nabla_{y_{n}}\left(\chi(\boldsymbol{z} ; \cdot, \cdot)-\chi\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right) \cdot \nabla_{y_{n}} \phi_{n} d y_{n} \\
& =\int_{Y_{n}} \nabla_{y_{n}} \cdot\left(\frac{\partial A(\boldsymbol{z} ; x, \boldsymbol{y})-\partial A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)}{\partial y_{k j}} e_{l}\right) \phi_{n} d y_{n} \\
& +\int_{Y_{n}} \nabla_{y_{n}} \cdot\left(\frac{\partial A(\boldsymbol{z} ; x, \boldsymbol{y})-\partial A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)}{\partial y_{k j}} \nabla_{y_{n}} w_{n l}(\boldsymbol{z} ; x, \boldsymbol{y})\right) \phi_{n} d y_{n} \\
& +\int_{Y_{n}} \nabla_{y_{n}} \cdot\left(\left(A(\boldsymbol{z} ; x, \boldsymbol{y})-A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right) \nabla_{y_{n}} \chi\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right) \phi_{n} d y_{n}
\end{aligned}
$$

Therefore there exists a constant $c>0$ such that for all $\boldsymbol{z}, \boldsymbol{z}^{\prime} \in U$

$$
\left\|\chi\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}, \cdot\right)-\chi\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}, \cdot\right)\right\|_{H^{2}\left(Y_{n}\right) / \mathbb{R}} \leq C\left\|A(\boldsymbol{z} ; \cdot, \cdot)-A\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}
$$

Performing a similar procedure for the derivatives of $\chi$, we deduce that

$$
\left\|w_{n l}(\boldsymbol{z} ; x, \boldsymbol{y})-w_{n l}\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right\|_{C_{n-1}^{1}\left(H^{2}\left(Y_{n}\right) / \mathbb{R}\right)} \leq C\left\|A(\boldsymbol{z} ; x, \boldsymbol{y})-A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}
$$

From this and (4.5),

$$
\left\|A_{n-1}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-1}\right)-A_{n-1}\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}_{n-1}\right)\right\|_{\left(\boldsymbol{C}_{n-1}^{1}\right)^{d \times d}} \leq C\left\|A(\boldsymbol{z} ; x, \boldsymbol{y})-A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}} .
$$

Inductively, we then show that for all $i=1, \ldots, n$ and all $l=1, \ldots, n$,

$$
\left\|w_{i l}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{i-1}, \cdot\right)-w_{i l}\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}_{i-1}, \cdot\right)\right\|_{\boldsymbol{C}_{i-1}^{1}\left(H^{2}\left(Y_{i}\right) / \mathbb{R}\right)} \leq C\left\|A(\boldsymbol{z} ; x, \boldsymbol{y})-A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}} .
$$

Therefore for the homogenized coefficient $A_{0}(\boldsymbol{z} ; x)$ holds

$$
\left\|A_{0}(\boldsymbol{z} ; x)-A_{0}\left(\boldsymbol{z}^{\prime} ; x\right)\right\|_{C^{1}(\bar{D})^{d \times d}} \leq C\left\|A(\boldsymbol{z} ; x, \boldsymbol{y})-A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}} .
$$

From (4.7), we obtain for all $\boldsymbol{z}, \boldsymbol{z}^{\prime} \in U$ and every $\phi \in H_{0}^{1}(D)$ that

$$
\int_{D} A_{0}(\boldsymbol{z} ; x) \nabla\left(u_{0}(\boldsymbol{z} ; x)-u_{0}\left(\boldsymbol{z}^{\prime} ; x\right)\right) \cdot \nabla \phi(x) d x=\int_{D}\left(A_{0}\left(\boldsymbol{z}^{\prime} ; x\right)-A_{0}(\boldsymbol{z} ; x)\right) \nabla u_{0}\left(\boldsymbol{z}^{\prime} ; x\right) \cdot \nabla \phi(x) d x
$$

From this identity and from the assumed $H^{2}(D)$ regularity for the Dirichlet problem in $D$ we conclude that there exists a constant $C>0$ such that for all $\boldsymbol{z}, \boldsymbol{z}^{\prime} \in U$ it holds

$$
\begin{align*}
\left\|u_{0}(\boldsymbol{z} ; \cdot)-u_{0}\left(\boldsymbol{z}^{\prime} ; \cdot\right)\right\|_{H^{2}(D)} & \leq C\left\|\nabla\left(\left(A_{0}\left(\boldsymbol{z}^{\prime} ; x\right)-A_{0}(\boldsymbol{z} ; x)\right) \nabla u_{0}\left(\boldsymbol{z}^{\prime} ; x\right)\right)\right\|_{L^{2}(D)} \\
& \leq C \sup _{x \in D}\left(\left|A_{0}(\boldsymbol{z} ; x)-A_{0}\left(\boldsymbol{z}^{\prime} ; x\right)\right|+\left|\nabla A_{0}(\boldsymbol{z} ; x)-\nabla A_{0}\left(\boldsymbol{z}^{\prime} ; x\right)\right|\right) \\
& \leq C\left\|A(\boldsymbol{z} ; \cdot, \cdot)-A\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}} \tag{4.17}
\end{align*}
$$

From (4.9) and the uniform boundedness of $w_{i}$ in $\boldsymbol{C}_{i-1}^{1}\left(H^{2}\left(Y_{i}\right)\right)$, we get (4.15). A similar argument as in the proof of Proposition 1.8 shows that $\boldsymbol{u}$ as a map from $U$ to $\mathcal{H}$ is measurable.

From Proposition 4.5, we deduce that $\boldsymbol{u} \in L^{2}(U, \rho ; \boldsymbol{\mathcal { H }})$, so the coefficients $\boldsymbol{u}_{\nu}$ in the expansion (2.2) are all in $\mathcal{H}$.

### 4.2 Regularity of the complex parametric, deterministic problems (3.4)

We show that the solution $\boldsymbol{u}(\boldsymbol{\zeta}, \cdot, \cdot)$ of the problem (3.4) belongs to $\mathcal{H}$ when the complex parameter $\boldsymbol{\zeta}$ is in a subset $\overline{\mathcal{U}}$ of the domain $\mathcal{U}$ defined in (3.2). We choose a constant $\bar{K}<1$ that satisfies

$$
\begin{equation*}
\bar{K} \sum_{j=1}^{\infty}\left(\beta_{j}+\left\|\Psi_{j}\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}\right)<\frac{\alpha}{2 M} . \tag{4.18}
\end{equation*}
$$

We then choose a constant $\bar{J}_{0}$ so that

$$
\begin{equation*}
\sum_{j>\bar{J}_{0}}\left(\beta_{j}+\left\|\Psi_{j}\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}\right)<\frac{\alpha \bar{K}}{6 M(1+\bar{K})} \tag{4.19}
\end{equation*}
$$

We then denote $\bar{E}=\left\{1,2, \ldots, \bar{J}_{0}\right\}, \bar{F}=\mathbb{N} \backslash \bar{E}$ and set

$$
\left|\nu_{\bar{F}}\right|=\sum_{j>\bar{J}_{0}}\left|\nu_{j}\right| .
$$

For each index $\nu \in \mathcal{F}$, we define

$$
\begin{equation*}
\bar{r}_{m}=\bar{K} \text { when } m \leq \bar{J}_{0}, \text { and } \bar{r}_{m}=\max \left\{1, \frac{\alpha \nu_{m}}{M\left|\nu_{\bar{F}}\right|\left(\beta_{m}+\left\|\Psi_{m}\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}\right)}\right\} \text { when } m>\bar{J}_{0}, \tag{4.20}
\end{equation*}
$$

where we again adopted the convention that $\left|\nu_{m}\right| /\left|\nu_{\bar{F}}\right|=0$ if $\left|\nu_{\bar{F}}\right|=0$. For $m \geq 1$, we define the set $\overline{\mathcal{U}}_{m} \subset \mathbb{C}$ as

$$
\begin{equation*}
[-1,1] \subset \overline{\mathcal{U}}_{m}:=\left\{\boldsymbol{\zeta}_{m} \in \mathbb{C}: \operatorname{dist}\left(\zeta_{m},[-1,1]\right) \leq \bar{r}_{m}\right\} \subset \mathbb{C} . \tag{4.21}
\end{equation*}
$$

We then consider the complex parametric domain $\overline{\mathcal{U}} \subset \mathcal{U}$ defined as

$$
\overline{\mathcal{U}}=\bigotimes_{m=1}^{\infty} \overline{\mathcal{U}}_{m} \subset \mathbb{C}^{\mathbb{N}}
$$

We consider the problem (3.4) for complex valued parameter vectors $\boldsymbol{\zeta} \in \overline{\mathcal{U}}$. For $\boldsymbol{\zeta} \in \overline{\mathcal{U}}$, we have

$$
\begin{align*}
&\|A(\boldsymbol{\zeta} ; x, \boldsymbol{y})\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}} \leq\|\bar{A}(x, \boldsymbol{y})\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}+\sum_{m=1}^{\infty}\left\|\Psi_{m}(x, \boldsymbol{y})\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}\left(1+\bar{r}_{m}\right) \\
& \leq\|\bar{A}(x, \boldsymbol{y})\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}+\sum_{m=1}^{\bar{J}_{0}}\left\|\Psi_{m}\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}(1+\bar{K}) \\
&+\sum_{m>\bar{J}_{0}}\left(2+\frac{\alpha \nu_{m}}{M\left|\nu_{\bar{F}}\right|\left(\beta_{m}+\left\|\Psi_{m}\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}\right)}\right)\left\|\Psi_{m}\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}} \\
& \leq\|\bar{A}(x, \boldsymbol{y})\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}+2 \sum_{m=1}^{\infty}\left\|\Psi_{m}\right\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}+\frac{\alpha}{M} . \tag{4.22}
\end{align*}
$$

As in (3.3), we have

$$
\begin{equation*}
\left|A_{i j}(\zeta ; x, y)\right| \leq\left\|\bar{A}_{i j}\right\|_{L^{\infty}(D \times Y)}+4 \frac{\kappa}{1+\kappa} \alpha_{0}+\frac{2 \alpha}{M} . \tag{4.23}
\end{equation*}
$$

Therefore $A(\boldsymbol{\zeta} ; x, \boldsymbol{y})$ is uniformly bounded in $\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}$ for all $\boldsymbol{\zeta} \in \overline{\mathcal{U}}$. We show next that the solution of the parametric problem is jointly holomorphic with respect to any finite set of parameters. For each index $\nu \in \mathcal{F}$, we define the (finite dimensional) domain

$$
\overline{\mathcal{U}}_{\nu}=\bigotimes_{j \in \operatorname{supp}(\nu)} \overline{\mathcal{U}}_{j}
$$

We have the following analyticity result.
Proposition 4.6 For $\nu \in \mathcal{F}$ and $\zeta \in \overline{\mathcal{U}}$, fixing $\zeta_{k}$ for $k \notin \operatorname{supp}(\nu)$, under Assumption 4.3, if the domain $D$ is convex then $\boldsymbol{u}$ is analytic as a map from $\overline{\mathcal{U}} \rightarrow \mathcal{H}$ when the constant $M$ in (4.18) and (4.19) is sufficiently large.

Proof Let $w_{n l}(\boldsymbol{\zeta})$ be the solution of problem (4.3) for the complex valued coefficient $A(\boldsymbol{\zeta} ; x, \boldsymbol{y})$. We show that $w_{n l}(\boldsymbol{\zeta})$ is holomorphic as a mapping from $\overline{\mathcal{U}}_{\nu}$ to $\boldsymbol{C}_{n-1}^{1}\left(H^{2}\left(Y_{n}\right) / \mathbb{R}\right)$. To this end, we establish complex differentiability by showing that certain difference quotient have limits.

For any $m$, we fix all coordinates $\zeta_{k}$ for $k \neq m$, and partition $\zeta \in \mathbb{C}^{\mathbb{N}}$ as $\zeta=\left(\zeta_{m}^{*}, \zeta_{m}\right)$. Let further $\delta \in \mathbb{C}$ denote the step size of the difference quotients

$$
\eta_{m n l}^{\delta}(\boldsymbol{\zeta} ; \cdot, \cdot):=\delta^{-1}\left(w_{n l}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; \cdot, \cdot\right)-w_{n l}(\boldsymbol{\zeta} ; \cdot \cdot \cdot)\right)
$$

The function $\eta_{m n l}^{\delta}$ satisfies

$$
\begin{aligned}
& \int_{Y_{n}} A(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \nabla_{y_{n}} \eta_{m n l}^{\delta}(\boldsymbol{\zeta} ; x, y) \cdot \overline{\nabla_{y_{n}} \phi_{n}\left(y_{n}\right)} d y_{n}=-\int_{Y_{n}} \Psi_{m} e_{l} \cdot \overline{\nabla_{y_{n}} \phi_{n}\left(y_{n}\right)} d y_{n} \\
&-\int_{Y_{n}} \Psi_{m} \nabla_{y_{n}} w_{n l}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}\right) \cdot \overline{\nabla_{y_{n}} \phi_{n}\left(y_{n}\right)} d y_{n}, \quad \forall \phi_{n} \in H_{\#}^{1}\left(Y_{n}\right)
\end{aligned}
$$

Let $\eta_{m n l}(\boldsymbol{\zeta} ; x, \boldsymbol{y})$ denote the solution of the problem

$$
\begin{aligned}
& \int_{Y_{n}} A(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \nabla_{y_{n}} \eta_{m n l}(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \cdot \overline{\nabla_{y_{n}} \phi_{n}\left(y_{n}\right)} d y_{n}=-\int_{Y_{n}} \Psi_{m} e_{l} \cdot \overline{\nabla_{y_{n}} \phi_{n}\left(y_{n}\right)} d y_{n} \\
&-\int_{Y_{n}} \Psi_{m} \nabla_{y_{n}} w_{n l}(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \cdot \overline{\nabla_{y_{n}} \phi_{n}\left(y_{n}\right)} d y_{n}, \quad \forall \phi_{n} \in H_{\#}^{1}\left(Y_{n}\right)
\end{aligned}
$$

We then have, for every $\zeta \in U$,
$\int_{Y_{n}} A(\boldsymbol{\zeta} ; x, \boldsymbol{y}) \nabla_{y_{n}}\left(\eta_{m n l}^{\delta}-\eta_{m n l}\right) \cdot \overline{\nabla_{y_{n}} \phi_{n}} d y_{n}=-\int_{Y_{n}} \Psi_{n} \nabla_{y_{n}}\left(w_{n l}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}\right)-w_{n l}(\boldsymbol{\zeta} ; x, \boldsymbol{y})\right) \cdot \overline{\nabla_{y_{n}} \phi_{n}} d y_{n}$.
Proceeding in the same fashion as in the proof of Proposition 4.2 we obtain

$$
\left\|\eta_{m n l}^{\delta}-\eta_{m n l}\right\|_{C_{n-1}^{1}\left(H^{2}\left(Y_{n}\right) / \mathbb{R}\right)} \leq c \| \nabla_{y_{n}} \cdot\left(\Psi_{n} \nabla_{y_{n}}\left(w_{n l}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; \cdot, \cdot\right)-w_{n l}\left(\boldsymbol{\zeta}^{\prime} ; \cdot, \cdot\right)\right) \|_{\boldsymbol{C}_{n-1}^{1}\left(L^{2}\left(Y_{n}\right)\right)}\right.
$$

which converges to 0 when $\delta \rightarrow 0$ as $w_{n l}$ is continuous as a map from $\overline{\mathcal{U}}_{m}$ to $C_{n-1}^{1}\left(H^{2}\left(Y_{n}\right) / \mathbb{R}\right)$ (which can be shown as in the proof of Proposition 4.5). Therefore $w_{n l}$ is complex differentiable with respect to $\zeta_{m}$ and therefore an analytic function of $\zeta_{m}$ taking values in $\boldsymbol{C}_{n-1}^{1}\left(H^{2}\left(Y_{n}\right) / \mathbb{R}\right)$. From Hartogs' theorem, we conclude that $w_{n l}$ is analytic as a function from $\overline{\mathcal{U}}_{\nu}$ to $\boldsymbol{C}_{n-1}^{1}\left(H^{2}\left(Y_{n}\right) / \mathbb{R}\right)$. By (4.5), $A_{n-1}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)$ is an analytic, $\boldsymbol{C}_{n-1}^{1}$-valued function of $\boldsymbol{\zeta}$ in $\overline{\mathcal{U}}_{\nu}$.

Next we consider $w_{(n-1) l}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)$. Again, we verify analyticity by showing complex differentiability via the difference quotients

$$
\eta_{m(n-1) l}^{\delta}=\frac{w_{(n-1) l}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; \cdot, \cdot\right)-w_{(n-1) l}(\boldsymbol{\zeta} ; \cdot, \cdot)}{\delta}
$$

For these difference quotients, we have for parameter vectors $\boldsymbol{\zeta}$ as above the equation

$$
\begin{aligned}
& \int_{\mathbf{Y}_{n-1}} A_{n-1}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right) \nabla_{y_{n-1}} \eta_{m(n-1) l}^{\delta}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right) \cdot \overline{\nabla_{y_{n-1}} \phi_{n-1}\left(y_{n-1}\right)} d y_{n-1} \\
& =-\int_{\mathbf{Y}_{n-1}} \frac{A_{n-1}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}_{n-1}\right)-A_{n-1}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)}{\delta} e_{l} \cdot \overline{\nabla_{y_{n-1}} \phi_{n-1}\left(y_{n-1}\right)} d \boldsymbol{y}_{n-1} \\
& -\int_{\mathbf{Y}_{n-1}} \frac{A_{n-1}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}_{n-1}\right)-A_{n-1}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)}{\delta} \nabla_{y_{n-1}} w_{(n-1) l}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}_{n-1}\right) \\
& \quad \bar{\nabla}_{y_{n-1} \phi_{n-1}\left(y_{n-1}\right)} d y_{n-1},
\end{aligned}
$$

forall $\phi_{n-1} \in H_{\#}^{1}\left(Y_{n-1}\right)$. We next let $\eta_{m(n-1) l}\left(\boldsymbol{z} ; x, \boldsymbol{y}_{n-2}, y_{n-1}\right) \in H_{\#}^{1}\left(Y_{n-1}\right) / \mathbb{R}$ satisfy

$$
\begin{aligned}
& \int_{Y_{n-1}} A_{n-1}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right) \nabla_{y_{n-1}} \eta_{m(n-1) l}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right) \cdot \overline{\nabla_{y_{n-1}} \phi_{n-1}\left(y_{n-1}\right)} d y_{n-1} \\
& =-\int_{Y_{n-1}} \frac{\partial A_{n-1}}{\partial \zeta_{m}} e_{l} \cdot \overline{\nabla_{y_{n-1}} \phi_{n-1}\left(y_{n-1}\right)} d y_{n-1} \\
& \quad-\int_{Y_{n-1}} \frac{\partial A_{n-1}}{\partial \zeta_{m}} \nabla_{y_{n-1}} w_{(n-1) l}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right) \cdot \overline{\nabla_{y_{n-1}} \phi_{n-1}\left(y_{n-1}\right)} d y_{n-1}
\end{aligned}
$$

We deduce

$$
\begin{aligned}
& \int_{Y_{n-1}} A_{n-1}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right) \nabla_{y_{n-1}}\left(\eta_{m(n-1) l}^{\delta}(\boldsymbol{\zeta} ; x, \boldsymbol{y})-\eta_{m(n-1) l}(\boldsymbol{\zeta} ; x, \boldsymbol{y})\right) \cdot \overline{\nabla_{y_{n-1}} \phi_{n-1}\left(y_{n-1}\right)} d y_{n-1} \\
& =-\int_{Y_{n-1}}\left(\frac{A_{n-1}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}_{n-1}\right)-A_{n-1}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)}{\delta}-\frac{\partial A_{n-1}}{\partial \zeta_{m}}\right) e_{l} \cdot \overline{\nabla_{y_{n-1}} \phi_{n-1}\left(y_{n-1}\right)} d y_{n-1} \\
& -\int_{Y_{n-1}}\left(\frac{A_{n-1}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}_{n-1}\right)-A_{n-1}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)}{\delta}-\right. \\
& \left.\quad \frac{\partial A_{n-1}}{\partial \zeta_{m}}\right) \nabla_{y_{n-1}} w_{(n-1) l}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}\right) \cdot \overline{\nabla_{y_{n-1}} \phi_{n-1}\left(y_{n-1}\right)} d y_{n-1} \\
& \quad-\int_{Y_{n-1}} \frac{\partial A_{n-1}}{\partial \zeta_{m}}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right) \nabla_{y_{n-1}}\left(w_{(n-1) l}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}_{n-1}\right)-\left(w_{(n-1) l}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)\right)\right) \\
& \quad \cdot \nabla_{y_{n-1}} \phi_{n-1}\left(y_{n-1}\right) \\
& d y_{n-1}
\end{aligned}
$$

Therefore

$$
\begin{aligned}
& \left\|\eta_{m(n-1) l}^{\delta}(\boldsymbol{\zeta} ; x, \boldsymbol{y})-\eta_{m(n-1) l}(\boldsymbol{\zeta} ; x, \boldsymbol{y})\right\|_{\boldsymbol{C}_{n-2}^{1}\left(H^{2}\left(Y_{n-1}\right) / \mathbb{R}\right)} \leq \\
& \| \nabla_{y_{n-1}} \cdot\left[\left(\frac{A_{n-1}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}_{n-1}\right)-A_{n-1}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)}{\delta}-\frac{\partial A_{n-1}}{\partial \zeta_{m}}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)\right) e_{l}+\right. \\
& \left(\frac{A_{n-1}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}_{n-1}\right)-A_{n-1}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)}{\delta}-\frac{\partial A_{n-1}}{\partial \zeta_{m}}\right) \nabla_{y_{n-1}} w_{(n-1) l}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}_{n-1}\right) \\
& \left.-\frac{\partial A_{n-1}}{\partial \zeta_{m}}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right) \nabla_{y_{n-1}}\left(w_{(n-1) l}\left(\zeta_{m}^{*}, \zeta_{m}+\delta ; x, \boldsymbol{y}_{n-1}\right)-\left(w_{(n-1) l}\left(\boldsymbol{\zeta} ; x, \boldsymbol{y}_{n-1}\right)\right)\right)\right] \|_{\boldsymbol{C}_{n-2}^{1}\left(L^{2}\left(Y_{n-1}\right)\right)}
\end{aligned}
$$

which converges to 0 as $A_{n-1}$ is holomorphic as a mapping from $\overline{\mathcal{U}}_{m}$ to $C_{n-1}^{1}$ and $w_{(n-1) l}$ is continuous as a mapping from $\overline{\mathcal{U}}_{m}$ to $\boldsymbol{C}_{n-1}^{1}\left(H^{2}\left(Y_{n-1}\right) / \mathbb{R}\right)$.

Similarly, $w_{i l}$ is analytic as a map from $\overline{\mathcal{U}}_{\nu}$ to $C_{i-1}^{2}\left(H^{2}\left(Y_{i}\right)\right)$ for other values of $i$. To show that $u_{i}$ is analytic from $\overline{\mathcal{U}}_{\nu}$ to $\mathcal{H}_{i}$ and $\boldsymbol{u}$ is analytic from $\overline{\mathcal{U}}_{\nu}$ to $\mathcal{H}$, it remains to establish the analyticity of $u_{0}$ as a map from $\overline{\mathcal{U}}$ to $H^{2}(D)$ where the domain $D$ is convex. We note that Theorem 3.2.1.2 of Grisvard [9] is not readily applicable to elliptic equations with complex coefficients in a convex domain.

As $w_{i l}$ are holomorphic as a map from $\overline{\mathcal{U}}_{\nu}$ to $C_{i-1}^{2}\left(H^{2}\left(Y_{i}\right)\right)$, the coefficient $A_{0}(\zeta ; x)$ of the complex parametric homogenized equation is analytic. As $\overline{\mathcal{U}}^{i-1} \subset \mathcal{U}$, from (3.5), $\Re\left(\xi^{H} A(\zeta ; x, \boldsymbol{y}) \xi\right) \geq \alpha|\xi|^{2} / 2$ for all $\zeta \in \overline{\mathcal{U}}, x \in D$ and $\boldsymbol{y} \in \mathbf{Y}$. Following the proof of Proposition 4.1, we deduce that

$$
\forall \xi \in \mathbb{C}^{d}, \zeta \in \bar{U}, x \in D: \quad \Re\left(\xi^{H} A_{0}(\zeta ; x) \xi\right) \geq \frac{\alpha}{2}|\xi|^{2}
$$

which implies that for all $\xi \in \mathbb{R}^{d}$,

$$
\xi^{\top} \Re A_{0}(\zeta ; x) \xi \geq \frac{\alpha}{2}|\xi|^{2}
$$

Further, there is a positive constant $\beta^{\prime \prime}$ that depends only on $\alpha, \sup _{i, j}\left\|A_{i j}(\zeta ; x)\right\|_{L^{\infty}(D)}, d$ and $n$ such that

$$
\xi^{\top} \Re A_{0}(\zeta ; x) \xi \leq \beta^{\prime \prime}|\xi|^{2}
$$

From (3.3). $\beta^{\prime \prime}$ can be chosen independently of $M$ (i.e. independently of the complex parametric domain $\overline{\mathcal{U}}$ when $M$ is sufficiently large; here we choose $M \geq 4)$. Let $D_{m}$ be a sequence of convex subdomains of $D$ with smooth boundary such that $\operatorname{dist}\left(\partial D_{m}, \partial D\right) \rightarrow 0$ as in the proof of Theorem 3.2.1.2 in [9]. Consider the Dirichet problems

$$
-\nabla \cdot\left(A_{0}(\zeta ; x) \nabla \phi_{m}(\zeta ; x)\right)=f(x), \quad \phi_{m} \in H_{0}^{1}\left(D_{m}\right)
$$

i.e.

$$
-\nabla \cdot\left(\left(\Re A_{0}(\zeta ; x)\right) \nabla \phi_{m}\right)=f(x)+i \nabla \cdot\left(\left(\Im A_{0}(\zeta ; x)\right) \nabla \phi_{m}\right) .
$$

As the boundary of $D_{m}$ is smooth, $\phi_{m} \in H^{2}\left(D_{m}\right)$ ([19], Section 20). Applying the proof of Lemma 3.1.3.2 in [9] for $\Re \phi_{m}$ and $\Im \phi_{m}$ respectively, we find that there is a constant $c_{1}$ which depends on $\left\|\Re A_{0}(z ; x)\right\|_{\left(C^{1}(\bar{D})\right)^{d \times d}}$, the diameter of $D, \alpha$ and $\beta^{\prime \prime}$ such that

$$
\begin{equation*}
\left\|\phi_{m}\right\|_{H^{2}\left(D_{m}\right)} \leq c_{1}\left(\left\|\nabla \cdot\left(A_{0}(\zeta ; x) \nabla \phi_{m}(\zeta ; x)\right)\right\|_{L^{2}\left(D_{m}\right)}+\left\|\nabla \cdot\left(\Im A_{0}(\zeta ; x) \nabla \phi_{m}\right)\right\|_{L^{2}\left(D_{m}\right)}\right) \tag{4.24}
\end{equation*}
$$

We note that

$$
\left\|\nabla \phi_{m}\right\|_{L^{2}\left(D_{m}\right)} \leq \frac{2}{\alpha}\left\|\nabla \cdot\left(A_{0}(\zeta ; x) \nabla \phi_{m}(\zeta ; x)\right)\right\|_{L^{2}\left(D_{m}\right)}
$$

Further, from (4.6), $\left\|A_{0}(z ; x)\right\|_{\left(C^{1}(\bar{D})\right)^{d \times d}}$ has an upper bound depending on an upper bound of $\|A(\boldsymbol{\zeta} ; x, \boldsymbol{y})\|_{\left(\boldsymbol{C}_{n}^{1}\right)^{d \times d}}$, which can be chosen independently of $M$ (from (4.22)). Therefore

$$
\left\|\phi_{m}\right\|_{H^{2}\left(D_{m}\right)} \leq c_{1}\left(c_{2}\left\|\nabla \cdot\left(A_{0}(\zeta ; x) \nabla \phi_{m}(\zeta ; x)\right)\right\|_{L^{2}\left(D_{m}\right)}+c_{3} \sup _{i, j}\left\|\Im A_{0 i j}(\zeta ; x)\right\|_{L^{\infty}\left(D_{m}\right)}\left\|\nabla \nabla \phi_{m}\right\|_{L^{2}\left(D_{m}\right)^{d \times d}}\right)
$$

where the constants $c_{1}$ and $c_{2}$ are independent of $M$, and $c_{3}$ only depends on the dimension $d$. Assume that $\sup _{i, j}\left\|\Im A_{0 i j}(\zeta ; x)\right\|_{L^{\infty}(D)}$ is sufficiently small so that

$$
\begin{equation*}
c_{3} \sup _{i, j}\left\|\Im A_{0 i j}(\zeta ; x)\right\|_{L^{\infty}(D)}<1 /\left(2 c_{1}\right) \tag{4.25}
\end{equation*}
$$

we then have

$$
\left\|\phi_{m}\right\|_{H^{2}\left(D_{m}\right)} \leq 2 c_{1} c_{2}\left\|\nabla \cdot\left(A_{0}(\zeta ; x) \nabla \phi_{m}(\zeta ; x)\right)\right\|_{L^{2}\left(D_{m}\right)}
$$

Therefore, $\phi_{m}$ is uniformly bounded in $H_{0}^{1}\left(D_{m}\right) \cap H^{2}\left(D_{m}\right)$. Arguing as in the proof of Theorem 3.2.1.2 in Grisvard [9], $\phi_{m}$ (extended to 0 outside $D_{m}$ ) converges weakly to $u_{0}$ in $H_{0}^{1}(D)$; and the weak limit in $L^{2}(D)$ of the second derivatives of $\phi_{m}$ (again extended to 0 outside $D_{m}$ ) must be the second derivative of $u_{0}$; thus $u_{0} \in H_{0}^{1}(D) \cap H^{2}(D)$, and

$$
\left\|u_{0}\right\|_{H^{2}(D)} \leq c\left\|\nabla \cdot\left(A_{0}(\zeta ; x) \nabla u_{0}(\zeta ; x)\right)\right\|_{L^{2}(D)}
$$

It remains to show that we can find the a constant $M$ in (4.18) and (4.19) so that (4.25) holds. We note that

$$
\begin{aligned}
\left|\Im A_{i j}\right| & \leq \sum_{m=1}^{\infty} \bar{r}_{m}\left\|\Psi_{m i j}\right\|_{L^{\infty}(D)} \leq \sum_{m=1}^{J_{0}} \bar{K} 2 \beta_{m}+\sum_{m>J_{0}}\left(1+\frac{\alpha \nu_{m}}{M\left|\nu_{\bar{F}}\right|\left(\beta_{m}+\left\|\psi_{m}\right\|_{\left.\left(\mathbf{C}_{n}^{1}\right)^{d \times d}\right)}\right.}\right) 2 \beta_{m} \\
& \leq \frac{\alpha}{M}+\frac{\alpha}{3 M}+\frac{2 \alpha}{M}
\end{aligned}
$$

which is small when $M$ is large.
Using the cell problem (4.10) for the complex parametric problem, with $\phi_{n}=\Im w_{n l}$, taking the imaginary part of both sides, we have

$$
\begin{aligned}
\int_{Y_{n}} \Re A(\zeta ; x, y) \nabla_{y_{n}} \Im w_{n l}(\zeta ; x, y) \cdot \nabla_{y_{n}} \Im w_{n l}(\zeta ; x, y)= & -\int_{Y_{n}} \Im A(\zeta ; x, y) e_{l} \cdot \nabla \Im w_{n l} d y_{n} \\
& -\int_{Y_{n}} \Im A(\zeta ; x, y) \nabla \Re w_{n l} \cdot \nabla \Im w_{n l} d y_{n}
\end{aligned}
$$

Therefore,

$$
\left\|\nabla \Im w_{n l}(\zeta ; x, \cdot)\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}, L^{2}\left(Y_{n}\right)\right)^{d}} \leq c(\alpha, d) \sup _{i, j}\left\|\Im A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\left(1+\left\|\nabla \Re w_{n l}\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}, L^{2}\left(Y_{n}\right)\right)^{d}}\right) .
$$

We note from (4.10) that

$$
\left\|\nabla \Re w_{n l}\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}, L^{2}\left(Y_{n}\right)\right)^{d}} \leq c(\alpha, d)\left(1+\sup _{i, j}\left\|A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right),
$$

so

$$
\left\|\nabla \Im w_{n l}(\zeta ; x, \cdot)\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}, L^{2}\left(Y_{n}\right)\right)^{d}} \leq c(\alpha, d) \sup _{i, j}\left\|\Im A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\left(1+\sup _{i, j}\left\|A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right)
$$

From (4.5), we have

$$
\begin{aligned}
\sup _{i, j}\left\|\Im A_{(n-1) i j}\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}\right)} \leq & \sup _{i, j}\left\|\Im A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}+c(d)\left(\sup _{i, j}\left\|A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\left\|\Im \nabla w_{n}\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}, L^{2}\left(Y_{n}\right)\right)^{d \times d}}\right. \\
& +\sup _{i, j}\left\|\Im A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\left\|\Re \nabla w_{n}\right\|_{\left.L^{\infty}\left(D \times \mathbf{Y}_{n-1}, L^{2}\left(Y_{n}\right)\right)^{d \times d}\right)} \\
\leq & c(\alpha, d) \sup _{i, j}\left\|\Im A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\left(1+\sup _{i, j}\left\|A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right)^{2} .
\end{aligned}
$$

Repeating this argument we have

$$
\begin{gathered}
\sup _{i, j}\left\|\Im A_{0 i j}\right\|_{L^{\infty}(D)} \leq c(\alpha, d)^{n} \sup _{i, j}\left\|\Im A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\left(1+\sup _{i, j}\left\|A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right)^{2} \cdot\left(1+\sup _{i, j}\left\|A_{(n-1) i j}\right\|_{L^{\infty}\left(D \times \mathbf{Y}_{n-1}\right)}\right)^{2} \\
\ldots\left(1+\sup _{i, j}\left\|A_{1 i j}\right\|_{L^{\infty}(D \times Y)}\right)^{2} \\
\leq \sup _{i, j}\left\|\Im A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\left[c(\alpha, d)\left(1+\sup _{i, j}\left\|A_{i j}\right\|_{L^{\infty}(D \times \mathbf{Y})}\right)\right]^{2^{n+1}-2},
\end{gathered}
$$

where the last estimate is obtained in a similar fashion as for (4.8). Thus, when the constant $M$ in (4.18) and (4.19) is sufficiently large, the condition (4.25) holds.

### 4.3 Summability of $\boldsymbol{u}_{\nu}$

We now study the summability of the $\mathcal{H}$ norms of $\boldsymbol{u}_{\nu}$. First, we have the following estimate
Proposition 4.7 The Legendre coefficients $\boldsymbol{u}_{\nu}$ in (3.7) of the parametric, deterministic solution ( $\left.\boldsymbol{z} ; x, \boldsymbol{y}\right)$ satisfy the estimate

$$
\forall \nu \in \mathcal{F}: \quad\left\|\boldsymbol{u}_{\nu}\right\|_{\mathcal{H}} \leq C\left(\prod_{m \in \operatorname{supp}(\nu)} \frac{2(1+\bar{K})^{\nu_{2}}}{\bar{K}} \bar{\eta}_{m}^{\nu_{m}}\right),
$$

where $\bar{\eta}_{m}:=\bar{r}_{m}+\sqrt{1+\bar{r}_{m}^{2}}$ with $\bar{r}_{m}$ defined in (4.20).
The proof of this Proposition is identical to that for Proposition 3.6.
To study the summability of the sequence $\left(\left\|u_{\nu}\right\|_{\mathcal{H}}\right)_{\nu \in \mathcal{F}}$, we make the following
Assumption 4.8 There is a constant $0<p<1$ such that

$$
\sum_{k=1}^{\infty}\left\|\Psi_{k}\right\|_{\left(C_{n}^{1}\right) d \times d}^{p}<\infty .
$$

Remark 4.9 Assumption 4.8 holds when in estimate (4.14),

$$
p\left(\frac{t}{d}-\varepsilon\right)\left(1-2 \frac{t^{*}}{t}\right)>1 .
$$

We note that if $\beta_{k}$ is taken as an upper bound for $\|$ trace $\Psi_{k} \|_{L^{\infty}(D \times \mathbf{Y})}$, then Assumption 4.8 implies Assumption 3.2.
Proposition 4.10 Under Assumption 4.8, $\left(\left\|\boldsymbol{u}_{\nu}\right\|_{\mathcal{H}}\right)_{\nu \in \mathcal{F}} \in \ell^{p}(\mathcal{F})$.
The proof of this Proposition is identical to that of Proposition 3.8 except that we use $\bar{K}, \bar{E}, \bar{F}$, in places of $K, E$ and $F$.

Remark 4.11 All of the above results hold if the domain $D$ is not convex but has a smooth boundary.

## 5 Correctors

### 5.1 Correctors for two scale problems

For two scale problems where the coefficient $A$ does not depend on the slow variable $x$, an estimate of the solution $u^{\varepsilon}$ in terms of the solution $u_{0}$ and the corrector $u_{1}$ of the homogenized, high dimensional one-scale problem (4.7) has been established under the provision of sufficient regularity. Specifically, assuming that $u_{0} \in C^{2}(\bar{D})$ and $w_{1 l} \in W^{1, \infty}(Y)$ (see e.g. Jikov et al. [12] page 28), we will now prove this result, under slightly weaker regularity requirements for $u_{0}$ than what was required in [12]. We give its full proof here to verify the regularity requirements and, more importantly, to show that the error estimate for the two scale parametric problem (1.16) holds uniformly for all $\boldsymbol{z} \in U$. As for two length scales there is only one fast variable, we denote in this case $\boldsymbol{y}$ by $y$ and $\mathbf{Y}$ simply by $Y$. For two scale problems we denote by $w^{l}(\boldsymbol{z} ; x, y)$ the functions $w_{1 l}(\boldsymbol{z} ; x, y)$.
Proposition 5.1 For the parametric two scale problem (1.16), assume that $A(z ; x, y) \in L^{\infty}\left(U ; C^{1}\left(\bar{D} ; C_{\#}^{1}(Y)\right)\right)$, that the function $u_{0}(\boldsymbol{z} ; x) \in L^{\infty}\left(U ; H^{2}(D)\right), w^{l}(\boldsymbol{z} ; x, \boldsymbol{y}) \in L^{\infty}\left(U ; C^{1}\left(\bar{D} ; H_{\#}^{2}(Y)\right)\right) \cap L^{\infty}\left(U \times \bar{D} ; C_{\#}^{1}(\bar{Y})\right)$, and that the domain $D$ has a Lipschitz boundary. Then there exists a constant $c>0$ such that for every $0<\varepsilon \leq 1$,

$$
\begin{equation*}
\sup _{z \in U}\left\|u^{\varepsilon}(\boldsymbol{z} ; x)-\left[u_{0}(\boldsymbol{z} ; x)+\varepsilon u_{1}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)\right]\right\|_{H^{1}(D)} \leq c \varepsilon^{1 / 2} . \tag{5.1}
\end{equation*}
$$

Proof For $\boldsymbol{z} \in U$, define

$$
u_{1}^{\varepsilon}(\boldsymbol{z} ; x)=u_{0}(\boldsymbol{z} ; x)+\varepsilon w^{l}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial u_{0}(\boldsymbol{z} ; x)}{\partial x_{l}} .
$$

We first show that

$$
\left\|\operatorname{div} A^{\varepsilon} \nabla u_{1}^{\varepsilon}-\operatorname{div} A_{0} \nabla u_{0}\right\|_{H^{-1}(D)} \leq c \varepsilon,
$$

where $c$ is independent of $z$. We adapt the argument of [12] page 28 for the case where $u_{0} \in H^{2}(D)$ (but not in $\left.C^{2}(\bar{D})\right)$. We note that

$$
\begin{aligned}
& \left(A^{\varepsilon}(\boldsymbol{z} ; x) \nabla u_{1}^{\varepsilon}(\boldsymbol{z} ; x)\right)_{i} \\
= & \left(A_{i j}^{\varepsilon}(\boldsymbol{z} ; x)+A_{i k}^{\varepsilon}(\boldsymbol{z} ; x) \frac{\partial w^{j}}{\partial y_{k}}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)\right) \frac{\partial u_{0}}{\partial x_{j}}(\boldsymbol{z} ; x)+\varepsilon A_{i j}^{\varepsilon}(\boldsymbol{z} ; x) w^{k}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial^{2} u_{0}}{\partial x_{j} \partial x_{k}}(\boldsymbol{z} ; x) \\
= & A_{0 i j}(\boldsymbol{z} ; x) \frac{\partial u_{0}}{\partial x_{j}}(\boldsymbol{z} ; x)+\left(A_{i j}^{\varepsilon}(\boldsymbol{z} ; x)+A_{i k}^{\varepsilon}(\boldsymbol{z} ; x) \frac{\partial w^{j}}{\partial y_{k}}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)-A_{0 i j}(\boldsymbol{z} ; x)\right) \frac{\partial u_{0}}{\partial x_{j}}(\boldsymbol{z} ; x) \\
+ & \varepsilon A_{i j}^{\varepsilon}(\boldsymbol{z} ; x) w^{k}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial^{2} u_{0}}{\partial x_{j} \partial x_{k}}(\boldsymbol{z} ; x) \\
= & A_{0 i j}(\boldsymbol{z} ; x) \frac{\partial u_{0}}{\partial x_{j}}(\boldsymbol{z} ; x)+g_{i}^{j}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial u_{0}}{\partial x_{j}}(\boldsymbol{z} ; x)+\varepsilon A_{i j}^{\varepsilon}(\boldsymbol{z} ; x) w^{k}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial^{2} u_{0}}{\partial x_{j} \partial x_{k}}(\boldsymbol{z} ; x),
\end{aligned}
$$

where the functions $g_{i}^{j}(\boldsymbol{z} ; x, y)$ (which are $Y$-periodic with respect to $y$ ) are defined as

$$
g_{i}^{j}(\boldsymbol{z} ; x, y)=A_{i j}(\boldsymbol{z} ; x, y)+A_{i k}(\boldsymbol{z} ; x, y) \frac{\partial w^{j}}{\partial y_{k}}(\boldsymbol{z} ; x, y)-A_{0 i j}(\boldsymbol{z} ; x) .
$$

By definition, for all $\boldsymbol{z} \in U$ and every $x \in D$ holds

$$
\int_{Y} g_{i}^{j}(\boldsymbol{z} ; x, y) d y=0 \quad \text { and } \quad \frac{\partial}{\partial y_{i}} g_{i}^{j}(\boldsymbol{z} ; x, y)=0 .
$$

Therefore, there exist functions $\alpha_{i j}^{k}(\boldsymbol{z} ; x, y)$ which are $Y$-periodic in $y$ such that $\alpha_{i j}^{k}=-\alpha_{j i}^{k}$ and

$$
\begin{equation*}
g_{i}^{k}(\boldsymbol{z} ; x, y)=\frac{\partial}{\partial y_{j}} \alpha_{i j}^{k}(\boldsymbol{z} ; x, y) . \tag{5.2}
\end{equation*}
$$

As $w^{j} \in L^{\infty}\left(U ; C^{1}\left(\bar{D} ; H_{\#}^{2}(Y)\right)\right)$ and $A \in L^{\infty}\left(U ; C^{1}\left(\bar{D} ; C_{\#}^{1}(\bar{Y})\right)_{s y m}^{d \times d}\right), g_{i}^{j}(x, y) \in L^{\infty}\left(U ; C^{1}\left(\bar{D} ; H_{\#}^{1}(Y)\right)\right)$. The functions $\alpha_{i j}^{k}$ in (5.2) are constructed as follows (see Jikov et al [12] page 7). We write $\boldsymbol{g}^{k}=\left(g_{i}^{k}\right)$ as a Fourier series as

$$
\boldsymbol{g}^{k}(\boldsymbol{z} ; x, y)=\sum_{l \in \mathbb{Z}^{d}, l \neq 0} \boldsymbol{g}_{l}^{k}(\boldsymbol{z} ; x) \exp (\sqrt{-1} l \cdot y)
$$

As $\boldsymbol{g}^{k}(\boldsymbol{z} ; x, y) \in L^{\infty}\left(U ; C^{1}\left(\bar{D} ; H_{\#}^{1}(Y)\right)\right)^{d}$, we have $\boldsymbol{g}_{l}^{k} \in L^{\infty}\left(U ; C^{1}(\bar{D})\right)^{d}$ and, for all $r=1, \ldots, d$, there exists a constant $C_{r}$ such that

$$
\begin{equation*}
\sup _{\boldsymbol{z} \in U} \sup _{x \in D} \sum_{l \in \mathbb{Z}^{d}, l \neq 0}\left|\boldsymbol{g}_{l}^{k}(\boldsymbol{z}, x)\right|^{2} l_{r}^{2} \leq C_{r} \tag{5.3}
\end{equation*}
$$

The functions $\alpha_{i j}^{k}$ are defined as

$$
\alpha_{i j}^{k}(\boldsymbol{z} ; x, y)=\sqrt{-1} \sum_{l \in \mathbb{Z}^{d}, l \neq 0} \frac{\left(\boldsymbol{g}_{l}^{k}(\boldsymbol{z} ; x)\right)_{j} l_{i}-\left(\boldsymbol{g}_{l}^{k}(\boldsymbol{z} ; x)\right)_{i} l_{j}}{|l|^{2}} \exp (\sqrt{-1} l \cdot \boldsymbol{y})
$$

From this definition and (5.3), it is then obvious that for $r, s=1, \ldots, d$

$$
\sup _{\boldsymbol{z} \in U} \sup _{x \in D} \sum_{l \in \mathbb{Z}^{d}, l \neq 0} \frac{\left|\left(\boldsymbol{g}_{l}^{k}(\boldsymbol{z} ; x)\right)_{j} l_{i}-\left(\boldsymbol{g}_{l}^{k}(\boldsymbol{z} ; x)\right)_{i} l_{j}\right|^{2}}{|l|^{4}} l_{r}^{2} l_{s}^{2} \leq C_{r s}
$$

Therefore $\alpha_{i j}^{k}(\boldsymbol{z} ; x, y) \in L^{\infty}\left(U ; C^{1}\left(\bar{D} ; H_{\#}^{2}(Y)\right)\right.$ and, by the embedding theorem, for $d \leq 3$ holds $\alpha_{i j}^{k}(\boldsymbol{z} ; x, y) \in$ $L^{\infty}\left(U ; C^{1}\left(\bar{D} ; C_{\#}(\bar{Y})\right)\right)$. Next, we observe that

$$
\left(A^{\varepsilon} \nabla u_{1}^{\varepsilon}(\boldsymbol{z} ; x)-A_{0} \nabla u_{0}(\boldsymbol{z} ; x)\right)_{i}=\varepsilon \frac{\partial}{\partial x_{j}}\left(\alpha_{i j}^{k}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial u_{0}(\boldsymbol{z} ; x)}{\partial x_{k}}\right)+\left(r_{\varepsilon}\right)_{i}(\boldsymbol{z} ; x)
$$

where
$\left(r_{\varepsilon}\right)_{i}(\boldsymbol{z} ; x)=-\left.\varepsilon \frac{\partial \alpha_{i j}^{k}(\boldsymbol{z} ; x, y)}{\partial x_{j}}\right|_{y=x / \varepsilon} \frac{\partial u_{0}(\boldsymbol{z} ; x)}{\partial x_{k}}-\varepsilon \alpha_{i j}^{k}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial^{2} u_{0}(\boldsymbol{z} ; x)}{\partial x_{k} \partial x_{j}}+\varepsilon w^{k}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) A_{i j}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial^{2} u_{0}(\boldsymbol{z} ; x)}{\partial x_{k} \partial x_{j}}$.
As $\alpha_{i j}^{k} \in L^{\infty}\left(U ; C^{1}(\bar{D} ; C(\bar{Y}))\right),\left\|\left(r_{\varepsilon}\right)_{i}(\boldsymbol{z}, \cdot)\right\|_{L^{2}(D)} \leq c \varepsilon$ for all $\boldsymbol{z} \in U$. As $\alpha_{i j}^{k}=-\alpha_{j i}^{k}$,

$$
\left\|\operatorname{div} A^{\varepsilon} \nabla u_{1}^{\varepsilon}(\boldsymbol{z} ; x)-\operatorname{div} A_{0} \nabla u_{0}(\boldsymbol{z} ; x)\right\|_{H^{-1}(D)}<c \varepsilon,
$$

where the constant $c$ does not depend on $\boldsymbol{z}$. As $\operatorname{div} A_{0} \nabla u_{0}=\operatorname{div} A^{\varepsilon} \nabla u^{\varepsilon}$, we find that

$$
\left\|\operatorname{div} A^{\varepsilon} \nabla u_{1}^{\varepsilon}(\boldsymbol{z} ; x)-\operatorname{div} A^{\varepsilon} \nabla u^{\varepsilon}(\boldsymbol{z} ; x)\right\|_{H^{-1}(D)} \leq c \varepsilon,
$$

where the constant $c$ does not depend on $z$. Let $\tau^{\varepsilon} \in C_{0}^{\infty}(D)$ such that $\tau^{\varepsilon}=1$ outside an $\varepsilon$ neighbourbood of $\partial D$ and such that $\varepsilon\left|\nabla \tau^{\varepsilon}(x)\right| \leq c$ for all $\varepsilon>0$. We consider the function

$$
w_{1}^{\varepsilon}(\boldsymbol{z} ; x)=u_{0}(\boldsymbol{z} ; x)+\varepsilon \tau^{\varepsilon}(x) w^{k}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial u_{0}}{\partial x_{k}}(\boldsymbol{z} ; x)=u_{1}^{\varepsilon}(x)-\varepsilon\left(1-\tau^{\varepsilon}(x)\right) w^{k}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial u_{0}(\boldsymbol{z} ; x)}{\partial x_{k}} .
$$

We then get

$$
\begin{array}{r}
\frac{\partial}{\partial x_{j}}\left(u_{1}^{\varepsilon}-w_{1}^{\varepsilon}\right)(\boldsymbol{z} ; x)=-\varepsilon \frac{\partial \tau^{\varepsilon}}{\partial x_{j}}(x) w^{k}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial u_{0}(\boldsymbol{z} ; x)}{\partial x_{k}}+\left(1-\tau^{\varepsilon}(x)\right) \frac{\partial w^{k}}{\partial y_{j}}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial u_{0}}{\partial x_{k}}(\boldsymbol{z} ; x)+ \\
\varepsilon\left(1-\tau^{\varepsilon}(x)\right) w^{k}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial^{2} u_{0}(\boldsymbol{z} ; x)}{\partial x_{k} \partial x_{j}} .
\end{array}
$$

For $\varepsilon>0$ sufficiently small, let $D^{\varepsilon} \subset D$ be an $\varepsilon$ neighbourhood of $\partial D$. As $\partial D$ is Lipschitz, for all functions $\phi \in C^{\infty}(\bar{D})$ it holds $\|\phi\|_{L^{2}\left(D^{\varepsilon}\right)}^{2} \leq c \varepsilon^{2}\|\phi\|_{H^{1}(D)}^{2}+c \varepsilon\|\phi\|_{L^{2}(\partial D)}^{2}$, so for all $\phi \in H^{1}(D)$ we have $\|\phi\|_{L^{2}\left(D^{\varepsilon}\right)} \leq$ $c \varepsilon^{1 / 2}\|\phi\|_{H^{1}(D)}$. Therefore, since $u_{0}(\boldsymbol{z} ; x) \in L^{\infty}\left(U ; H^{2}(D)\right)$ and $w^{l}(\boldsymbol{z} ; x, \boldsymbol{y}) \in L^{\infty}\left(U \times \bar{D} ; C^{1}(\bar{Y})\right)$

$$
\sup _{\boldsymbol{z} \in U}\left\|u_{1}^{\varepsilon}(\boldsymbol{z} ; \cdot)-w_{1}^{\varepsilon}(\boldsymbol{z} ; \cdot)\right\|_{H^{1}(D)} \leq c \varepsilon^{1 / 2}
$$

where the constant $c$ does not depend on $\varepsilon$. Thus,

$$
\sup _{\boldsymbol{z} \in U}\left\|\operatorname{div}\left(A^{\varepsilon}\left(\nabla u_{1}^{\varepsilon}(\boldsymbol{z} ; \cdot)-\nabla w_{1}^{\varepsilon}(\boldsymbol{z} ; \cdot)\right)\right)\right\|_{H^{-1}(D)} \leq c \varepsilon^{1 / 2}
$$

so

$$
\sup _{\boldsymbol{z} \in U}\left\|\operatorname{div}\left(A^{\varepsilon}\left(\nabla u^{\varepsilon}(\boldsymbol{z} ; \cdot)-\nabla w_{1}^{\varepsilon}(\boldsymbol{z} ; \cdot)\right)\right)\right\|_{H^{-1}(D)} \leq c \varepsilon^{1 / 2}
$$

From Assumption 1.1, we get

$$
\sup _{\boldsymbol{z} \in U}\left\|u^{\varepsilon}(\boldsymbol{z} ; \cdot)-w_{1}^{\varepsilon}(\boldsymbol{z} ; \cdot)\right\|_{H_{0}^{1}(D)} \leq c \varepsilon^{1 / 2} .
$$

Hence we have proved that there exists $c>0$ such that for every $0<\varepsilon \leq 1$ holds

$$
\sup _{\boldsymbol{z} \in U}\left\|u^{\varepsilon}(\boldsymbol{z} ; \cdot)-u_{1}^{\varepsilon}(\boldsymbol{z} ; \cdot)\right\|_{H^{1}(D)} \leq c \varepsilon^{1 / 2} .
$$

We then have the following estimate of the homogenization error for the two scale problems (1.5).
Theorem 5.2 Assume that $A(\boldsymbol{z} ; x, y) \in L^{\infty}\left(U ; C^{1}\left(\bar{D} ; C_{\#}^{1}(Y)\right)_{s y m}^{d \times d}\right)$, that $f \in L^{2}(D)$ and that the domain $D$ is convex. There exists a constant $c>0$ independent of $\boldsymbol{z} \in U$ such that

$$
\left\|u^{\varepsilon}(\boldsymbol{z} ; \cdot)-\left[\nabla u_{0}(\boldsymbol{z} ; \cdot)+\nabla_{y} u_{1}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)\right]\right\|_{L^{2}\left(U ; H^{1}(D)\right)} \leq c \varepsilon^{1 / 2} .
$$

Proof Under the assumptions of the theorem, we have shown that $u_{0}(\boldsymbol{z} ; \cdot) \in L^{\infty}\left(U ; H^{2}(D)\right)$ and $w^{l}(\boldsymbol{z} ; x, \cdot) \in L^{\infty}\left(U ; C^{1}\left(\bar{D}, H_{\#}^{2}(Y) / \mathbb{R}\right)\right) \subset L^{\infty}\left(U ; C^{1}\left(\bar{D} ; C_{\#}(\bar{Y})\right)\right)$. To apply Proposition 5.1 we show that $\left.w^{l}(\boldsymbol{z} ; x, \cdot) \in L^{\infty}\left(U \times \bar{D}, C_{\#}^{1}(\bar{Y})\right)\right)$. The functions $w^{l}$ satisfy

$$
\begin{equation*}
\int_{Y} A(\boldsymbol{z} ; x, y) \nabla_{y} w^{l}(\boldsymbol{z} ; x, y) \cdot \nabla_{y} \phi(y) d y=\int_{Y} \nabla \cdot\left(A(\boldsymbol{z} ; x, y) e_{l}\right) \phi(y) d y, \quad \forall \phi \in H_{\#}^{1}(Y) . \tag{5.4}
\end{equation*}
$$

As $d \leq 3$, we have the continuous embedding

$$
w^{l} \in L^{\infty}\left(U ; C^{1}\left(\bar{D} ; H^{2}(Y)\right)\right) \subset L^{\infty}\left(U ; C^{1}\left(\bar{D} ; W^{1,5}(Y)\right)\right)
$$

Therefore, there exists a constant $c>0$ independent of $\boldsymbol{z} \in \mathcal{U}$ and of $x \in D$ such that

$$
\forall \boldsymbol{z} \in U \forall x \in D: \quad\left\|w^{l}(\boldsymbol{z} ; x, \cdot)\right\|_{W^{1,5}(\mathbf{Y}) / \mathbb{R}} \leq c
$$

Let $D^{\prime}$ be a smooth and bounded domain that contains the unit cube $Y$ and let $\tau \in \mathcal{D}\left(D^{\prime}\right)$ be such that $\tau(y)=1$ when $y \in Y$. Applying Theorem 6 of [13, pg. 177], we find that there is a constant $c>0$ (which only depends on the Lipschitz constant of $A, \alpha, \beta$ and $\tau)$ such that

$$
\left\|w^{l}\right\|_{W^{2,5}(Y)} \leq c\left(\left\|\operatorname{div}_{y}\left(A(\boldsymbol{z}, x, \cdot) e_{l}\right)\right\|_{L^{5}(Y)}+\left\|w^{l}\right\|_{W^{1,5}(Y)}\right) .
$$

As $d \leq 3$ implies the embedding $W^{2,5}(Y) \subset C^{1}(\bar{Y})$, it holds $w^{l}(\boldsymbol{z} ; x, \cdot) \in L^{\infty}\left(U ; L^{\infty}\left(D ; C_{\#}^{1}(\bar{Y})\right)\right)$. Therefore Proposition 5.1 holds.

It remains to show that $\nabla_{y} u_{1}(z ; x, x / \varepsilon)$ as a function from $U$ to $L^{2}(D)$ is measurable. To this end, we note that for every $\boldsymbol{z}, \boldsymbol{z}^{\prime} \in U$ holds

$$
\begin{aligned}
\int_{Y} A(\boldsymbol{z} ; x, \boldsymbol{y}) \nabla_{\boldsymbol{y}}\left(w^{l}(\boldsymbol{z} ; x, \boldsymbol{y})-w^{l}\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right) & \cdot \nabla_{\boldsymbol{y}} \phi(\boldsymbol{y}) d \boldsymbol{y}=\int_{Y} \nabla \cdot\left(\left(A(\boldsymbol{z} ; x, \boldsymbol{y})-A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)\right) e_{l}\right) \phi(\boldsymbol{y}) d \boldsymbol{y} \\
& +\int_{Y}\left(A\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right)-A(\boldsymbol{z} ; x, \boldsymbol{y})\right) \nabla_{y} w^{l}\left(\boldsymbol{z}^{\prime} ; x, \boldsymbol{y}\right) \cdot \nabla_{\boldsymbol{y}} \phi(\boldsymbol{y}) d \boldsymbol{y}
\end{aligned}
$$

Thus for all $x \in \bar{D}$ and every $\boldsymbol{z}, \boldsymbol{z}^{\prime} \in U$

$$
\left\|w^{l}(\boldsymbol{z} ; x, \cdot)-w^{l}\left(\boldsymbol{z}^{\prime} ; x, \cdot\right)\right\|_{W^{2,5}(Y) / \mathbb{R}} \leq c\left\|A(\boldsymbol{z} ; \cdot, \cdot)-A\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right\|_{C^{1}(\bar{D} \times \bar{Y})^{d \times d}}\left(1+\left\|w^{l}\left(\boldsymbol{z}^{\prime}\right)\right\|_{W^{2,5}(Y) / \mathbb{R}}\right)
$$

Therefore,

$$
\left\|w^{l}(\boldsymbol{z} ; \cdot, \cdot)-w^{l}\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right\|_{L^{\infty}\left(\bar{D}, C^{1}(\bar{Y})\right)} \leq c\left\|A(\boldsymbol{z} ; \cdot, \cdot)-A\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right\|_{C^{1}(\bar{D} \times \bar{Y})^{d \times d}} .
$$

From (4.17),

$$
\begin{aligned}
\left\|\nabla_{y} u_{1}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)-\nabla_{y} u_{1}\left(\boldsymbol{z}^{\prime} ; x, \frac{x}{\varepsilon}\right)\right\|_{L^{2}(D)} & =\left\|\nabla_{y} w^{l}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right) \frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x)-\nabla_{y} w^{l}\left(\boldsymbol{z}^{\prime} ; x, \frac{x}{\varepsilon}\right) \frac{\partial u_{0}}{\partial x_{l}}\left(\boldsymbol{z}^{\prime} ; x\right)\right\| \\
& \leq\left\|A(\boldsymbol{z} ; \cdot, \cdot)-A\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right\|_{C^{1}(\bar{D} \times \bar{Y})} .
\end{aligned}
$$

As in the proof of Proposition 1.8, this shows that for every $0<\varepsilon \leq 1$, the function $\left.\nabla_{y} u_{1}(\boldsymbol{z} ; x, y)\right|_{y=x / \varepsilon}$ as a map from $U$ to $L^{2}(D)$ is strongly measurable. This completes the proof.

Following [6, Def. 2.16], we define a "folding" or averaging operator
Definition 5.3 For $\Phi \in L^{1}(D \times Y)$, we define

$$
\mathcal{U}^{\varepsilon}(\Phi)(x)=\int_{Y} \Phi\left(\varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t,\left\{\frac{x}{\varepsilon}\right\}\right) d t
$$

where $[x / \varepsilon]$ denotes the "integer" part of $x / \varepsilon$ with respect to the period $Y$ and $\{x / \varepsilon\}=x / \varepsilon-[x / \varepsilon]$, where $\Phi(x)=0$ when $x \notin D$.

We shall use the following result from [6, Prop. 2.18]. As we will use it later, we present its proof.
Lemma 5.4 For $\Phi \in L^{1}(D \times Y)$,

$$
\int_{D^{\varepsilon}} \mathcal{U}^{\varepsilon}(\Phi)(x) d x=\int_{D \times Y} \Phi(x, y) d y d x
$$

where $D^{\varepsilon}$ denotes the $2 \varepsilon$ neighbourhood of $D$ and where $\Phi(x)=0$ when $x \notin D$.

Proof Let $I$ be a subset of $\mathbb{Z}^{d}$ such that $D \subset \bigcup_{m \in I} \varepsilon(m+\bar{Y})$. Then

$$
\begin{array}{r}
\int_{D^{\varepsilon}} \int_{Y} \Phi\left(\varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t,\left\{\frac{x}{\varepsilon}\right\}\right) d t d x=\sum_{m \in I} \int_{Y} \int_{\varepsilon(m+Y)} \Phi\left(\varepsilon m+\varepsilon t,\left\{\frac{x}{\varepsilon}\right\}\right) d x d t \\
=\varepsilon^{d} \sum_{m \in I} \int_{Y} \int_{Y} \Phi(\varepsilon m+\varepsilon t, y) d z d y=\int_{D \times Y} \Phi(x, y) d x d y
\end{array}
$$

For the corrector function $u_{1}(\boldsymbol{z} ; x, y)$, we have the following result.
Lemma 5.5 If $u_{0} \in L^{\infty}\left(U ; H^{2}(D)\right)$ and $w^{l} \in L^{\infty}\left(U ; C_{1}^{1}\right)$ for all $l=1, \ldots, d$, then there exists a constant $c$ independent of $\boldsymbol{z}$ such that for every $0<\varepsilon \leq 1$

$$
\sup _{\boldsymbol{z} \in U} \int_{D} \left\lvert\, \nabla_{y} u_{1}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)-\mathcal{U}^{\varepsilon}\left(\left.\nabla_{y}\left(u_{1}(\boldsymbol{z} ; \cdot, \cdot)\right)(x)\right|^{2} d x \leq c \varepsilon^{2} .\right.\right.
$$

Proof As

$$
u_{1}(\boldsymbol{z} ; x, \boldsymbol{y})=\sum_{l=1}^{d} \frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x) w^{l}(\boldsymbol{z} ; x, \boldsymbol{y}),
$$

it is sufficient to show that for each $l=1, \ldots, d$

$$
\int_{D}\left|\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x) \nabla_{y} w^{l}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)-\int_{Y} \frac{\partial u_{0}}{\partial x_{l}}\left(\boldsymbol{z} ; \varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t\right) \nabla_{y} w^{l}\left(\boldsymbol{z} ; \varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t, \frac{x}{\varepsilon}\right) d t\right|^{2} d x \leq c \varepsilon^{2}
$$

The expression on the left hand side is bounded by

$$
\begin{array}{r}
\int_{D} \int_{Y}\left|\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x) \nabla_{y} w^{l}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)-\frac{\partial u_{0}}{\partial x_{l}}\left(\boldsymbol{z} ; \varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t\right) \nabla_{y} w^{l}\left(\boldsymbol{z} ; \varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t, \frac{x}{\varepsilon}\right)\right|^{2} d t d x \\
\leq 2 \int_{D} \int_{Y}\left|\left(\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x)-\frac{\partial u_{0}}{\partial x_{l}}\left(\boldsymbol{z} ; \varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t\right)\right) \nabla_{y} w^{l}\left(\boldsymbol{z} ; \varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t, \frac{x}{\varepsilon}\right)\right|^{2} d t d x+ \\
2 \int_{D} \int_{Y}\left|\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x)\right|^{2}\left|\nabla_{y} w^{l}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)-\nabla_{\boldsymbol{y}} w^{l}\left(\boldsymbol{z} ; \varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t, \frac{x}{\varepsilon}\right)\right|^{2} d t d x
\end{array}
$$

As $w^{l} \in L^{\infty}\left(U ; \boldsymbol{C}_{1}^{1}\right)$, there exists a constant $c$ independent of $\boldsymbol{z} \in U$ such that for $\varepsilon$ sufficiently small

$$
\operatorname{ess} \sup _{\boldsymbol{z} \in U} \sup _{t \in Y}\left|\nabla_{y} w^{l}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)-\nabla_{y} w^{l}\left(\boldsymbol{z} ; \varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t, \frac{x}{\varepsilon}\right)\right| \leq c \varepsilon .
$$

Therefore for all $\boldsymbol{z} \in U$

$$
\begin{array}{r}
\int_{D}\left|\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x) \nabla_{y} w^{l}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)-\mathcal{U}^{\varepsilon}\left(\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z}, \cdot) \nabla_{y} w^{l}(\boldsymbol{z} ; \cdot, \cdot)\right)(x)\right|^{2} d x \\
\quad \leq c \int_{D} \int_{Y}\left|\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x)-\frac{\partial u_{0}}{\partial x_{l}}\left(\boldsymbol{z} ; \varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t\right)\right|^{2} d t d x+c \varepsilon^{2}
\end{array}
$$

Next we claim that for all $\boldsymbol{z} \in U$

$$
\int_{D} \int_{Y}\left|\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x)-\frac{\partial u_{0}}{\partial x_{l}}\left(\boldsymbol{z} ; \varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t\right)\right|^{2} d t d x \leq c \varepsilon^{2}
$$

where $c$ is independent of $\boldsymbol{z}$. To prove this, let $\phi(x)$ be a smooth function. Then

$$
\begin{aligned}
\int_{D}\left(\phi(x)-\mathcal{U}^{\varepsilon}(\phi)(x)\right)^{2} d x \leq & \int_{D} \int_{Y}\left|\phi(x)-\phi\left(\varepsilon\left[\frac{x}{\varepsilon}\right]+\varepsilon t\right)\right|^{2} d t \\
\leq & \sum_{i=1}^{d} \int_{D} \int_{Y} \left\lvert\, \phi\left(\varepsilon\left[\frac{x_{1}}{\varepsilon}\right]+\varepsilon t_{1}, \ldots, \varepsilon\left[\frac{x_{i-1}}{\varepsilon}\right]+\varepsilon t_{i-1}, x_{i}, \ldots, x_{d}\right)-\right. \\
& \left.\phi\left(\varepsilon\left[\frac{x_{1}}{\varepsilon}\right]+\varepsilon t_{1}, \ldots, \varepsilon\left[\frac{x_{i}}{\varepsilon}\right]+\varepsilon t_{i}, x_{i+1}, \ldots, x_{d}\right)\right|^{2} d t d x \\
\leq & \sum_{i=1}^{d} \int_{D} \int_{Y}\left|\varepsilon \int_{t_{i}}^{\left\{x_{i} / \varepsilon\right\}} \frac{\partial \phi}{\partial x_{i}}\left(\varepsilon\left[\frac{x_{1}}{\varepsilon}\right]+\varepsilon t_{1}, \ldots, \varepsilon\left[\frac{x_{i}}{\varepsilon}\right]+\varepsilon \zeta_{i}, x_{i+1}, \ldots, x_{d}\right) d \zeta_{i}\right|^{2} d t d x \\
\leq & \varepsilon^{2} \sum_{i=1}^{d} \int_{D} \int_{Y} \int_{0}^{1}\left|\frac{\partial \phi}{\partial x_{i}}\left(\varepsilon\left[\frac{x_{1}}{\varepsilon}\right]+\varepsilon t_{1}, \ldots, \varepsilon\left[\frac{x_{i}}{\varepsilon}\right]+\varepsilon \zeta_{i}, x_{i+1}, \ldots, x_{d}\right)\right|^{2} d \zeta_{i} d t d x \\
\leq & \varepsilon^{2} \sum_{i=1}^{d} \int_{D}\left|\frac{\partial \phi}{\partial x_{i}}\right|^{2} d x .
\end{aligned}
$$

The last inequality is derived from Lemma 5.4, freezing the variables $x_{i+1}, \ldots, x_{d}$. Fix $\boldsymbol{z} \in U$ and $0<\varepsilon \leq 1$ arbitrary, and consider a sequence $\left\{\phi_{n}\right\}_{n} \subset C^{\infty}(\bar{D})$ which converges to $\partial u_{0}(\boldsymbol{z} ; x) / \partial x_{l}$ in $H^{1}(D)$. As $n \rightarrow \infty$,

$$
\begin{aligned}
\int_{D}\left|\mathcal{U}^{\varepsilon}\left(\phi_{n}\right)(x)-\mathcal{U}^{\varepsilon}\left(\frac{\partial u_{0}(\boldsymbol{z} ; x)}{\partial x_{l}}\right)(x)\right|^{2} d x \leq & \int_{D} \mathcal{U}^{\varepsilon}\left(\left(\phi_{n}(x)-\frac{\partial u_{0}(\boldsymbol{z} ; x)}{\partial x_{l}}\right)^{2}\right)(x) d x \\
& \leq \int_{D}\left(\phi_{n}(x)-\frac{\partial u_{0}(\boldsymbol{z} ; x)}{\partial x_{l}}\right)^{2} d x \rightarrow 0
\end{aligned}
$$

Therefore

$$
\begin{aligned}
& \int_{D}\left(\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x)-\mathcal{U}^{\varepsilon}\left(\frac{\partial u_{0}(\boldsymbol{z} ; \cdot)}{\partial x_{l}}\right)(x)\right)^{2} d x \\
& \leq 3 \int_{D}\left(\frac{\partial u_{0}}{\partial x_{l}}-\phi_{n}\right)^{2} d x+3 \int_{D}\left(\phi_{n}-\mathcal{U}^{\varepsilon}\left(\phi_{n}\right)\right)^{2} d x+3 \int_{D}\left(\mathcal{U}^{\varepsilon}\left(\phi_{n}\right)-\mathcal{U}^{\varepsilon}\left(\frac{\partial u_{0}}{\partial x_{l}}\right)\right)^{2} d x \\
& \leq 6 \int_{D}\left(\frac{\partial u_{0}}{\partial x_{l}}-\phi_{n}\right)^{2} d x+3 \varepsilon^{2} \sum_{i=1}^{d} \int_{D}\left|\frac{\partial \phi_{n}}{\partial x_{i}}\right|^{2} d x
\end{aligned}
$$

As $n \rightarrow \infty$, we have

$$
\int_{D}\left(\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x)-\mathcal{U}^{\varepsilon}\left(\frac{\partial u_{0}(\boldsymbol{z} ; \cdot)}{\partial x_{l}}\right)(x)\right)^{2} \leq 3 \varepsilon^{2} \sum_{i=1}^{d} \int_{D}\left|\frac{\partial^{2} u_{0}}{\partial x_{i} x_{l}}\right|^{2} d x
$$

Thus

$$
\int_{D}\left|\frac{\partial u_{0}}{\partial x_{l}}(\boldsymbol{z} ; x) \nabla_{y} w^{l}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)-\mathcal{U}^{\varepsilon}\left(\frac{\partial u_{0}(\boldsymbol{z} ; x)}{\partial x_{l}} \nabla_{y} w^{l}(\boldsymbol{z} ; x, \boldsymbol{y})\right)(x)\right|^{2} d x \leq c \varepsilon^{2}
$$

for a constant $c$ which does not depend on $\boldsymbol{z} \in U$.
Lemma 5.6 With Assumption 4.3, $\mathcal{U}^{\varepsilon}\left(\nabla_{\boldsymbol{y}} u_{1}(\boldsymbol{z} ; \cdot, \cdot)\right)$ as a map from $U$ to $L^{2}(D)$ is measurable.
Proof First we note that $\left(\mathcal{U}^{\varepsilon}(\Phi)(x)\right)^{2} \leq \mathcal{U}^{\varepsilon}\left(\Phi^{2}\right)(x)$ for a.e. $x$ for all functions $\Phi \in L^{2}(\Omega \times Y)$. Thus

$$
\forall \boldsymbol{z} \in U: \quad\left\|\mathcal{U}^{\varepsilon}\left(\nabla_{y} u_{1}(\boldsymbol{z} ; \cdot, \cdot)-\nabla_{y} u_{1}\left(\boldsymbol{z}^{\prime}, \cdot, \cdot\right)\right)(\cdot)\right\|_{L^{2}(D)} \leq\left\|\nabla_{y} u_{1}(\boldsymbol{z} ; \cdot, \cdot)-\nabla_{y} u_{1}\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right\|_{L^{2}(D \times Y)}
$$

From the proof of Proposition 4.5, there exists $c>0$ independent of $\boldsymbol{z} \in U$ such that

$$
\left\|\mathcal{U}^{\varepsilon}\left(\nabla_{\boldsymbol{y}} u_{1}(\boldsymbol{z} ; \cdot, \cdot)(\cdot)-\nabla_{y} u_{1}\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)(\cdot)\right)\right\|_{L^{2}(D)} \leq c\left\|A(\boldsymbol{z} ; \cdot, \cdot)-A\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right\|_{\left(\boldsymbol{C}_{2}^{1}\right) d^{d \times d}}
$$

An argument similar to that in the proof of Proposition 1.8 shows that $\mathcal{U}^{\varepsilon}\left(\nabla_{y} u_{1}\right)$ as a map from $U$ to $L^{2}(D)$ is measurable.

Proposition 5.7 If $A(\boldsymbol{z} ; x, y) \in L^{\infty}\left(U ; \boldsymbol{C}_{1}^{1}\right), f \in L^{2}(D)$ and if the domain $D$ is convex, then there is a constant $c>0$ such that, for $\varepsilon>0$ sufficiently small, there holds

$$
\left\|\nabla_{y} u_{1}\left(\boldsymbol{z} ; x, \frac{x}{\varepsilon}\right)-\mathcal{U}^{\varepsilon}\left(\nabla_{y} u_{1}(\boldsymbol{z} ; \cdot, \cdot)\right)(x)\right\|_{L^{2}\left(U ; L^{2}(D)\right)}^{2} \leq c \varepsilon
$$

Proof We need to show that $w^{l} \in L^{\infty}\left(U ; \boldsymbol{C}_{1}^{1}\right)$. We will do this by analyzing a suitable difference quotient. To define it, we introduce for $x \in \bar{D}$ and for $\underline{\delta} \in \mathbb{R}^{d}$, the translation operator

$$
\tau_{\underline{\delta}}(x):=x+\underline{\delta} \in \bar{D} .
$$

We then have

$$
\begin{array}{rl}
\int_{Y} & A(\boldsymbol{z} ; x, y) \nabla_{y}\left(\frac{w^{l}\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right)-w^{l}(\boldsymbol{z} ; x, y)}{|\underline{\delta}|}\right) \cdot \nabla_{y} \phi(y) d y \\
& =\int_{Y} \nabla_{y} \cdot\left(\frac{A\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right)-A(\boldsymbol{z} ; x, y)}{|\underline{\delta}|} e_{l}\right) \phi(y) d y \\
& -\int_{Y} \frac{A\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right)-A(\boldsymbol{z} ; x, y)}{|\underline{\delta}|} \nabla_{y} w^{l}\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right) \cdot \nabla_{\boldsymbol{y}} \phi(y) d y \tag{5.5}
\end{array}
$$

for all $\phi \in H_{\#}^{1}(Y)$. Let $\psi(\boldsymbol{z} ; x, \cdot) \in H_{\#}^{1}(Y) / \mathbb{R}$ denote the solution of

$$
\begin{align*}
& \int_{Y} A(\boldsymbol{z} ; x, y) \nabla_{y} \psi(\boldsymbol{z} ; x, y) \cdot \nabla_{y} \phi(y) d y=\int_{Y} \nabla_{y} \cdot\left(\frac{\partial A(\boldsymbol{z} ; x, y)}{\partial x_{i}} e_{l}\right) \phi(y) d y \\
&-\int_{Y} \frac{\partial A(\boldsymbol{z} ; x, y)}{\partial x_{i}} \nabla_{y} w^{l}(\boldsymbol{z} ; x, y) \cdot \nabla_{y} \phi(y) d y \quad \forall \phi \in H_{\#}^{1}(Y) \tag{5.6}
\end{align*}
$$

We have

$$
\begin{aligned}
& \int_{Y} A(\boldsymbol{z} ; x, y) \nabla_{y}\left(\frac{w^{l}\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right)-w^{l}(\boldsymbol{z} ; x, y)}{|\underline{\delta}|}-\psi(\boldsymbol{z} ; x, y)\right) \cdot \nabla_{y} \phi(y) d y \\
& =\int_{Y} \nabla_{y} \cdot\left(\left(\frac{A\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right)-A(\boldsymbol{z} ; x, y)}{|\underline{\delta}|}-\frac{\partial A(\boldsymbol{z} ; x, y)}{\partial x_{i}}\right) e_{l}\right) \phi(y) d y \\
& -\int_{Y}\left(\frac{A\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right)-A(\boldsymbol{z} ; x, y)}{\mid \underline{\delta \mid}}-\frac{\partial A(\boldsymbol{z}, x, y)}{\partial x_{i}}\right) \nabla_{y} w^{l}\left(\boldsymbol{z}, \tau_{\underline{\delta}}(x), y\right) \cdot \nabla_{y} \phi(y) d y \\
& -\int_{Y} \frac{\partial A}{\partial x_{i}}(\boldsymbol{z} ; x, y) \nabla_{y}\left(w^{l}\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right)-w^{l}(\boldsymbol{z}, x, y)\right) \cdot \nabla_{y} \phi(y) d y
\end{aligned}
$$

Therefore

$$
\begin{align*}
& \left\|\frac{w^{l}\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), \cdot\right)-w^{l}(\boldsymbol{z} ; x, \cdot)}{|\underline{\delta}|}-\psi(\boldsymbol{z} ; x, \cdot)\right\|_{W^{2,5}(Y) / \mathbb{R}} \leq \\
& c\left\|\frac{A\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), \cdot\right)-A(\boldsymbol{z} ; x, \cdot)}{|\underline{\delta}|}-\frac{\partial A(\boldsymbol{z} ; x, \cdot)}{\partial x_{i}}\right\|_{\left.C^{1}(\bar{Y})\right)}\left(1+\left\|w^{l}\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), \cdot\right)\right\|_{W^{2,5}(Y) / \mathbb{R}}\right) \\
& +c\left\|w^{l}\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right)-w^{l}(\boldsymbol{z} ; x, y)\right\|_{W^{2,5}(\mathbf{Y}) / \mathbb{R}} . \tag{5.7}
\end{align*}
$$

From (5.5), we find

$$
\left\|w^{l}\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right)-w^{l}(\boldsymbol{z} ; x, y)\right\|_{W^{2,5}(Y) / \mathbb{R}} \leq c\left\|A\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), y\right)-A(\boldsymbol{z} ; x, y)\right\|_{C^{1}(\bar{D} \times \bar{Y})}\left(1+\left\|w^{l}\left(\boldsymbol{z} ; \tau_{\underline{\delta}}(x), \cdot\right)\right\|_{W^{2,5}(Y)}\right)
$$ which converges to 0 when $\delta \rightarrow 0$. Thus the right hand side of (5.7) converges to 0 as $\delta \rightarrow 0$. Therefore

$$
\frac{\partial w^{l}(\boldsymbol{z} ; x, y)}{\partial x_{i}}=\psi(\boldsymbol{z} ; x, y) \text { in } W^{2,5}(Y) / \mathbb{R} \subset C^{1}(\bar{Y})
$$

From (5.6), we deduce that $\|\psi\|_{W^{2,5}(Y) / \mathbb{R}}$ is uniformly bounded for all $\boldsymbol{z} \in U$ and is continuous with respect to $x \in D$ due to the continuity of $\nabla_{x} A$ and $w^{l}$ in $W^{2,5}(Y) / \mathbb{R}$ so $w^{l}(\boldsymbol{z} ; x, y) \in L^{\infty}\left(U, C^{1}\left(\bar{D}, W^{2,5}(Y)\right)\right) \subset$ $L^{\infty}\left(U, \boldsymbol{C}_{1}^{1}\right)$. Therefore Lemma 5.5 holds uniformly for all $\boldsymbol{z} \in U$. Lemma 5.6 implies the assertion.

For the solution $\boldsymbol{u}_{\Lambda_{N}}$ of the semidiscrete Galerkin approximation (2.3), we denote $\boldsymbol{u}_{\Lambda_{N}}=\left(u_{0 \Lambda_{N}}, u_{1 \Lambda_{N}}\right)$. We have the following corrector result for the approximation (2.3).

Theorem 5.8 Assume that $A \in L^{\infty}\left(U ;\left(\boldsymbol{C}_{1}^{1}\right)^{d \times d}\right), f \in L^{2}(D)$ and that $D$ is convex. If Assumptions 1.1, 1.2 and 3.2 hold for some $0<p<1$, there exists a constant $c>0$ (independent of $\varepsilon$ and of $N$ ) such that for $\varepsilon>0$ sufficiently small and for $N$ sufficiently large, it holds that

$$
\left\|u^{\varepsilon}-\left[\nabla u_{0}+\mathcal{U}^{\varepsilon}\left(\nabla_{y} u_{1 \Lambda_{N}}\right)\right]\right\|_{L^{2}\left(U ; L^{2}(D)\right)} \leq c\left(\varepsilon^{1 / 2}+N^{-\sigma}\right)
$$

where $\sigma=1 / p-1 / 2>1 / 2$ and the sets $\Lambda_{N}$ are as in Theorem 3.9
Proof From Theorem 5.2 and Proposition 5.7, we get

$$
\left\|u^{\varepsilon}-\left[\nabla u_{0}+\mathcal{U}^{\varepsilon}\left(\nabla_{y} u_{1}\right)\right]\right\|_{L^{\infty}\left(U ; L^{2}(D)\right)} \leq c \varepsilon^{1 / 2}
$$

We note that

$$
\begin{aligned}
& \int_{U} \int_{D}\left|\mathcal{U}^{\varepsilon}\left(\nabla_{y} u_{1}(\boldsymbol{z} ; \cdot, \cdot)-\nabla_{y} u_{1 \Lambda_{N}}(\boldsymbol{z} ; \cdot, \cdot)\right)\right|^{2} d x d \rho(\boldsymbol{z}) \\
& \leq \int_{U} \int_{D} \mathcal{U}^{\varepsilon}\left(\left|\nabla_{y} u_{1}(\boldsymbol{z} ; \cdot, \cdot)-\nabla_{y} u_{1 \Lambda_{N}}(\boldsymbol{z} ; \cdot, \cdot)\right|^{2}\right) d x d \rho(\boldsymbol{z}) \\
& \leq \int_{U} \int_{D} \int_{Y}\left|\nabla_{y} u_{1}(\boldsymbol{z} ; x, y)-\nabla_{y} u_{1 \Lambda_{N}}(\boldsymbol{z} ; x, y)\right|^{2} d y d x d \rho(\boldsymbol{z}) \\
& \leq c N^{-2 \sigma}
\end{aligned}
$$

where the last estimate is deduced from Theorem 3.9. We then get the conclusion.

### 5.2 Correctors for multiple scale problems

For problems with more than two scales, an error estimate analogous to (5.1) appears not to be available. For such problems we will now prove a corrector result; however, we will not give an explicit order of convergence. Moreover, this result does not require any extra regularity beyond the smoothness required for the existence of the $n+1$-scale limit. We start our analysis with the definition of a corrector.

Definition 5.9 The $n+1$-scale "unfolding" operator $\mathcal{T}_{n}^{\varepsilon}: L^{1}(D) \rightarrow L^{1}(D \times \mathbf{Y})$ is defined by (see also [6]),

$$
\mathcal{T}_{n}^{\varepsilon}(\phi)(x, \boldsymbol{y})=\phi\left(\varepsilon_{1}\left[\frac{x}{\varepsilon_{1}}\right]+\varepsilon_{2}\left[\frac{y_{1}}{\varepsilon_{2} / \varepsilon_{1}}\right]+\ldots+\varepsilon_{n}\left[\frac{y_{n-1}}{\varepsilon_{n} / \varepsilon_{n-1}}\right]+\varepsilon_{n} y_{n}\right)
$$

where the function $\phi$ is understood as 0 outside $D$.
Denoting for $\varepsilon>0$ sufficiently small by $D^{\varepsilon}$ the $2 \varepsilon$ neighbourbood of $D$, we have

$$
\begin{equation*}
\int_{D} \phi d x=\int_{D^{\varepsilon}} \int_{\mathbf{Y}} \mathcal{T}_{n}^{\varepsilon}(\phi) d \boldsymbol{y} d x \tag{5.8}
\end{equation*}
$$

Fixing $\boldsymbol{z} \in U$, as $\varepsilon \rightarrow 0$, we can show that

$$
\begin{equation*}
\mathcal{T}_{n}^{\varepsilon}\left(\nabla u^{\varepsilon}(\boldsymbol{z})\right) \rightharpoonup \boldsymbol{\nabla} \boldsymbol{u} \text { in } L^{2}(D \times \mathbf{Y}) \tag{5.9}
\end{equation*}
$$

where $\boldsymbol{u}$ is as defined in Theorem 1.12. Following [6], we next define the "folding" operator $\mathcal{U}_{n}^{\varepsilon}$ by
Definition 5.10 For $\Phi \in L^{1}(D \times \mathbf{Y})$ (understood to vanish when $x \notin D$ ) and for $\varepsilon>0$ sufficiently small, the "folding" operator $\mathcal{U}_{n}^{\varepsilon}(\Phi) \in L^{1}(D)$ is defined as

$$
\begin{aligned}
\mathcal{U}_{n}^{\varepsilon}(\Phi)(x)=\int_{Y_{1}} & \ldots \int_{Y_{n}} \Phi\left(\varepsilon_{1}\left[\frac{x}{\varepsilon_{1}}\right]+\varepsilon_{1} t_{1}, \frac{\varepsilon_{2}}{\varepsilon_{1}}\left[\frac{\varepsilon_{1}}{\varepsilon_{2}}\left\{\frac{x}{\varepsilon_{1}}\right\}\right]+\frac{\varepsilon_{2}}{\varepsilon_{1}} t_{2}, \ldots\right. \\
& \left.\frac{\varepsilon_{n}}{\varepsilon_{n-1}}\left[\frac{\varepsilon_{n-1}}{\varepsilon_{n}}\left\{\frac{x}{\varepsilon_{n-1}}\right\}\right]+\frac{\varepsilon_{n}}{\varepsilon_{n-1}} t_{n},\left\{\frac{x}{\varepsilon_{n}}\right\}\right) d t_{n} \ldots d t_{1}
\end{aligned}
$$

We have the following measurability result.
Lemma 5.11 Under Assumption 1.2, for the solution $\boldsymbol{u}(\boldsymbol{z})$ of the parametric, deterministic problem (1.19), the function $\mathcal{U}_{n}^{\varepsilon}(\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z}))(x)$ (with $\boldsymbol{\nabla} \boldsymbol{u}$ as in (1.18)) as a map from $U$ to $L^{2}(D)^{d}$ is measurable.

Proof For the functions $\Phi \in L^{1}(D \times \mathbf{Y})$ (which are understood as 0 when $x \notin D$ ),

$$
\int_{D^{\varepsilon}} \mathcal{U}_{n}^{\varepsilon}(\Phi)(x) d x=\int_{D} \int_{\mathbf{Y}} \Phi d \boldsymbol{y} d x
$$

We note further that for a.e. $x \in D$,

$$
\left(\mathcal{U}_{n}^{\varepsilon}(\Phi)(x)\right)^{2} \leq \mathcal{U}_{n}^{\varepsilon}\left(\Phi^{2}\right)(x)
$$

From this we obtain, for every $\boldsymbol{z} \in U$ and with $\boldsymbol{\nabla} \boldsymbol{u}$ as in (1.18),
$\int_{D}\left|\mathcal{U}_{n}^{\varepsilon}\left(\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z})-\boldsymbol{\nabla} \boldsymbol{u}\left(\boldsymbol{z}^{\prime}\right)\right)(x)\right|^{2} d x \leq \int_{D} \mathcal{U}_{n}^{\varepsilon}\left(\left|\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z})-\boldsymbol{\nabla} \boldsymbol{u}\left(\boldsymbol{z}^{\prime}\right)\right|^{2}\right)(x) d x \leq \int_{D} \int_{\mathbf{Y}}\left|\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z})-\boldsymbol{\nabla} \boldsymbol{u}\left(\boldsymbol{z}^{\prime}\right)\right|^{2} d \boldsymbol{y} d x$.
We then get from (1.22) that there exists a constant $c$ such that

$$
\left\|\mathcal{U}_{n}^{\varepsilon}\left(\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z})-\boldsymbol{\nabla} \boldsymbol{u}\left(\boldsymbol{z}^{\prime}\right)\right)\right\|_{L^{2}(D)} \leq c\left\|A(\boldsymbol{z} ; \cdot, \cdot)-A\left(\boldsymbol{z}^{\prime} ; \cdot, \cdot\right)\right\|_{L^{\infty}(D \times \mathbf{Y})} .
$$

The proof then follows along the lines of the proof of Proposition 1.8.
We are now in position to prove a corrector result for the best $N$ term approximation. It states that the gpc approximation of the high dimensional limit problem describes $\mathbb{P}$-a.s. all oscillations of the physical solution at small $\varepsilon$.

Theorem 5.12 Under Assumption 1.2,

$$
\lim _{\substack{\varepsilon \rightarrow 0 \\ N \rightarrow \infty}}\left\|\nabla u^{\varepsilon}-\mathcal{U}_{n}^{\varepsilon}\left(\boldsymbol{\nabla} \boldsymbol{u}_{\Lambda_{N}}\right)\right\|_{L^{2}\left(U ; L^{2}(D)\right)}=0,
$$

where $\Lambda_{N}$ is a subset of $\mathcal{F}$ corresponding to $N$ largest terms of the sequence $\left(\left\|\boldsymbol{u}_{\nu}\right\|_{\mathrm{v}}\right)_{\nu \in \mathcal{F}}$.
Proof We consider

$$
\int_{D} \int_{\mathbf{Y}} \mathcal{T}_{n}^{\varepsilon}\left(A^{\varepsilon}\right)\left(\mathcal{T}_{n}^{\varepsilon}\left(\nabla u^{\varepsilon}(\boldsymbol{z})\right)(x, \boldsymbol{y})-\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z} ; x, \boldsymbol{y})\right) \cdot\left(\mathcal{T}_{n}^{\varepsilon}\left(\nabla u^{\varepsilon}(\boldsymbol{z})\right)(x, \boldsymbol{y})-\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z} ; x, \boldsymbol{y})\right) d \boldsymbol{y} d x
$$

From (1.5), (1.19), (5.8), (5.9), this expression converges to 0 . Therefore the convergence in (5.9) is indeed strong. Fixing $\boldsymbol{z} \in U$, we obtain as $\varepsilon \rightarrow 0$

$$
\left\|\nabla u^{\varepsilon}(\boldsymbol{z})-\mathcal{U}_{n}^{\varepsilon}(\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z}))\right\|_{L^{2}(D)}=\| \mathcal{U}_{n}^{\varepsilon}\left(\mathcal{T}_{n}^{\varepsilon}\left(\nabla u^{\varepsilon}(\boldsymbol{z})\right)-\mathcal{U}_{n}^{\varepsilon}(\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z}))\left\|_{L^{2}(D)} \leq\right\| \mathcal{T}_{n}^{\varepsilon}\left(\nabla u^{\varepsilon}(\boldsymbol{z})\right)-\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z}) \|_{L^{2}(D \times \mathbf{Y})} \rightarrow 0 .\right.
$$

As $\left\|\nabla u^{\varepsilon}(\boldsymbol{z})-\mathcal{U}_{n}^{\varepsilon}(\boldsymbol{\nabla} \boldsymbol{u}(\boldsymbol{z}))\right\|_{L^{2}(D)}$ is uniformly bounded for all $z \in U$, we conclude that

$$
\lim _{\varepsilon \rightarrow 0}\left\|\nabla u^{\varepsilon}-\mathcal{U}_{n}^{\varepsilon}(\boldsymbol{\nabla} \boldsymbol{u})\right\|_{L^{2}\left(U ; L^{2}(D)\right)}=0 .
$$

Furthermore, for each $\boldsymbol{z} \in U$

$$
\left\|\mathcal{U}_{n}^{\varepsilon}\left(\boldsymbol{\nabla}\left(\boldsymbol{u}_{\Lambda_{N}}(\boldsymbol{z})-\boldsymbol{u}(\boldsymbol{z})\right)\right)\right\|_{L^{2}(D)} \leq\left\|\boldsymbol{\nabla}\left(\boldsymbol{u}_{\Lambda_{N}}(\boldsymbol{z})-\boldsymbol{u}(\boldsymbol{z})\right)\right\|_{L^{2}(D \times \mathbf{Y})} .
$$

Therefore from Theorem 2.1

$$
\lim _{N \rightarrow \infty}\left\|\mathcal{U}_{n}^{\varepsilon}\left(\boldsymbol{\nabla}\left(\boldsymbol{u}_{\Lambda_{N}}-\boldsymbol{u}\right)\right)\right\|_{L^{2}\left(U ; L^{2}(D)\right)}=0
$$

uniformly for all $\varepsilon$. Thus

$$
\lim _{\substack{\varepsilon \rightarrow 0 \\ N \rightarrow \infty}}\left\|\nabla u^{\varepsilon}-\mathcal{U}_{n}^{\varepsilon}\left(\boldsymbol{\nabla} \boldsymbol{u}_{\Lambda_{N}}\right)\right\|_{L^{2}\left(U ; L^{2}(D)\right)}=0 .
$$

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