# A posteriori error estimation for partial differential equations with random input data 

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To my beloved family:
René, Christine, Simon, Aline.

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## Abstract

This thesis is devoted to the derivation of error estimates for partial differential equations with random input data, with a focus on a posteriori error estimates which are the basis for adaptive strategies. Such procedures aim at obtaining an approximation of the solution with a given precision while minimizing the computational costs. If several sources of error come into play, it is then necessary to balance them to avoid unnecessary work.

We are first interested in problems that contain small uncertainties approximated by finite elements. The use of perturbation techniques is appropriate in this setting since only few terms in the power series expansion of the exact random solution with respect to a parameter characterizing the amount of randomness in the problem are required to obtain an accurate approximation. The goal is then to perform an error analysis for the finite element approximation of the expansion up to a certain order. First, an elliptic model problem with random diffusion coefficient with affine dependence on a vector of independent random variables is studied. We give both a priori and a posteriori error estimates for the first term in the expansion for various norms of the error. The results are then extended to higher order approximations and to other sources of uncertainty, such as boundary conditions or forcing term. Next, the analysis of nonlinear problems in random domains is proposed, considering the onedimensional viscous Burgers' equation and the more involved incompressible steady-state Navier-Stokes equations. The domain mapping method is used to transform the equations in random domains into equations in a fixed reference domain with random coefficients. We give conditions on the mapping and the input data under which we can prove the well-posedness of the problems and give a posteriori error estimates for the finite element approximation of the first term in the expansion. Finally, we consider the heat equation with random Robin boundary conditions. For this parabolic problem, the time discretization brings an additional source of error that is accounted for in the error analysis.

The second part of this work consists in the analysis of a random elliptic diffusion problem that is approximated in the physical space by the finite element method and in the stochastic space by the stochastic collocation method on a sparse grid. Considering a random diffusion coefficient with affine dependence on a vector of independent random variables, we derive a residual-based a posteriori error estimate that controls the two sources of error. The stochastic error estimator is then used to drive an adaptive sparse grid algorithm which aims at alleviating the so-called curse of dimensionality inherent to tensor grids. Several numerical examples are given to illustrate the performance of the adaptive procedure.

Key words: PDEs with random inputs, uncertainty quantification, a priori and a posteriori error analysis, finite elements, perturbation techniques, stochastic collocation, elliptic equations, steady Navier-Stokes equations, heat equation

## Résumé

Cette thèse est consacrée à la dérivation d'estimations d'erreur pour des équations aux dérivées partielles contenant des données aléatoires. Un accent particuliers est mis sur les estimateurs a posteriori qui sont à la base d'algorithmes adaptatifs. Ces derniers visent à obtenir une approximation de la solution avec une certaine précision tout en minimisant le coût du calcul. Lorsque plusieurs sources d'erreurs entrent en jeu, il est judicieux de les équilibrer afin d'éviter tout travail inutile.

Nous nous intéressons pour commencer à des problèmes contenant de petites incertitudes résolus par la méthode des éléments finis. Dans ce cas, l'utilisation de méthodes dites de perturbation est indiquée car une bonne approximation de la solution peut être obtenue avec peu de termes dans le développement en série de puissances de la solution exacte par rapport à un paramètre controllant le niveau d'incertitude du problème. Le but principal de ce travail est d'effectuer une analyse d'erreur pour l'approximation par éléments finis du développement à un certain ordre. Nous considérons pour commencer un problème modèle elliptique avec un coefficient de diffusion aléatoire qui dépend de manière affine d'un vecteur de variables aléatoires indépendantes. Des estimations d'erreur a priori et $a$ posteriori sont données pour le premier terme dans le développement de la solution en considérant différentes normes de l'erreur. Les résultats obtenus sont alors généralisés pour des approximations d'ordres supérieurs ainsi que pour des problèmes contenant d'autres sources d'incertitudes, comme par exemple les conditions au bord ou le terme de force. L'étude se poursuit en considérant des problèmes non-linéaires définis sur des domaines aléatoires, tout d'abord l'équation de Burgers à une dimension d'espace puis les équations de Navier-Stokes stationnaires incompressibles. Les problèmes sont reformulés sur un domaine fixe de reference à l'aide d'une transformation introduisant alors des coefficients aléatoires dans les équations. Nous donnons des conditions sur la transformation et les données sous lesquelles les problèmes sont bien posés et nous donnons des estimations d'erreur pour le premier terme du développement. Finalement, nous considérons le problème de la chaleur avec des conditions au bord de type Robin qui contiennent des incertitudes. Pour ce problème parabolique, la discrétisation temporelle ajoute une source supplémentaire d'erreur qui est prise en compte dans l'analyse d'erreur.

Dans la deuxième partie de ce travail, nous analysons un problème de diffusion elliptique avec coefficient aléatoire résolu approximativement par la méthode des éléments finis en espace physique et par la méthode de collocation stochastique avec grille fine en espace stochastique. En considérant un coefficient de diffusion dépendant de manière affine d'un vecteur de
variables aléatoires indépendantes, nous donnons un estimateur d'erreur a posteriori basé sur le résidu qui contrôle les deux sources d'erreur. L'estimateur controlant l'erreur stochastique est ensuite utilisé dans un algorithme construisant de manière adaptative une grille peu dense, permettant ainsi de palier au problème du fléau de la dimension dont souffrent les grilles de type tensiorel. Plusieurs exemples numériques sont donnés pour illustrer les performances de l'algorithme adaptatif.

Mots clefs : EDP avec données aléatoires, quantification des incertitudes, analyse d'erreur a priori et a posteriori, éléments finis, technique de perturbation, collocation stochastique, équations elliptiques, Navier-Stokes stationaire, équation de la chaleur

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## Introduction

Partial differential equations (PDEs) are widely used for modelling problems in many fields such as physics, biology or engineering. Nowadays, uncertainty is often included in mathematical models arising from the simulation of complex systems. The uncertainty can reflect an intrinsic variability of the system (aleatory uncertainty) or our inability to adequately characterize all the inputs (epistemic uncertainty), due for instance to experimental measurements. It can occur in the coefficients, the forcing term, the geometry, the boundary conditions, the initial condition or combinations of them. A possible way to describe the uncertainties present in the model is to use a probability framework. In such a setting, the uncertain input data are characterized with random variables, or more generally random fields, yielding PDEs with random input. In a forward uncertainty quantification (UQ) problem, the goal is then to determine the effect of the uncertainty on the solution or a specific quantity of interest.

Several methods have been developed to tackle the numerical approximation of such problems in both the deterministic and, more recently, the stochastic variables. We give a short overview of the available methods, pointing to some references for an in-depth description, but we have no pretension to be exhaustive.

The best known and most commonly used methods for solving deterministic problems numerically are the finite difference [112,117], the finite element [31,49,61] and the finite volume [85] methods, for which the theory is at a mature stage. Many other methods have been developed, either new methods or extension of the ones mentioned above, such as discontinuous Galerkin [105], boundary element [30], meshless [86] or extended finite element methods [81]. The selection of the method depends upon the type of problems to solve: elliptic, parabolic or hyperbolic.

For the approximation of random PDEs, the most popular method is certainly the Monte-Carlo method (see [63] for instance) which consists in solving the equations for i.i.d. realizations of the random input. The main drawback of this method is its well-known slow convergence rate with respect to the sample size $K$, namely of $\mathscr{O}(1 / \sqrt{K})$. However, the convergence is independent of the dimension of the random space and this method is very easy to use in practice. To improve the convergence rate of the method, some extensions have been introduced such as the quasi-Monte Carlo [54,55] and the multi-level Monte-Carlo [68] methods. Other than MC type methods, we mention the stochastic spectral methods comprising the Stochastic

Galerkin (SG) $[67,90]$ and the Stochastic Collocation (SC) $[7,97,124]$ methods. These methods exploit the possible regularity of the solution with respect to the random input combining the generalized Polynomial Chaos (gPC) expansion of the solution with a Galerkin projection or an interpolation procedure. Finally, in the framework of PDEs with small uncertainties, the perturbation or Neumann series expansion methods [6,37,82,127] appear to be an appropriate choice. For all these methods, an approximation in the physical space can be obtained using any deterministic method mentioned above. In particular, in this thesis we focus on the finite element method.

When a numerical method is used to solve a problem for which the exact solution is not at hand, approximation errors are introduced. An error analysis should then be done to appropriately estimate the various sources of error. In an a priori error analysis, the convergence of the method is assessed under suitable regularity assumptions on the exact solution. The a priori error estimate gives useful information about the asymptotic behaviour of the numerical approximation when the various discretization parameters vary. However, this theoretical bound usually depend on the unknown solution and is thus not a computable quantity. In a posteriori error analysis, the goal is to provide computable error estimators that depend only on the numerical approximation and the input data and that are localized in space. Having such error estimators available can be necessary in many situations. Indeed, if the solution presents local features evolving at fine scale, such as shocks, boundary layers or singularities due to re-entrant corners in physical space, very fine approximation spaces are required to capture them. However, this becomes quickly numerically unaffordable due to the limitations in computer power and memory. A remedy is then to use adaptive strategies based on a (reliable and efficient) a posteriori error estimator, refining only where needed, to get satisfactory accuracy in the approximation while limiting the computational effort. When several sources of error are affecting the numerical solution, the estimator should also furnish an estimation of the contribution of each error component to the total error, so that it can be used for balancing the errors.

The derivation of a posteriori error estimate controlling the finite element error started in the late seventies with the work by Babuška and Rheinboldt [8], where a residual-based error estimate is derived. Since then, many different types of a posteriori error estimates for the FEM have been introduced, such as error estimators obtained by solving local problems [1,41,83] or hierarchical [14], post-processed [128] and goal-oriented [13,22, 100] error estimators, just to mention a few. We refer to Verfürth [118], Ainsworth and Oden [3] or Grätsch and Bathe [73] for a review of these different a posteriori error estimation techniques. Concerning the error estimation of methods for solving random PDEs, a posteriori error estimators in the energy norm for the stochastic Galerkin finite element method (SG-FEM) are derived in $[24,58,59]$, where adaptive refinement algorithms are proposed for both stochastic and physical spaces. In the algorithm proposed in [59], the refined mesh is the same for all generalized polynomial chaos (gPC) modes, contrary to the one in [58] where the refinement procedure is applied independently for each mode. In [24], the adaptive procedure is driven by the two-sided estimates the authors obtained for the error reduction when the finite element subspace,
respectively the stochastic approximation space, is enriched. Concerning the stochastic collocation finite element method (SC-FEM), a priori error estimates are given in [7,20] but, to our knowledge, no a posteriori error estimate for the whole solution in suitable norms has been derived yet. Recently, a posteriori error estimates for a specific quantity of interest have been developed. Goal-oriented error estimates can be found in $[33,35,92]$ for the SG method and in [4] for the SC method.

We can distinguish two parts in this thesis. In the first part, which encompasses Chapters 1, 2,3 and 4, we consider PDEs with small uncertainties affecting the coefficients, the forcing term, the physical domain, the boundary conditions or combinations of them. The assumption of small uncertainties naturally leads to the choice of perturbation techniques for the approximation of the stochastic space. Indeed, if the level of uncertainty is small, then only few terms in the power series expansion of the solution with respect to a parameter $\varepsilon$ characterizing the amount of randomness of the problem will be needed to obtain an accurate approximation. With this technique, we are reduced to solve only deterministic problems whose solutions can be computed approximately with for instance the finite element method. The main goal of this thesis is then to derive error estimates that control the two sources of error: the stochastic error due to the truncation in the expansion of the solution and the spatial error coming from the finite element approximation of the continuous deterministic problems.

To have a general idea of the methodology, let us consider an abstract problem of the form: find $u(\cdot, \mathbf{Y}(\omega)) \in V$ such that almost surely

$$
\mathscr{A}(u, v ; \mathbf{Y}(\omega))=F(v ; \mathbf{Y}(\omega)) \quad \forall v \in V
$$

where $\mathbf{Y}$ is a random vector used to characterize the randomness in the input data, whose variability is controlled by a (small) parameter $\varepsilon$. Here, $V$ is a given Hilbert space, $\mathscr{A}$ is a bilinear form on $V \times V$ and $F$ is a linear functional on $V$, the latter two being parametrized by the random vector $\mathbf{Y}$. The solution $u$ of this problem also depends on $\mathbf{Y}$ and, adopting a perturbation approach, it is then expanded as

$$
u(\mathbf{x}, \mathbf{Y}(\omega))=u_{0}(\mathbf{x})+\varepsilon u_{1}(\mathbf{x}, \mathbf{Y}(\omega))+\varepsilon^{2} u_{2}(\mathbf{x}, \mathbf{Y}(\omega))+\ldots
$$

Considering a finite element space $V_{h} \subset V$, the first term in the expansion is approximated by $u_{0, h} \in V_{h}$, the solution of

$$
\mathscr{A}\left(u_{0, h}, v_{h} ; \mathbf{y}_{0}\right)=F\left(v_{h}, \mathbf{y}_{0}\right) \quad \forall v_{h} \in V_{h}
$$

with $\mathbf{y}_{0}=\mathbb{E}[\mathbf{Y}]$. Defining the residual for $u_{0, h}$ by

$$
R(\nu ; \mathbf{Y}(\omega)):=F(\nu ; \mathbf{Y}(\omega))-\mathscr{A}\left(u_{0, h}, \nu ; \mathbf{Y}(\omega)\right)
$$

the first step in the residual-based error estimation, that separates the two sources of error, is
then

$$
\mathscr{A}\left(u-u_{0, h}, v ; \mathbf{Y}(\omega)\right)=F(\nu ; \mathbf{Y}(\omega))-\mathscr{A}\left(u_{0, h}, v ; \mathbf{Y}(\omega)\right)=\mathrm{I}+\mathrm{II}
$$

with

$$
\begin{aligned}
\mathrm{I} & :=F\left(\nu ; \mathbf{y}_{0}\right)-\mathscr{A}\left(u_{0, h}, \nu ; \mathbf{y}_{0}\right)=R\left(\nu ; \mathbf{y}_{0}\right) \\
\text { II } & :=F(v ; \mathbf{Y}(\omega))-F\left(v ; \mathbf{y}_{0}\right)-\mathscr{A}\left(u_{0, h}, v ; \mathbf{Y}(\omega)\right)+\mathscr{A}\left(u_{0, h}, v ; \mathbf{y}_{0}\right)=R(\nu ; \mathbf{Y}(\omega))-R\left(\nu ; \mathbf{y}_{0}\right)
\end{aligned}
$$

The two terms can then be bounded separately. The first term I is nothing else than the residual for $u_{0, h}$ that can be bounded using a standard procedure as described by Verfürth in [118]. It yields an a posteriori error estimator that is localized on each element of the spatial mesh which can be used for mesh refinement. The second term is the one controlling the randomness. In this work, we will apply this methodology to a wide range of problems, as detailed in the thesis outline given below.

A different perspective is considered in the second part of this thesis, constituted of Chapter 5. Dropping the assumption of small uncertainty, and thus making perturbation techniques unsuitable, we use the stochastic collocation method to solve the random PDE. For the abstract problem considered above, this method, combined with the finite element method for the physical space discretization, consists in solving

$$
\mathscr{A}\left(u_{h}\left(\cdot, \mathbf{y}_{k}\right), v_{h} ; \mathbf{y}_{k}\right)=F\left(v_{h} ; \mathbf{y}_{k}\right) \quad \forall v_{h} \in V_{h}
$$

for a given set of collocation points $\mathbf{y}_{k}, k=1, \ldots, N_{c}$, in the stochastic space and building a global polynomial approximation

$$
u_{h, N_{c}}(\mathbf{x}, \mathbf{Y}(\omega))=\sum_{k=1}^{N_{c}} u_{h}\left(\mathbf{x}, \mathbf{y}_{k}\right) L_{k}(\mathbf{Y}(\omega))
$$

for suitable multivariate polynomials $L_{k}$. The goal is then to estimate the error due to this method when combined with the finite element method for the spatial discretization. We propose a residual-based a posteriori error estimate for an elliptic diffusion problem. It consists of two terms controlling each source of error, the SC and the FE error. The stochastic estimator is then used to drive an adaptive sparse grid algorithm.

The precise outline of this thesis is as follows.

## Thesis outline

We start in Chapter 1 with an in-depth analysis of a second order elliptic differential equation with random diffusion coefficient. We present the methodology we are using, namely a perturbation technique for the stochastic space approximation and the finite element method for the physical space discretization. We provide then a priori and a posteriori error analysis in various norms and for several approximations. Extension to some class of nonlinear problems
and a comparison in terms of computational costs with the stochastic collocation method are also provided. Many numerical experiments are presented to illustrate the theoretical findings.

The results are then extended in Chapter 2 where other sources of uncertainty are considered. More precisely, we consider first the case of random Neumann boundary conditions and then the combination of two uncertain inputs, the diffusion coefficient and the forcing term, described by two independent sets of random variables.

In Chapter 3, we consider nonlinear partial differential equations defined in random domains. Using the so-called domain mapping method, we use a random mapping to transform these equations into PDEs on a fixed reference domain with random coefficients. We start with the analysis of the one-dimensional steady-state viscous Burgers' equation in random intervals and consider then the more involved steady-state incompressible Navier-Stokes equations in random domains. We show the well-posedness of these problems, under suitable conditions on the mapping and the input data, and perform a posteriori error estimation for the finite element approximation of the first term in the expansion.

A time dependent parabolic problem is analysed in Chapter 4, considering the heat equation with random Robin boundary conditions. For the stochastic space, physical space and time discretizations, we use a perturbation technique, the finite element method and the (implicit) backward Euler scheme, respectively. We give an a posteriori error estimate for the first order approximation, which is here constituted of three parts controlling each source of error.

We conclude this thesis with an adaptive sparse grid algorithm for the stochastic collocation finite element method in Chapter 5. Considering again the diffusion model problem with random diffusion coefficient, that is assumed to depend affinely on a finite number of random variables, we derive an a posteriori error estimate for the total error that provides a guaranteed upper bound for the error. We propose then an algorithm that adaptively construct the multiindex set underlying the sparse grid and give numerical results to illustrate its performances.

Note: all the one-dimensional numerical experiments have been carried out using MATLAB Released R2012a, while the 2D numerical results have been obtained using either FreeFem++ 3.21 [78] or MATLAB.

# 1 Elliptic model problems with random diffusion coefficient 


#### Abstract

This chapter is mainly based on the paper [74] with respect to which we have done minor changes in the notation, essentially the distinction between a random vector $\mathbf{Y}: \Omega \rightarrow \Gamma \subset \mathbb{R}^{L}$ and a realization $\mathbf{y} \in \Gamma$. Moreover, we have added the following complements. First, a general statement of the model problem under consideration in Section 1.1. Additional numerical results are provided in Section 1.7. In particular, we present adaptive algorithms with nonuniform refinement which balances the two sources of error, namely the physical space discretization and the uncertainty. We give in Appendix some details about the derivation of the various deterministic problems for the first three terms in the expansion of the random solution, and state a precise link between each component of such terms and the derivatives of $u$ with respect to the stochastic space variable. Finally, a detailed proof of the upper and lower bounds of a certain error estimator and estimates of the interpolation constant closes this chapter.


## Introduction

In this chapter, we are focusing on PDEs with small uncertainties (for instance the linear model problem $-\operatorname{div}(a \nabla u)=f$ with $a=a_{0}+\varepsilon\left(a_{1} Y_{1}+\ldots+a_{L} Y_{L}\right)$ where $\varepsilon$ is small and $Y_{1}, \ldots, Y_{L}$ are random variables). Following a different path than Monte-Carlo type, stochastic Galerkin or stochastic collocation methods, we adopt a perturbation approach, see e.g. [37,82], which is appropriate for problems with small variability. We thus expand the stochastic solution $u$ as

$$
\begin{equation*}
u(\mathbf{x}, \omega)=u_{0}(\mathbf{x})+\varepsilon u_{1}(\mathbf{x}, \omega)+\mathscr{O}\left(\varepsilon^{2}\right) \tag{1.1}
\end{equation*}
$$

where $\varepsilon$ is a parameter controlling the magnitude of uncertainty in the input which is assumed to be small. Uncoupled problems can be derived to find the deterministic part $u_{0}$ and the stochastic one $u_{1}$ (and higher order terms), the error analysis being performed in various norms. The main goal is then to derive $a$ posteriori estimates for the error between the exact (random) solution $u$ and certain approximations to be defined. For instance, if we write $u_{0, h}$ for the FE approximation of $u_{0}$, then we will show that the error $u-u_{0, h}$ splits into two parts.

More precisely, we will derive an a posteriori error estimator $\eta$ composed of two deterministic computable quantities $\eta_{1}$ and $\eta_{2}$ such that the following upper bound for the error holds

$$
\left\|u-u_{0, h}\right\| \leq C \eta, \quad \eta=\left(\eta_{1}^{2}+\eta_{2}^{2}\right)^{\frac{1}{2}}
$$

with the norm $\|\cdot\|$ to be defined and where $C$ is a constant depending only on the domain $D$, the mesh and a (deterministic) ellipticity constant. Therefore, by solving only one deterministic problem we can obtain an upper bound of the error due to space discretization $\left(\eta_{1}\right)$ and the error due to uncertainty $\left(\eta_{2}\right)$. This estimator can then be used to determine a mesh size yielding comparable accuracy in $h$ and $\varepsilon$. The same kind of results can be obtained for $\left\|u-\left(u_{0, h}+\varepsilon u_{1, h}\right)\right\|$, yielding a better accuracy in $\varepsilon$, and then for higher order terms.

We mention that the a posteriori error estimator that we obtain for $u-u_{0, h}$ for the elliptic model problem (1.2) has similarities with the one derived in [26], although the context of this paper is quite different from the one considered here. In [26] the authors derive an adaptive finite element method (AFEM) for elliptic PDEs with discontinuous coefficients. The proposed algorithm takes into account the error due to FE approximation but also the effect of replacing the discontinuous input data by some piecewise polynomial approximation, which plays the same role as $a_{0}$ in our setting. More precisely, before applying a standard AFEM to the problem, the mesh is first refined so that the discontinuous input are approximated by piecewise polynomials with a prescribed accuracy. The specific form of the uncertain input we consider here, see (1.12), allows us to increase the accuracy in $\varepsilon$ by adding terms in the expansion (1.1) of $u$.

This chapter is organized as follows. The model problem, a second-order elliptic diffusion problem with homogeneous Dirichlet boundary conditions and random diffusion coefficient, is stated in Section 1.1. The diffusion coefficient is assumed, among others, to be expanded as a finite sum which depends on independent random variables with zero mean and unit variance. The methodology we are using to approximation the solution is given in Section 1.2. Error analysis in the $H_{0}^{1}$ and $L^{2}$ norms in the physical space, as well as goal-oriented error estimation, is performed in Section 1.3 where the exact (random) solution $u$ is approximated by the (deterministic) FE approximation of $u_{0}$. In Section 1.4, we consider the error between $u$ and the FE approximation of $u_{0}+\varepsilon u_{1}$, before giving a generalization for an approximation of arbitrary order in $\varepsilon$. The theory is then extended to nonlinear problems in Section 1.5. In Section 1.6, a comparison of the computational costs for the stochastic collocation method and the one presented here is performed. Section 1.7 is devoted to numerical examples used to illustrate and validate the theoretical results. Finally, a few complements are given in Appendix.

### 1.1 Problem statement

We start with a general and detailed description of the problem under consideration in this chapter, namely an elliptic diffusion PDE with random diffusion coefficient. In this description, we will make some distinctions in notation that will no longer be used in the next sections for ease of presentation.

## General problem statement

Let $D$ be a bounded polyhedral domain in $\mathbb{R}^{d}, d=1,2,3$, and $(\Omega, \mathscr{F}, P)$ a complete probability space, where $\Omega$ is the set of outcomes, $\mathscr{F} \subset 2^{\Omega}$ is the $\sigma$-algebra of events and $P: \mathscr{F} \rightarrow[0,1]$ is a probability measure. For any $p \in[1, \infty)$, let $L_{P}^{p}(\Omega)$ be the space of real-valued random variables $Y$ on $(\Omega, \mathscr{F}, P)$ that are $p$-integrable with respect to $P$, i.e. such that $\int_{\Omega}|Y(\omega)|^{p} d P(\omega)<\infty$. Moreover, if $Y \in L_{P}^{1}(\Omega)$, we denote its expected value (or mean) by $\mathbb{E}[Y]=\int_{\Omega} Y(\omega) d P(\omega)$. The following problem is considered.

Find $u: D \times \Omega \rightarrow \mathbb{R}$ such that $P$-almost everywhere in $\Omega$ (in other words almost surely in $\Omega$ ):

$$
\left\{\begin{array}{rlll}
-\operatorname{div}(a(\mathbf{x}, \omega) \nabla u(\mathbf{x}, \omega)) & = & f(\mathbf{x}) & \mathbf{x} \in D  \tag{1.2}\\
u(\mathbf{x}, \omega) & =0 & \mathbf{x} \in \partial D
\end{array}\right.
$$

where $a$ is a random field on $(\Omega, \mathscr{F}, P)$ over $L^{\infty}(D)$. For simplicity, the right-hand side $f$ is assumed to be deterministic, $f \in L^{2}(D)$, but the case of stochastic forcing term could be considered as well adding no real difficulty, see Chapter 2. Note that the divergence and gradient operators apply only on $\mathbf{x}$, the physical space variable. Let $H_{0}^{1}(D)$ be endowed with the following norm

$$
\|v\|_{H_{0}^{1}(D)}:=\|\nabla v\|_{L^{2}(D)}=\left(\int_{D}|\nabla \nu|^{2}\right)^{\frac{1}{2}} .
$$

The problem (1.2) can be written in weak form as:
find $u \in L_{P}^{2}(\Omega) \otimes H_{0}^{1}(D)$ such that

$$
\begin{equation*}
\mathbb{E}\left[\int_{D} a \nabla u \cdot \nabla v d \mathbf{x}\right]=\mathbb{E}\left[\int_{D} f v d \mathbf{x}\right] \quad \forall v \in L_{P}^{2}(\Omega) \otimes H_{0}^{1}(D) . \tag{1.3}
\end{equation*}
$$

Since the tensor product space $L_{P}^{2}(\Omega) \otimes H_{0}^{1}(D)$ is isomorphic (see for instance [10]) to the Bochner space

$$
\begin{equation*}
L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right):=\left\{v: \Omega \rightarrow H_{0}^{1}(D) \mid v \text { is strongly measurable and }\|v\|_{L_{p}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}<\infty\right\} \tag{1.4}
\end{equation*}
$$

where

$$
\|v\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}^{2}:=\int_{\Omega}\|\nabla v(\cdot, \omega)\|_{L^{2}(D)}^{2} d P(\omega)=\mathbb{E}\left[\|\nabla v\|_{L^{2}(D)}^{2}\right]
$$

we can see the weak solution $u$ of problem (1.2) as a function $u: \Omega \rightarrow H_{0}^{1}(D)$. The correspond-
ing pointwise weak formulation, equivalent to (1.3), is then given by:
find $u(\cdot, w) \in H_{0}^{1}(D)$ such that

$$
\begin{equation*}
\int_{D} a(\mathbf{x}, \omega) \nabla u(\mathbf{x}, \omega) \cdot \nabla v(\mathbf{x}) d \mathbf{x}=\int_{D} f(\mathbf{x}) v(\mathbf{x}) d \mathbf{x} \quad \forall v \in H_{0}^{1}(D), P \text {-a.e. in } \Omega . \tag{1.5}
\end{equation*}
$$

If the diffusion coefficient $a$ is (uniformly) bounded from below and from above, namely

$$
\begin{equation*}
\exists 0<a_{\min } \leq a_{\max }<\infty: \quad P\left(\omega \in \Omega: a_{\min } \leq a(\mathbf{x}, \omega) \leq a_{\max } \forall \mathbf{x} \in \bar{D}\right)=1 \tag{1.6}
\end{equation*}
$$

then we can show, by a straightforward application of Lax-Milgram's Lemma, that problem (1.5) is well-posed. More precisely, there exists a unique solution $u \in L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ which satisfies the a priori estimate

$$
\|u\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)} \leq \frac{C_{P}}{a_{\min }}\|f\|_{L^{2}(D)}
$$

with $C_{P}=C_{P}(D)$ the Poincaré constant.

Remark 1.1.1. With the above assumptions, the solution belongs to $L_{P}^{k}\left(\Omega ; H_{0}^{1}(D)\right)$ for any $k \in[1, \infty]$. This is also true in the more general case $f \in L_{P}^{k p}\left(\Omega ; H_{0}^{1}(D)\right)$ and $a(x, \omega) \geq a_{\text {min }}(\omega)>0$ a.e. in $D$ and a.s. in $\Omega$ with $\frac{1}{a_{\min }} \in L_{P}^{k q}(\Omega)$, where $\frac{1}{p}+\frac{1}{q}=1$ (see [7]).

We further assume that the random coefficient is well approximated by the finite expansion

$$
\begin{equation*}
a(\mathbf{x}, \omega) \approx a_{L}(\mathbf{x}, \omega)=a_{0}(\mathbf{x})+\varepsilon \sum_{j=1}^{L} a_{j}(\mathbf{x}) Y_{j}(\omega) \quad \text { with } \quad a_{0}(\mathbf{x})=\mathbb{E}[a(\mathbf{x}, \cdot)], \tag{1.7}
\end{equation*}
$$

where $\left\{Y_{j}\right\}_{j=1}^{L}$ are independent random variables with zero mean and unit variance.

Remark 1.1.2. The characterization (1.7) of the random input can be achieved using for instance a truncated Karhunen-Loève type expansion (see [87, 88]) if the mean and the twopoint correlation (or equivalently the covariance) of a is known. In this case, the functions $a_{j}$, $j=1, \ldots, L$, in (1.7) write $a_{j}(\mathbf{x})=\sqrt{\lambda_{j}} \varphi_{j}(\mathbf{x})$ with $\left\{\lambda_{j}, \varphi_{j}\right\}$ the eigenpairs of the (compact and self-adjoint) integral operator associated with the covariance kernel $V: D \times D \rightarrow \mathbb{R}$ given by

$$
V\left(\mathbf{x}, \mathbf{x}^{\prime}\right):=\frac{1}{\varepsilon^{2}} \mathbb{E}\left[\left(a(\mathbf{x}, \omega)-a_{0}(\mathbf{x})\right)\left(a\left(\mathbf{x}^{\prime}, \omega\right)-a_{0}\left(\mathbf{x}^{\prime}\right)\right)\right] .
$$

Notice that, in general, the family of random variables appearing in the KL expansion of an arbitrary random field a are only uncorrelated (see [111]), but not necessarily independent.

The problem (1.2) is then approximated by:
find $u_{L}: D \times \Omega \rightarrow \mathbb{R}$ such that $P$-a.e. in $\Omega$ the following equation holds

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{L}(\mathbf{x}, \omega) \nabla u_{L}(\mathbf{x}, \omega)\right) & =f(\mathbf{x}) & & \mathbf{x} \in D  \tag{1.8}\\
u_{L}(\mathbf{x}, \omega) & =0 & \mathbf{x} \in \partial D
\end{array}\right.
$$

which admits a unique weak solution $u_{L} \in L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ under the assumption

$$
\exists 0<a_{L, \min } \leq a_{L, \max }<\infty: \quad P\left(\omega \in \Omega: a_{L, \min } \leq a_{L}(\mathbf{x}, \omega) \leq a_{L, \max }, \forall \mathbf{x} \in \bar{D}\right)=1
$$

The stochasticity of the problem (1.8) for $u_{L}$ can therefore be parametrized by the random vector $\mathbf{Y}=\left(Y_{1}, \ldots, Y_{L}\right)$. Indeed, with the definition of $a_{L}$ given in (1.7) we have $a_{L}(\mathbf{x}, \omega)=$ $\tilde{a}_{L}\left(\mathbf{x}, Y_{1}(\omega), \ldots, Y_{L}(\omega)\right)$ and thus $u_{L}(\mathbf{x}, \omega)=\tilde{u}_{L}\left(\mathbf{x}, Y_{1}(\omega), \ldots, Y_{L}(\omega)\right)$ thanks to the Doob-Dynkin Lemma (see [6, p.6] for instance). We can therefore derive a parametric deterministic weak formulation of (1.8). Let $\Gamma=\Gamma_{1} \times \Gamma_{2} \times \ldots \times \Gamma_{L}$ where $\Gamma_{j}$ denotes the bounded image in $\mathbb{R}$ of the random variable $Y_{j}$, i.e. $\Gamma_{j}:=Y_{j}(\Omega)$, for $j=1, \ldots, L$. Moreover, let $\rho_{j}: \Gamma_{j} \rightarrow \mathbb{R}^{+}$denote the probability density function of $Y_{j}, j=1, \ldots, L$. Thanks to the independence of the random variables, the joint density function $\rho: \Gamma \rightarrow \mathbb{R}^{+}$of the random vector $\mathbf{Y}$ factorizes as $\rho(\mathbf{y})=$ $\prod_{j=1}^{L} \rho_{j}\left(y_{j}\right)$ for all $\mathbf{y}=\left(y_{1}, \ldots, y_{L}\right) \in \Gamma$. We can thus replace the probability space $(\Omega, \mathscr{F}, P)$ by its image ( $\Gamma, B(\Gamma), \rho(\mathbf{y}) d \mathbf{y})$, where $B(\Gamma)$ denotes the Borel $\sigma$-algebra defined on $\Gamma$ and $\rho(\mathbf{y}) d \mathbf{y}$ the probability measure of $\mathbf{Y}$. For any measurable function $\tilde{g}_{L}: \Gamma \rightarrow \mathbb{R}$ defined on $(\Gamma, B(\Gamma), \rho(\mathbf{y}) d \mathbf{y})$, the expectation of the random variable $g_{L}=\tilde{g}_{L} \circ \mathbf{Y}: \Omega \rightarrow \mathbb{R}$ is then given by

$$
\mathbb{E}\left[g_{L}\right]=\int_{\Omega} g_{L}(\omega) d P(\omega)=\int_{\Omega} \tilde{g}_{L}(\mathbf{Y}(\omega)) d P(\omega)=\int_{\Gamma} \tilde{g}_{L}(\mathbf{y}) \rho(\mathbf{y}) d \mathbf{y}
$$

Remark 1.1.3. The error analysis for $u-u_{0}$ with $u_{0}$ the first term in the expansion, see (1.1), is exactly the same as the one performed in Section 1.3 if the random variables are assumed uncorrelated instead of independent, i.e. such that $\mathbb{E}\left[Y_{i} Y_{j}\right]=\mathbb{E}\left[Y_{i}\right] \mathbb{E}\left[Y_{j}\right]$ for any $i, j=1, \ldots, L$ with $i \neq j$. For the higher order approximations, however, few changes have to be made to the analysis given in Section 1.4. Moreover, the definitions given above are not restricted to continuous random variables but also hold for discrete random variables. In such a case, we consider a generalized probability density function defined via Dirac delta functions. For instance, the density function of a random variable $Y_{j}$ taking value $\pm 1$ with probability $\frac{1}{2}$ would be

$$
\rho_{j}\left(y_{j}\right)=\frac{1}{2}\left(\delta\left(y_{j}+1\right)+\delta\left(y_{j}-1\right)\right)
$$

Such random variables will be considered in the numerical results of Section 1.7.

The (parametric, pointwise) weak formulation of problem (1.8) reads:
find $\tilde{u}_{L}: \Gamma \rightarrow H_{0}^{1}(D)$ such that

$$
\begin{equation*}
\int_{D} \tilde{a}_{L}(\mathbf{x}, \mathbf{y}) \nabla \tilde{u}_{L}(\mathbf{x}, \mathbf{y}) \cdot \nabla v(\mathbf{x}) d \mathbf{x}=\int_{D} f(\mathbf{x}) v(\mathbf{x}) d \mathbf{x} \quad \forall v \in H_{0}^{1}(D), \rho \text {-a.e. in } \Gamma \tag{1.9}
\end{equation*}
$$

where $\tilde{a}_{L}(\mathbf{x}, \mathbf{y})=a_{0}(\mathbf{x})+\varepsilon \sum_{j=1}^{L} a_{j}(\mathbf{x}) y_{j}$. Thanks again to Lax-Milgram's lemma, we know that there exists a unique solution $\tilde{u}_{L} \in L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)$ of problem (1.9) which satisfies

$$
\left\|\tilde{u}_{L}\right\|_{L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)} \leq \frac{C_{P}}{a_{\min }}\|f\|_{L^{2}(D)}
$$

where similarly to (1.4) we define

$$
\begin{equation*}
L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right):=\left\{v: \Gamma \rightarrow H_{0}^{1}(D) \mid v \text { is strongly measurable and }\|v\|_{L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)}<\infty\right\} \tag{1.10}
\end{equation*}
$$

with

$$
\|\nu\|_{L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)}^{2}:=\int_{\Gamma}\|\nabla v(\cdot, \mathbf{y})\|_{L^{2}(D)}^{2} \rho(\mathbf{y}) d \mathbf{y} .
$$

Notice that the weak solution $u_{L}$ of problem (1.8) and the solution $\tilde{u}_{L}$ of problem (1.9) are related by

$$
u_{L}(x, \omega)=\tilde{u}_{L}\left(x, Y_{1}(\omega), \ldots, Y_{L}(\omega)\right) \quad \text { a.s. in } \Omega
$$

and we have

$$
\left\|u_{L}\right\|_{L_{p}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}=\left\|\tilde{u}_{L}\right\|_{L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)} .
$$

For the sake of presentation $\tilde{a}_{L}$ and $\tilde{u}_{L}$ will be denoted again $a_{L}$ and $u_{L}$, respectively, i.e. we write $a_{L}(\mathbf{x}, \omega)=a_{L}\left(\mathbf{x}, Y_{1}(\omega), \ldots, Y_{L}(\omega)\right)$ and $u_{L}(\mathbf{x}, \omega)=u_{L}\left(\mathbf{x}, Y_{1}(\omega), \ldots, Y_{L}(\omega)\right)$, when no ambiguity arises. Moreover, the goal here is not to analyse the error committed when replacing $a$ by $a_{L}$, i.e. when the random input is approximated via $L$ random variables. Therefore, we assume from now on that $a=a_{L}$, i.e. $u=u_{L}$. We mention that a complete analysis of the (strong, weak) error $u-u_{L}$ can be found in [44].

## Specific problem statement

We give now a short statement of the problem that will be analysed in the subsequent sections, indicating only the necessary assumptions and using the shorthand notation described above. We consider the following problem.

Find $u: D \times \Omega \rightarrow \mathbb{R}$ such that a.s. in $\Omega$ :

$$
\left\{\begin{array}{rll}
-\operatorname{div}(a(\mathbf{x}, \omega) \nabla u(\mathbf{x}, \omega)) & =f(\mathbf{x}) &  \tag{1.11}\\
\mathbf{x} \in D \\
u(\mathbf{x}, \omega) & =0 & \\
\mathbf{x} \in \partial D
\end{array}\right.
$$

where $f \in L^{2}(D)$ is deterministic and $a$ is a random field on $(\Omega, \mathscr{F}, P)$ over $L^{\infty}(D)$ which satisfies the following assumptions (see $[6,7,10]$ for instance) that ensure, among others, the well-posedness of the problem:
(A1) coercivity and continuity: $a$ is bounded and uniformly coercive, i.e. there exist two real
constants $0<a_{\min } \leq a_{\max }<\infty$ such that

$$
P\left(\omega \in \Omega: a_{\min } \leq a(\mathbf{x}, \omega) \leq a_{\max }, \forall \mathbf{x} \in \bar{D}\right)=1
$$

(A2) finite dimensional noise: $a$ is parametrized by $L$ mutually independent random variables $a(\mathbf{x}, \omega)=a\left(\mathbf{x}, Y_{1}(\omega), Y_{2}(\omega), \ldots, Y_{L}(\omega)\right)$. More precisely, we assume that $a$ can be expanded as

$$
\begin{equation*}
a(\mathbf{x}, \omega)=a_{0}(\mathbf{x})+\varepsilon \sum_{j=1}^{L} a_{j}(\mathbf{x}) Y_{j}(\omega), \tag{1.12}
\end{equation*}
$$

where the $\left\{Y_{j}\right\}_{j=1}^{L}$ are independent random variables with zero mean and unit variance, $a_{j} \in W^{1, \infty}(D)$ for $j=0, \ldots, L$ and $\varepsilon \in\left[0, \varepsilon_{\text {max }}\right]$ with $\varepsilon_{\max }$ the maximum value such that property (A1) is satisfied. The functions $a_{j}, j=0,1, \ldots, L$, and the random variables $Y_{j}$, $j=1, \ldots, L$, are assumed to be independent of $\varepsilon$.

Notice that assuming $a_{j} \in L^{\infty}(D)$ for $j=0,1, \ldots, L$ is enough to ensure the well-posedness of the problem. We impose here more regularity in order to avoid difficulties that are beyond the scope of this work. We refer to [23] for a derivation of a posteriori error estimation in the case of discontinuous coefficients. Moreover, as a consequence of assumption (A1), the random variables $Y_{j}, j=1, \ldots, L$, have to be bounded almost surely. In particular, they have finite moment of any order. Finally, from assumption (A2) it follows that the mean and variance of $a$ are given by $\mathbb{E}[a](\mathbf{x})=a_{0}(\mathbf{x})$ and $\operatorname{Var}[a](\mathbf{x})=\varepsilon^{2} \sum_{j=1}^{L} a_{j}^{2}(\mathbf{x})$, respectively. Therefore, for fixed functions $a_{j}$, we can modify the variance of $a$ by changing the value of $\varepsilon$. From assumption (A2), the solution $u$ is a function of the random variables $Y_{j}$, i.e. $u(\mathbf{x}, \omega)=u\left(\mathbf{x}, Y_{1}(\omega), \ldots, Y_{L}(\omega)\right)$. Replacing $(\Omega, \mathscr{F}, P)$ by $(\Gamma, B(\Gamma), \rho(\mathbf{y}) d \mathbf{y})$, the stochastic elliptic boundary value problem (1.11) can equivalently be written in the following deterministic parametric form:
find $u: D \times \Gamma \rightarrow \mathbb{R}$ such that $\rho$-a.e. in $\Gamma$ we have

$$
\left\{\begin{array}{rlll}
-\operatorname{div}(a(\mathbf{x}, \mathbf{y}) \nabla u(\mathbf{x}, \mathbf{y})) & =f(\mathbf{x}) & \mathbf{x} \in D  \tag{1.13}\\
u(\mathbf{x}, \mathbf{y}) & =0 & & \mathbf{x} \in \partial D .
\end{array}\right.
$$

The (parametric, pointwise) weak form of problem (1.13) then reads:
find $u(\cdot, \mathbf{y}) \in H_{0}^{1}(D)$ such that

$$
\begin{equation*}
\mathscr{A}(u(\cdot, \mathbf{y}), v ; \mathbf{y})=\mathscr{F}(v) \quad \forall v \in H_{0}^{1}(D), \rho \text {-a.e. in } \Gamma . \tag{1.14}
\end{equation*}
$$

where

$$
\begin{align*}
\mathscr{A}(u(\cdot, \mathbf{y}), v ; \mathbf{y}) & =\int_{D} a(\mathbf{x}, \mathbf{y}) \nabla u(\mathbf{x}, \mathbf{y}) \cdot \nabla v(\mathbf{x}) d \mathbf{x},  \tag{1.15}\\
\mathscr{F}(v) & =\int_{D} f(\mathbf{x}) v(\mathbf{x}) d \mathbf{x} . \tag{1.16}
\end{align*}
$$

Again, thanks to Lax-Milgram's lemma the coercivity and continuity assumptions on $a$ ensure the well-posedness of problem (1.14), namely there exists a unique solution $u \in L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)$. Indeed, since $a$ is bounded from below and above almost surely, the bilinear form $\mathscr{A}$ is continuous and coercive with constant of continuity and coercivity given respectively by $a_{\max }$ and $a_{\min }$. Furthermore, the linear (deterministic) functional $\mathscr{F}$ is continuous, with constant of continuity equal to $C_{P}\|f\|_{L^{2}(D)}$, where $C_{P}$ denotes the constant in the Poincaré inequality. Therefore, the solution $u$ of problem (1.14) satisfies

$$
\begin{equation*}
\|\nabla u(\cdot, \mathbf{y})\|_{L^{2}(D)} \leq \frac{C_{P}}{a_{\min }}\|f\|_{L^{2}(D)} \quad \rho \text {-a.e. in } \Gamma . \tag{1.17}
\end{equation*}
$$

Notice that the weak solution of problem (1.11) is then given by $u(\cdot, \mathbf{Y}(\omega))$ with $u$ the parametric solution of problem (1.14) and it satisfies

$$
\begin{equation*}
\|\nabla u(\cdot, \mathbf{Y}(\omega))\|_{L^{2}(D)} \leq \frac{C_{P}}{a_{m i n}}\|f\|_{L^{2}(D)} \quad \text { a.s. in } \Omega \tag{1.18}
\end{equation*}
$$

Moreover, it has been proved (see for instance [7]) that solution $u=u(\mathbf{x}, \mathbf{y})$ of (1.14) is analytic with respect to each variable $y_{j}, j=1, \ldots, L$.

For ease of presentation, the dependence of the random variables $Y_{j}$ with respect to $\omega \in \Omega$ will not necessarily be indicated in the subsequent analysis.

### 1.2 Methodology

In this section, we present the method we use to approximate the random (weak) solution $u$ of problem (1.11). We use first a perturbation technique for the stochastic space approximation, yielding a collection of deterministic problems. The physical space approximation of each problem is then performed using the finite element method. More precisely, we assume from now on that $\varepsilon$ in (1.12) is small enough that (A2) holds and expand the solution $u=u(\mathbf{x}, \mathbf{Y}(\omega))$ with respect to $\varepsilon$ up to a certain order $N \in \mathbb{N}$

$$
\begin{equation*}
u(\mathbf{x}, \mathbf{Y}(\omega))=u_{0}(\mathbf{x})+\varepsilon u_{1}(\mathbf{x}, \mathbf{Y}(\omega))+\ldots+\varepsilon^{N} u_{N}(\mathbf{x}, \mathbf{Y}(\omega))+\mathscr{O}\left(\varepsilon^{N+1}\right) \tag{1.19}
\end{equation*}
$$

Inserting the latter expansion into (1.11) with $a$ defined in (1.12) and keeping the $\mathscr{O}(1)$ term with respect to $\varepsilon$ yields the problem:
find $u_{0}: D \rightarrow \mathbb{R}$ such that

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{0}(\mathbf{x}) \nabla u_{0}(\mathbf{x})\right) & =f(\mathbf{x}) & \mathbf{x} \in D  \tag{1.20}\\
u_{0}(\mathbf{x}) & =0 & \mathbf{x} \in \partial D .
\end{array}\right.
$$

Then, writing $u_{1}(\mathbf{x}, \mathbf{Y}(\omega))=\sum_{j=1}^{L} U_{j}(\mathbf{x}) Y_{j}(\omega)$ and keeping the $\mathscr{O}(\varepsilon)$ terms in (1.11) yields the $L$ problems:
find $U_{j}: D \rightarrow \mathbb{R}$ such that

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{j}(\mathbf{x}) \nabla u_{0}(\mathbf{x})+a_{0}(\mathbf{x}) \nabla U_{j}(\mathbf{x})\right) & = & 0 & \mathbf{x} \in D  \tag{1.21}\\
U_{j}(\mathbf{x}) & = & 0 & \mathbf{x} \in \partial D
\end{array} \quad j=1, \ldots, L,\right.
$$

in which the solution $u_{0}$ of problem (1.20) is needed. Notice that for $j=1, \ldots, L$, the function $U_{j}$ is related to $\frac{\partial u\left(\mathbf{x}, \mathbf{y}_{0}\right)}{\partial y_{j}}$ with $\mathbf{y}_{0}=\mathbb{E}[\mathbf{Y}]=\mathbf{0}$. Similarly, we can use the solutions $U_{j}, j=1, \ldots, L$, of problem (1.21) to compute the deterministic part of the next term in the expansion (1.19), which in turn is related to the second derivatives $\frac{\partial^{2} u\left(\mathbf{x}, \mathbf{y}_{0}\right)}{\partial y_{k} y_{j}}, j, k=1, \ldots, L$. Indeed, if we write $u_{2}(\mathbf{x}, \mathbf{Y}(\omega))=\sum_{j, k=1}^{L} U_{j k}(\mathbf{x}) Y_{j}(\omega) Y_{k}(\omega)$, keeping the $\mathscr{O}\left(\varepsilon^{2}\right)$ terms in (1.11), we get the $L^{2}$ problems:
find $U_{j k}: D \rightarrow \mathbb{R}$ such that

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{j}(\mathbf{x}) \nabla U_{k}(\mathbf{x})+a_{0}(\mathbf{x}) \nabla U_{j k}(\mathbf{x})\right) & = & 0 & \mathbf{x} \in D  \tag{1.22}\\
U_{j k}(\mathbf{x}) & = & 0 & \mathbf{x} \in \partial D
\end{array} \quad j, k=1, \ldots, L .\right.
$$

More details about the derivation of problems (1.20), (1.21) and (1.22) are given in Appendix 1.A.

Remark 1.2.1. We will prove in the sections 1.3, 1.4.1 and 1.4.2 that

$$
u-u_{0}=\mathscr{O}(\varepsilon), \quad u-\left(u_{0}+\varepsilon u_{1}\right)=\mathscr{O}\left(\varepsilon^{2}\right) \text { and } u-\left(u_{0}+\varepsilon u_{1}+\varepsilon^{2} u_{2}\right)=\mathscr{O}\left(\varepsilon^{3}\right) .
$$

The solution to the deterministic problems (1.20), (1.21) and (1.22) can be approximated using for instance the finite element method. For any $h>0$, let $\mathscr{T}_{h}$ be a family of partitions of $D$ into $d$-simplices (intervals, triangles, tetrahedra) $K$ of diameter $h_{K} \leq h$. Unless otherwise stated, we will always consider shape regular (see [49]) meshes of $D$, i.e. decompositions such that there exists a constant $c>0$ satisfying

$$
\begin{equation*}
\frac{h_{K}}{\rho_{K}} \leq c \quad \forall K \in \mathscr{T}_{h}, \forall h>0 \tag{1.23}
\end{equation*}
$$

where $\rho_{K}=\sup \{\operatorname{diam}(B): B$ is a ball contained in $K\}$. The condition (1.23) is equivalent to a minimal angle condition, namely that there exists a constant $\alpha_{0}$ such that $\alpha_{K} \geq \alpha_{0}>0$ for all $K \in \mathscr{T}_{h}$ with $\alpha_{K}$ the smallest angle of $K$. Let $V_{h} \subset H_{0}^{1}(D)$ be the space of continuous, piecewise linear finite element functions associated to $\mathscr{T}_{h}$ that vanish on $\partial D$, that is

$$
V_{h}:=\left\{v_{h} \in C^{0}(\bar{D}): v_{h \mid K} \in \mathbb{P}_{1} \forall K \in \mathscr{T}_{h}\right\} \cap H_{0}^{1}(D),
$$

where $\mathbb{P}_{1}$ is the set of polynomials of degree less than or equal to 1 .
In the derivation of a priori and a posteriori error estimates, we will need an interpolation operator which maps $H_{0}^{1}(D)$ to $V_{h}$, along with interpolation error bounds. We distinguish the
cases $d=1$ and $d=2,3$. For the one-dimensional case, any function of $H_{0}^{1}(D)$ is continuous thanks to the Sobolev embedding theorem. Therefore, the Lagrange interpolant operator $r_{h}: C^{0}(\bar{D}) \rightarrow V_{h}$, which requires point evaluations, is well-defined and satisfies the following error bounds: there exists a constant $C>0$ such that $\forall h>0, \forall K \in \mathscr{T}_{h}$ and all $v \in H_{0}^{1}(D)$ we have

$$
\begin{equation*}
\left\|v-r_{h} v\right\|_{L^{2}(K)} \leq C h_{K}\left\|v^{\prime}\right\|_{L^{2}(K)} \tag{1.24}
\end{equation*}
$$

and for all $v \in H^{2}(D)$

$$
\left\|\nu-r_{h} \nu\right\|_{L^{2}(K)}+h_{K}\left\|\nu^{\prime}-\left(r_{h} \nu\right)^{\prime}\right\|_{L^{2}(K)} \leq h_{K}^{2}\left\|v^{\prime \prime}\right\|_{L^{2}(K)}
$$

For the case $d=2,3$, the functions of $H^{2}(D)$ are continuous and we have the following error bound (see $[31,49]$ for instance) based on the Bramble-Hilbert lemma: there exists a constant $C>0$ such that $\forall h>0, \forall K \in \mathscr{T}_{h}$ and all $v \in H^{2}(K)$ we have

$$
\begin{equation*}
\left\|v-r_{h} v\right\|_{L^{2}(K)}+h_{K}\left\|\nabla\left(v-r_{h} v\right)\right\|_{L^{2}(K)} \leq C h_{K}^{2}|v|_{H^{2}(K)} \tag{1.25}
\end{equation*}
$$

In general however, such regularity might not be reached by the solution of problem (1.14), since we are seeking a solution in $H_{0}^{1}(D)$ in the physical space. In that case, we will use the Clément interpolant [50] operator $\mathscr{I}_{h}: H^{1}(D) \rightarrow V_{h}$ which satisfies the following interpolation results: there exists a constant $C>0$ such that $\forall h>0, \forall K, e \in \mathscr{T}_{h}$ and all $v \in H^{1}(D)$ we have

$$
\begin{align*}
& \left\|v-\mathscr{I}_{h} v\right\|_{L^{2}(K)} \leq C h_{K}|v|_{H^{1}(N(K))},  \tag{1.26}\\
& \left\|\nabla\left(v-\mathscr{I}_{h} v\right)\right\|_{L^{2}(K)} \leq C|v|_{H^{1}(N(K))} \tag{1.27}
\end{align*}
$$

and

$$
\begin{equation*}
\left\|v-\mathscr{I}_{h} v\right\|_{L^{2}(e)} \leq C h_{e}^{\frac{1}{2}}|v|_{H^{1}\left(N\left(K_{e}\right)\right)}, \tag{1.28}
\end{equation*}
$$

where, for an internal edge $e, K_{e}$ is the union of the two elements touching $e$ and $N(K)$ (respectively $N\left(K_{e}\right)$ ) denotes the patch of elements associated to $K$ (respectively $K_{e}$ ). Notice that the constant $C$ in (1.26), (1.27) and (1.28) depends on the constant in (1.23) characterizing the mesh aspect ratio.

We will now derive a priori and a posteriori error estimates in various norms, the error being the difference between the exact solution and a certain approximate solution to be defined. We first start by giving error estimates between the exact solution $u$ and $u_{0, h}$, the FE approximation of $u_{0}$. Our goal is to decompose the error into two parts, the error due to the finite element approximation $(h)$ and the error due to the uncertainty $(\varepsilon)$.

### 1.3 Error analysis for the first order approximation

We consider $u$ the (weak) solution of (1.11) and $u_{0}$ that of (1.20), i.e. the case $N=0$ in the expansion (1.19). The error due to the stochastic truncation is of order $\varepsilon$. Indeed, for any
$v \in H_{0}^{1}(D)$ and a.s. in $\Omega$ we have

$$
\begin{equation*}
\int_{D} a \nabla\left(u(\cdot, \mathbf{Y}(\omega))-u_{0}\right) \cdot \nabla v=\int_{D} f v-\int_{D} a(\cdot, \mathbf{Y}(\omega)) \nabla u_{0} \cdot \nabla v=-\varepsilon \sum_{j=1}^{L} Y_{j}(\omega) \int_{D} a_{j} \nabla u_{0} \cdot \nabla v \tag{1.29}
\end{equation*}
$$

Using the FEM, the unknown solution $u_{0}$ of problem (1.20) is approximated by $u_{0, h}$, the solution of:

$$
\begin{equation*}
\text { find } u_{0, h} \in V_{h} \text { such that } \int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v_{h}=\int_{D} f v_{h} \quad \forall v_{h} \in V_{h} \text {. } \tag{1.30}
\end{equation*}
$$

In what follows, we will derive a priori and a posteriori error estimates for $u-u_{0, h}$ in various norms. In particular, the a posteriori error estimators, which are computable quantities, yield useful information about the two sources of error by computing only one deterministic problem.

### 1.3.1 A priori error analysis

This section is devoted to a priori error estimation for the strong and weak errors, which gives information on the asymptotic behaviour of the error. In particular, we will show that the order of the error of the mean in $\varepsilon$ is twice the order of the strong error, while the order of the error in $h$ is the same for both. Sections 1.3.2, 1.3.2 and 1.3.2 are instead devoted to a posteriori error estimates in different norms.

## Strong error estimate

Let us first give error estimates on the strong error, i.e. on the error between $u$ and $u_{0, h}$ in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm. Our goal is to prove that there exists a constant $C>0$ independent of $h$ and $\varepsilon$ such that

$$
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq C(h+\varepsilon)
$$

Proposition 1.3.1. Let $u$ and $u_{0}$ be the (weak) solutions of problems (1.11) and (1.20), respectively, and let $u_{0, h}$ be the solution of problem (1.30). If $u_{0} \in H^{2}(D)$, then we have the a priori error estimate

$$
\begin{equation*}
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \sqrt{2}\left[\frac{a_{0, \text { max }}}{a_{0, \text { min }}} C^{2} h^{2}\left|u_{0}\right|_{H^{2}(D)}^{2}+L \frac{\varepsilon^{2} C_{P}^{2}}{a_{0, \text { min }}^{2} a_{\text {min }}^{2}}\|f\|_{L^{2}(D)}^{2} \sum_{j=1}^{L}\left\|a_{j}^{2}\right\|_{L^{\infty}(D)}\right]^{\frac{1}{2}} \tag{1.31}
\end{equation*}
$$

where $C>0$ is the constant that appears in (1.25). Moreover, if we assume that for a fixed value $\alpha>\frac{1}{2}$, there exists a constant $M_{\alpha}$ such that for any $L$ we have $\sum_{j=1}^{L}\left\|a_{j}^{2}\right\|_{L^{\infty}(D)} j^{2 \alpha} \leq M_{\alpha}$, then we
also have

$$
\begin{equation*}
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \sqrt{2}\left[\frac{a_{0, \max }}{a_{0, \min }} C^{2} h^{2}\left|u_{0}\right|_{H^{2}(D)}^{2}+M_{\alpha} \frac{\varepsilon^{2} C_{P}^{2}}{a_{0, \text { min }}^{2} a_{\min }^{2}}\|f\|_{L^{2}(D)}^{2} \sum_{j=1}^{\infty} j^{-2 \alpha}\right]^{\frac{1}{2}} . \tag{1.32}
\end{equation*}
$$

Remark 1.3.2. The a priori error estimate (1.31) blows up when $L$ tends to infinity since the second part of the estimate depends linearly on L. If we add a constraint on the functions $a_{j}, j=1, \ldots, L$, for instance that $a_{j}$ decays as $j^{-\beta}$ with $\beta>\alpha+\frac{1}{2}$, then (1.32) holds with $M_{\alpha}$ independent of $L$.

Proof. Using the fact that almost surely it holds

$$
\int_{D} a_{0} \nabla u_{0} \cdot \nabla v=\int_{D} f v=\int_{D} a \nabla u \cdot \nabla v \quad \forall v \in H_{0}^{1}(D),
$$

we have for any $v \in V$

$$
\begin{align*}
\int_{D} a_{0} \nabla\left(u-u_{0, h}\right) \cdot \nabla v & =\int_{D} a_{0} \nabla\left(u-u_{0}\right) \cdot \nabla v+\int_{D} a_{0} \nabla\left(u_{0}-u_{0, h}\right) \cdot \nabla v  \tag{1.33}\\
& =-\int_{D}\left(a-a_{0}\right) \nabla u \cdot \nabla v+\int_{D} a_{0} \nabla\left(u_{0}-u_{0, h}\right) \cdot \nabla v \\
& \leq\left[\left(\int_{D} \frac{\left(a_{0}-a\right)^{2}}{a_{0}}|\nabla u|^{2}\right)^{\frac{1}{2}}+\left(\int_{D} a_{0} \mid \nabla\left(u_{0}-u_{0, h}\right)^{2}\right)^{\frac{1}{2}}\right] \cdot\left(\int_{D} a_{0}|\nabla v|^{2}\right)^{\frac{1}{2}} .
\end{align*}
$$

Thanks to the inequality $(a+b)^{2} \leq 2\left(a^{2}+b^{2}\right), v=u(\cdot, \mathbf{Y}(\omega))-u_{0, h} \in V$ a.s. in $\Omega$ in the last inequality yields

$$
\begin{equation*}
\left(\int_{D} a_{0}\left|\nabla\left(u-u_{0, h}\right)\right|^{2}\right)^{\frac{1}{2}} \leq \sqrt{2}\left[\frac{1}{a_{0, \min }} \int_{D}\left(a-a_{0}\right)^{2}|\nabla u|^{2}+\int_{D} a_{0}\left|\nabla\left(u_{0}-u_{0, h}\right)\right|^{2}\right]^{\frac{1}{2}} . \tag{1.34}
\end{equation*}
$$

The second term of the right-hand side of (1.34) can be bounded in a standard manner as follows. Using the Galerkin orthogonality property

$$
\int_{D} a_{0} \nabla\left(u_{0}-u_{0, h}\right) \cdot \nabla v_{h}=0 \quad \forall v_{h} \in V_{h},
$$

we easily get

$$
\int_{D} a_{0}\left|\nabla\left(u_{0}-u_{0, h}\right)\right|^{2} \leq a_{0, \max }\left\|\nabla\left(u_{0}-\mathscr{I}_{h} u_{0}\right)\right\|_{L^{2}(D)}^{2} .
$$

Since $u_{0} \in H^{2}(D)$ by assumption, thanks to the interpolation result (1.25) we get

$$
\begin{equation*}
\int_{D} a_{0}\left|\nabla\left(u_{0}-u_{0, h}\right)\right|^{2} \leq a_{0, \max } C^{2} h^{2}\left|u_{0}\right|_{H^{2}(D)}^{2} \tag{1.35}
\end{equation*}
$$

Therefore, using this last relation and the lower bound for $a_{0}$ in (1.34) yields a.s. in $\Omega$

$$
\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2} \leq 2\left[\frac{a_{0, \text { max }}}{a_{0, \text { min }}} C^{2} h^{2}\left|u_{0}\right|_{H^{2}(D)}^{2}+\frac{1}{a_{0, \text { min }}^{2}} \int_{D}\left(a-a_{0}\right)^{2}|\nabla u|^{2}\right]
$$

Then, we take the expected value on both sides of the last inequality to get

$$
\begin{equation*}
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right] \leq 2\left[\frac{a_{0, \max }}{a_{0, \text { min }}} C^{2} h^{2}\left|u_{0}\right|_{H^{2}(D)}^{2}+\frac{1}{a_{0, \text { min }}^{2}} \mathbb{E}\left[\int_{D}\left(a-a_{0}\right)^{2}|\nabla u|^{2}\right]\right] \tag{1.36}
\end{equation*}
$$

To complete the proof, we finally bound the expected value that appears on the right-hand side of (1.36). First, using the relation $\left(\sum_{j=1}^{L} x_{j}\right)^{2} \leq L \sum_{j=1}^{L} x_{j}^{2}$, we easily get

$$
\mathbb{E}\left[\int_{D}\left(a-a_{0}\right)^{2}|\nabla u|^{2}\right] \leq L \frac{\varepsilon^{2} C_{P}^{2}}{a_{m i n}^{2}}\|f\|_{L^{2}(D)}^{2} \sum_{j=1}^{L}\left\|a_{j}^{2}\right\|_{L^{\infty}(D)}
$$

which proves (1.31). For (1.32), we use the additional assumption and the relation $\sum_{i} a_{i} b_{i} \leq$ $\left(\sum_{i} a_{i}^{2}\right)^{\frac{1}{2}}\left(\sum_{i} b_{i}^{2}\right)^{\frac{1}{2}}$ to obtain

$$
\left(a-a_{0}\right)^{2}=\varepsilon^{2}\left(\sum_{j=1}^{L} a_{j} j^{\alpha} j^{-\alpha} Y_{j}\right)^{2} \leq \varepsilon^{2}\left(\sum_{j=1}^{L} a_{j}^{2} j^{2 \alpha}\right)\left(\sum_{j=1}^{L} Y_{j}^{2} j^{-2 \alpha}\right) \leq M_{\alpha} \varepsilon^{2} \sum_{j=1}^{L} Y_{j}^{2} j^{-2 \alpha}
$$

Therefore, thanks to (1.18) and the fact that $\mathbb{E}\left[Y_{j}^{2}\right]=1$, we obtain

$$
\mathbb{E}\left[\int_{D}\left(a-a_{0}\right)^{2}|\nabla u|^{2}\right] \leq M_{\alpha} \frac{\varepsilon^{2} C_{P}^{2}}{a_{\min }^{2}}\|f\|_{L^{2}(D)}^{2} \sum_{j=1}^{L} j^{-2 \alpha} \leq M_{\alpha} \frac{\varepsilon^{2} C_{P}^{2}}{a_{\min }^{2}}\|f\|_{L^{2}(D)}^{2} \sum_{j=1}^{\infty} j^{-2 \alpha}
$$

Since $\alpha>\frac{1}{2}$, the series $\sum_{j=1}^{\infty} j^{-2 \alpha}$ converges which concludes the proof.

## Mean of the error estimate

We are now interested in the error on the law of $u$. We restrict ourselves, in particular, to the $H_{0}^{1}(D)$ norm of the expected value of $u-u_{0, h}$. In this case, the statistical error is of order 2 , to be compared to the order 1 of the strong error. Under the same regularity condition on $u_{0}$, we can show the following a priori error estimate.

Proposition 1.3.3. Let $u$ and $u_{0}$ be the (weak) solutions of problems (1.11) and (1.20), respectively, and let $u_{0, h}$ be the solution of problem (1.30). If $u_{0} \in H^{2}(D)$, then we have the a priori error estimate

$$
\begin{equation*}
\left\|\mathbb{E}\left[u-u_{0, h}\right]\right\|_{H_{0}^{1}(D)} \leq \sqrt{\frac{a_{0, \max }}{a_{0, \min }}} C_{1} h\left|u_{0}\right|_{H^{2}(D)}+\frac{\varepsilon^{2} C_{P}}{a_{0, \text { min }}^{3}}\|f\|_{L^{2}(D)} \sum_{j=1}^{L}\left\|a_{j}\right\|_{L^{\infty}(D)}^{2}+C_{2} \varepsilon^{3} \tag{1.37}
\end{equation*}
$$

where $C_{1}>0$ is the constant in (1.25) and $C_{2}$ is a constant independent of $u, h$ and $\varepsilon$. Therefore,
there exists a constant $\tilde{C}>0$ independent of $h$ and $\varepsilon$ such that

$$
\left\|\mathbb{E}\left[u-u_{0, h}\right]\right\|_{H_{0}^{1}(D)} \leq \tilde{C}\left(h+\varepsilon^{2}\right) .
$$

Proof. Let us define $u_{1}=\sum_{j=1}^{L} U_{j} Y_{j}$, where $U_{j}$ is the solution of problem (1.21) for $j=1, \ldots, L$. First, the expected value of the error $u(\cdot, \mathbf{Y})-u_{0, h}$ naturally splits into two parts

$$
\mathbb{E}\left[u-u_{0, h}\right]=\mathbb{E}\left[u-u_{0}\right]+\left(u_{0}-u_{0, h}\right)
$$

and thus, thanks to the triangle inequality, we get

$$
\left\|\mathbb{E}\left[u-u_{0, h}\right]\right\|_{H_{0}^{1}(D)} \leq\left\|\mathbb{E}\left[u-u_{0}\right]\right\|_{H_{0}^{1}(D)}+\left\|u_{0}-u_{0, h}\right\|_{H_{0}^{1}(D)}
$$

From (1.35), we deduce a bound for the second term given by

$$
\left\|u_{0}-u_{0, h}\right\|_{H_{0}^{1}(D)} \leq \sqrt{\frac{a_{0, \max }}{a_{0, \min }}} C_{1} h\left|u_{0}\right|_{H^{2}(D)}
$$

where $C_{1}$ is the constant that appears in (1.25). Let us bound the term $\left\|\mathbb{E}\left[u-u_{0}\right]\right\|_{H_{0}^{1}(D)}$, which is due to the uncertainty in the diffusion coefficient. Proceeding as in (1.29) and using the fact that $\int_{D}\left(a_{j} \nabla u_{0}+a_{0} \nabla U_{j}\right) \cdot \nabla v=0$ for all $v \in V$, the following equalities hold for any $v \in V$ and a.s. in $\Omega$

$$
\begin{align*}
\int_{D} a \nabla\left(u-u^{1}\right) \cdot \nabla v & =-\varepsilon \int_{D} a_{0} \nabla u_{1} \cdot \nabla v-\int_{D}\left(a-a_{0}\right) \nabla u^{1} \cdot \nabla v \\
& =-\varepsilon \sum_{j=1}^{L} Y_{j} \int_{D}\left(a_{0} \nabla U_{j}+a_{j} \nabla u_{0}\right) \cdot \nabla v-\varepsilon^{2} \int_{D} \sum_{i, j=1}^{L} Y_{i} Y_{j} a_{j} \nabla U_{i} \cdot \nabla v \\
& =-\varepsilon^{2} \int_{D} \sum_{i, j=1}^{L} Y_{i} Y_{j} a_{j} \nabla U_{i} \cdot \nabla v \tag{1.38}
\end{align*}
$$

Therefore, we have

$$
\int_{D} a_{0} \nabla\left(u-\left(u_{0}+\varepsilon u_{1}\right)\right) \cdot \nabla v=-\int_{D}\left(a-a_{0}\right) \nabla\left(u-\left(u_{0}+\varepsilon u_{1}\right)\right) \cdot \nabla v-\varepsilon^{2} \sum_{i, j=1}^{L} Y_{i} Y_{j} \int_{D} a_{i} \nabla U_{j} \cdot \nabla v
$$

Since $\mathbb{E}\left[u_{1}\right]=0$ and $\mathbb{E}\left[Y_{i} Y_{j}\right]=\delta_{i j}$, where $\delta_{i j}$ denotes the Kronecker delta, taking the expected value on both sides of last equality yields

$$
\int_{D} a_{0} \nabla \mathbb{E}\left[u-u_{0}\right] \cdot \nabla v=\mathbb{E}\left[-\int_{D}\left(a-a_{0}\right) \nabla\left(u-\left(u_{0}+\varepsilon u_{1}\right)\right) \cdot \nabla v\right]-\varepsilon^{2} \sum_{j=1}^{L} \int_{D} a_{j} \nabla U_{j} \cdot \nabla v
$$

### 1.3. Error analysis for the first order approximation

Thanks to Jensen's inequality (see e.g. [89]), we obtain

$$
\begin{aligned}
\int_{D} a_{0} \nabla \mathbb{E}\left[u-u_{0}\right] \cdot \nabla v \leq & \mathbb{E}\left[\left\|a-a_{0}\right\|_{L^{\infty}(D)}\left\|\nabla\left(u-\left(u_{0}+\varepsilon u_{1}\right)\right)\right\|_{L^{2}(D)}\right]\|\nabla v\|_{L^{2}(D)} \\
& +\varepsilon^{2}\|\nabla v\|_{L^{2}(D)} \sum_{j=1}^{L}\left\|a_{j}\right\|_{L^{\infty}(D)}\left\|\nabla U_{j}\right\|_{L^{2}(D)}
\end{aligned}
$$

If we take $v=\mathbb{E}\left[u-u_{0}\right]$ in the last inequality, we get
$\left\|\mathbb{E}\left[u-u_{0}\right]\right\|_{H_{0}^{1}(D)} \leq \frac{1}{a_{0, \text { min }}}\left\{\mathbb{E}\left[\left\|a-a_{0}\right\|_{L^{\infty}(D)}\left\|\nabla\left(u-\left(u_{0}+\varepsilon u_{1}\right)\right)\right\|_{L^{2}(D)}\right]+\varepsilon^{2} \sum_{j=1}^{L}\left\|a_{j}\right\|_{L^{\infty}(D)}\left\|\nabla U_{j}\right\|_{L^{2}(D)}\right\}$.
We now give a bound on $\left\|\nabla U_{j}\right\|_{L^{2}(D)}, j=1, \ldots, L$. First, using standard techniques (CauchySchwarz, Poincaré inequalities, lower bound for $a_{0}$ ), we get the following bound on the solution of problem (1.20)

$$
\left\|\nabla u_{0}\right\|_{L^{2}(D)} \leq \frac{C_{P}}{a_{0, \text { min }}}\|f\|_{L^{2}(D)}
$$

Then, taking $v=U_{j}$ as test function in the weak formulation of problem (1.21) yields

$$
a_{0, \min }\left\|\nabla U_{j}\right\|_{L^{2}(D)}^{2} \leq \int_{D} a_{0}\left|\nabla U_{j}\right|^{2}=-\int_{D} a_{j} \nabla u_{0} \cdot \nabla U_{j} \leq\left\|a_{j}\right\|_{L^{\infty}(D)}\left\|\nabla u_{0}\right\|_{L^{2}(D)}\left\|\nabla U_{j}\right\|_{L^{2}(D)}
$$

and thus

$$
\left\|\nabla U_{j}\right\|_{L^{2}(D)} \leq \frac{C_{P}}{a_{0, \text { min }}^{2}}\|f\|_{L^{2}(D)}\left\|a_{j}\right\|_{L^{\infty}(D)}
$$

Inserting this result in (1.39), we get
$\left\|\mathbb{E}\left[u-u_{0}\right]\right\|_{H_{0}^{1}(D)} \leq \frac{1}{a_{0, \text { min }}}\left\{\mathbb{E}\left[\left\|a-a_{0}\right\|_{L^{\infty}(D)}\left\|\nabla\left(u-\left(u_{0}+\varepsilon u_{1}\right)\right)\right\|_{L^{2}(D)}\right]+\frac{\varepsilon^{2} C_{P}}{a_{0, \text { min }}^{2}}\|f\|_{L^{2}(D)} \sum_{j=1}^{L}\left\|a_{j}\right\|_{L^{\infty}(D)}^{2}\right\}$.

To conclude the proof, we show that the first term of the right-hand side of the last inequality is of higher order in $\varepsilon$, namely of order $\varepsilon^{3}$. Indeed, we have

$$
\left\|a-a_{0}\right\|_{L^{\infty}(D)}=\varepsilon \sum_{j=1}^{L}\left|Y_{j}\right|\left\|a_{j}\right\|_{L^{\infty}(D)} \leq c_{1} \varepsilon
$$

and, taking $v=u-\left(u_{0}+\varepsilon u_{1}\right)$ in (1.38),

$$
\begin{equation*}
\left\|\nabla\left(u-\left(u_{0}+\varepsilon u_{1}\right)\right)\right\|_{L^{2}(D)} \leq \frac{1}{a_{m i n}} \varepsilon^{2} \sum_{i, j=1}^{L}\left|Y_{i} Y_{j}\right|\left\|a_{i}\right\|_{L^{\infty}(D)}\left\|\nabla U_{j}\right\|_{L^{2}(D)} \leq c_{2} \varepsilon^{2} \tag{1.41}
\end{equation*}
$$

with $c_{1}, c_{2}$ two (deterministic) constants independent of $u, h$ and $\varepsilon$. Therefore, we have

$$
\mathbb{E}\left[\left\|a-a_{0}\right\|_{L^{\infty}(D)}\left\|\nabla\left(u-\left(u_{0}+\varepsilon u_{1}\right)\right)\right\|_{L^{2}(D)}\right] \leq C_{2} \varepsilon^{3}
$$

with $C_{2}=c_{1} c_{2}$.

Remark 1.3.4. A bound for $\left\|\mathbb{E}\left[u-u_{0}\right]\right\|_{H_{0}^{1}(D)}$ can also be obtained using Jensen's inequality, the fact that the term $u_{1}$ is mean-free and (1.41) as follows

$$
\begin{aligned}
\left\|\mathbb{E}\left[u-u_{0}\right]\right\|_{H_{0}^{1}(D)} & =\left\|\mathbb{E}\left[u-u_{0}-\varepsilon u_{1}\right]\right\|_{H_{0}^{1}(D)} \\
& \leq \mathbb{E}\left[\left\|\nabla\left(u-\left(u_{0}+\varepsilon u_{1}\right)\right)\right\|_{L^{2}(D)}\right] \\
& \leq \frac{\varepsilon^{2} C_{P}}{a_{\text {min }} a_{0, \text { min }}^{2}}\|f\|_{L^{2}(D)}\left(\sum_{j=1}^{L}\left\|a_{j}\right\|_{L^{\infty}(D)}\right)^{2} .
\end{aligned}
$$

Compared to (1.37), there is no additional higher order term here but the constant for the term of order $\varepsilon^{2}$ is larger since the cross terms do not vanish and $a_{0, \text { min }}^{-1}$ is replaced by $a_{\min }^{-1}$.

### 1.3.2 A posteriori error analysis

A posteriori error estimate in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm
The goal is now to obtain an estimate of the error between $u$ and $u_{0, h}$ which does not depend on the exact (unknown) solution. Let us define the jump of a function $\varphi$ across an edge $e \in \mathscr{T}_{h}$ in the direction of $\mathbf{n}_{e}$ by

$$
[\varphi]_{\mathbf{n}_{e}}(\mathbf{x}):=\left\{\begin{aligned}
\lim _{t \rightarrow 0^{+}}\left(\varphi\left(\mathbf{x}+t \mathbf{n}_{e}\right)-\varphi\left(\mathbf{x}-t \mathbf{n}_{e}\right)\right) & \text { if } e \not \subset \partial D \\
0 & \text { if } e \subset \partial D
\end{aligned}\right.
$$

where $\mathbf{n}_{e}$ denotes a normal vector to $e$ of arbitrary (but fixed) direction for internal edges and the outwards normal to $\partial D$ if $e \in \partial D$. Notice that the quantity $\left[\nabla \varphi \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}$ is independent of the choice of the direction of the normal vector $\mathbf{n}_{e}$. We obtain the following residual type error upper bound, proceeding as in [118], which is based on the relation

$$
\mathscr{A}\left(u-u_{0, h}, v ; \mathbf{y}\right)=\mathscr{R}\left(v ; \mathbf{y}_{0}\right)+\left[\mathscr{R}(v ; \mathbf{y})-\mathscr{R}\left(v ; \mathbf{y}_{0}\right)\right] \quad \forall v \in H_{0}^{1}(D), \rho \text {-a.e. in } \Gamma
$$

with

$$
\mathscr{R}(\nu ; \mathbf{y}):=F(\nu)-\mathscr{A}\left(u_{0, h}, \nu ; \mathbf{y}\right),
$$

where $\mathscr{A}$ and $F$ are defined in (1.15) and (1.16), respectively, and $\mathbf{y}_{0}=\mathbb{E}[\mathbf{Y}]=\mathbf{0}$.

Proposition 1.3.5. Let $u$ be the weak solution of problem (1.11) and let $u_{0, h}$ be the solution of problem (1.30), respectively. There exists a constant $C>0$ depending only on the constants in (1.26) and (1.28) such that

$$
\begin{equation*}
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{\sqrt{2}}{a_{\min }}\left[C \eta_{1}^{2}+\eta_{2}^{2}\right]^{\frac{1}{2}}, \tag{1.42}
\end{equation*}
$$

with

$$
\begin{align*}
& \eta_{1}^{2}:=\sum_{K \in \mathscr{T}_{h}} h_{K}^{2} \int_{K}\left(f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{\mathbf{e}}\right]_{\mathbf{n}_{\mathbf{e}}}^{2}  \tag{1.43}\\
& \eta_{2}^{2}:=\varepsilon^{2} \int_{D} \sum_{j=1}^{L} a_{j}^{2}\left|\nabla u_{0, h}\right|^{2} . \tag{1.44}
\end{align*}
$$

Remark 1.3.6. We mention that the analysis is similar to the one given below if we consider the error in the energy norm $\left\|a_{0}^{1 / 2} \nabla\left(u-u_{0, h}\right)\right\|_{L_{p}^{2}\left(\Omega ; L^{2}(D)\right)}$ instead of $\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L_{p}^{2}\left(\Omega ; L^{2}(D)\right)}$. The former should be preferred if the deterministic part $a_{0}$ of the diffusion coefficient a varies widely over $D$.

Proof. In the sequel, $C$ will denote a constant whose value might change from one line to another. Let $v$ be any function in $H_{0}^{1}(D)$. We have a.s. in $\Omega$

$$
\begin{align*}
\int_{D} a \nabla\left(u-u_{0, h}\right) \cdot \nabla v & =\int_{D} a \nabla u \cdot \nabla v-\int_{D} a \nabla u_{0, h} \cdot \nabla v \\
& =\underbrace{\int_{D}\left(f v-a_{0} \nabla u_{0, h} \cdot \nabla v\right)}_{=: A_{1}}+\underbrace{\int_{D}\left(a_{0}-a\right) \nabla u_{0, h} \cdot \nabla v}_{=: A_{2}}, \tag{1.45}
\end{align*}
$$

where $A_{1}$ and $A_{2}$ correspond respectively to the error due to the finite element approximation of $u_{0}$, solution to problem (1.20), and the error due to the truncation in the expansion (1.19) of $u$. We bound now each term separately, starting with $A_{2}$. Using the expansion of $a$ given by (1.12), we have

$$
\begin{equation*}
A_{2} \leq\left(\int_{D}\left(a-a_{0}\right)^{2}\left|\nabla u_{0, h}\right|^{2}\right)^{\frac{1}{2}}\left(\int_{D}|\nabla \nu|^{2}\right)^{\frac{1}{2}}=\varepsilon\left(\int_{D}\left(\sum_{j=1}^{L} a_{j} Y_{j}\right)^{2}\left|\nabla u_{0, h}\right|^{2}\right)^{\frac{1}{2}}\|\nabla v\|_{L^{2}(D)} \tag{1.46}
\end{equation*}
$$

For the first term $A_{1}$, we use the relation $\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v_{h}=\int_{D} f v_{h}$ for all $v_{h} \in V_{h}$ with $v_{h}$ the Clément interpolant of $v$ together with interpolation results (1.26) to get

$$
\begin{align*}
A_{1} \leq & \sum_{K \in \mathscr{T}_{h}}\left(\left.\int_{K}\left|f+\nabla \cdot\left(\left.a_{0} \nabla u_{0, h}\right|^{2}\right)^{\frac{1}{2}} C h_{K}\right| \nu\right|_{H^{1}(N(K))}\right. \\
& +\sum_{e \in \mathscr{T}_{h}}\left(\int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{\mathrm{e}}\right]_{\mathbf{n}_{\mathrm{e}}}^{2}\right)^{\frac{1}{2}} C h_{e}^{\frac{1}{2}}|\nu|_{H^{1}\left(N\left(K_{e}\right)\right)} \\
\leq & \left.\sqrt{2} C\left[\sum_{K \in \mathscr{T}_{h}} h_{K}^{2} \int_{K}\left|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right|^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{\mathbf{e}}\right]_{\mathbf{n}_{\mathbf{e}}}^{2}\right]\right]^{\frac{1}{2}}\|\nabla v\|_{L^{2}(D)} . \tag{1.47}
\end{align*}
$$

We have used the fact that

$$
\sum_{K \in \mathscr{T}_{h}}\|\nabla v\|_{L^{2}(N(K))}^{2} \leq C_{0}\|\nabla v\|_{L^{2}(D)}^{2} \quad \text { and } \quad \sum_{e \in \mathscr{T}_{h}}\|\nabla v\|_{L^{2}\left(N\left(K_{e}\right)\right)}^{2} \leq C_{0}\|\nabla v\|_{L^{2}(D)}^{2}
$$

where $C_{0}$ depends on the maximum number of neighbours of each element in $\mathscr{T}_{h}$, which in turn depends on the constant in (1.23). Since $a_{\min }$ is a lower bound for $a$, we deduce from (1.45) with $v=u(\cdot, \mathbf{Y}(\omega))-u_{0, h} \in H_{0}^{1}(D)$ that a.s. in $\Omega$ we have

$$
\int_{D}\left|\nabla\left(u-u_{0, h}\right)\right|^{2} \leq \frac{1}{a_{\min }}\left[A_{1}+A_{2}\right] .
$$

Combining this last inequality with the bounds for $A_{1}$ and $A_{2}$ given by (1.47) and (1.46) respectively, we obtain a.s. in $\Omega$

$$
\begin{align*}
\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)} \leq & \frac{1}{a_{\min }}\left\{\sqrt { 2 } C \left[\sum_{K \in \mathscr{T}_{h}} h_{K}^{2} \int_{K}\left(f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)^{2}\right.\right. \\
& \left.\left.+\sum_{e \in \mathscr{T}_{h}} h_{e} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{\mathbf{e}}\right]_{\mathbf{n}_{\mathbf{e}}}^{2}\right]^{\frac{1}{2}}+\varepsilon\left(\int_{D}\left(\sum_{j=1}^{L} a_{j} Y_{j}\right)^{2}\left|\nabla u_{0, h}\right|^{2}\right)^{\frac{1}{2}}\right\} \tag{1.48}
\end{align*}
$$

and thus, taking the square of this last equation and using again $(a+b)^{2} \leq 2\left(a^{2}+b^{2}\right)$ yields

$$
\begin{aligned}
\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2} \leq & \frac{2}{a_{\min }^{2}}\left\{2 C ^ { 2 } \left(\sum_{K \in \mathscr{T}_{h}} h_{K}^{2} \int_{K}\left|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right|^{2}\right.\right. \\
& \left.\left.+\sum_{e \in \mathscr{T}_{h}} h_{e} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{\mathbf{e}}\right]_{\mathbf{n}_{\mathbf{e}}}^{2}\right)+\varepsilon^{2} \int_{D}\left(\sum_{j=1}^{L} a_{j} Y_{j}\right)^{2}\left|\nabla u_{0, h}\right|^{2}\right\} .
\end{aligned}
$$

The a posteriori error estimate (1.42) is obtained taking the square root of the expected value on both sides of the last inequality and exploiting the independence of the random variables, namely that $\mathbb{E}\left[Y_{i} Y_{j}\right]=\delta_{i j}$ for $i, j=1, \ldots, L$.

Remark 1.3.7. In the one-dimensional case, we can take $v_{h}=r_{h} v$ the Lagrange interpolant of $v$ and the sum over the edges (the discrete nodes here) vanishes. Indeed, any function and its Lagrange interpolant coincide at each node $x_{i}, i=0, \ldots, N_{h}$, of the considered discretization, or more precisely $v\left(x_{i}\right)-r_{h} v\left(x_{i}\right)=0$ for all $i=0, \ldots, N_{h}$. Since (1.24) holds for e.g. $C=2$, we can show that we have the following a posteriori error estimate

$$
\begin{equation*}
\mathbb{E}\left[\left\|u^{\prime}-u_{0, h}^{\prime}\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{\sqrt{2}}{a_{\min }}\left(4 \sum_{i=0}^{N_{h}-1} h_{i}^{2} \int_{x_{i}}^{x_{i+1}}\left(f+\left(a_{0} u_{0, h}^{\prime}\right)^{\prime}\right)^{2}+\varepsilon^{2} \int_{D} \sum_{j=1}^{L} a_{j}^{2}\left(u_{0, h}^{\prime}\right)^{2}\right)^{\frac{1}{2}} \tag{1.49}
\end{equation*}
$$

where $u^{\prime}$ denotes the spatial derivative $\frac{\partial u(x, \omega)}{\partial x}$.
Remark 1.3.8. The computable quantity $\eta=\left(\eta_{1}^{2}+\eta_{2}^{2}\right)^{\frac{1}{2}}$ can be used as an a posteriori error
estimator, which is reliable thanks to (1.42). It can be used to determine a mesh yielding comparable accuracy in $h$ and $\varepsilon$, i.e. for balancing the error due to physical space discretization and the error due to the uncertainty. The spatial error estimator $\eta_{1}$ is efficient in the sense that it provides (up to a multiplicative constant depending only on $a_{\text {max }}$ and the regularity of the mesh) a lower bound for the error plus the other contribution $\eta_{2}$ and oscillation terms, the proof being similar to the one given in Appendix 1.B. Even though we have not been able to prove that $\eta_{2}$ in (1.44) also provides a similar lower bound, the estimator $\eta$ appears to be efficient for all the numerical experiments we have considered.

We give below an a posteriori error estimator for the error $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ for which both upper and lower bounds can be shown. The spatial error estimator is the same, namely $\eta_{1}$ given in (1.43), while the stochastic error estimator is obtained by computing (approximately) the dual norm of the residual $r(\nu ; \mathbf{y}):=\mathscr{R}(\nu ; \mathbf{y})-\mathscr{R}\left(\nu ; \mathbf{y}_{0}\right)$. Here, we only give the statement of the error estimator and we refer to Appendix 1.B for more details including the proof of the bounds. Let $W_{j, h} \in V_{h}$ be the solution of the problem

$$
\int_{D} \nabla W_{j, h} \cdot \nabla v_{h}=-\int_{D} a_{j} \nabla u_{0, h} \cdot \nabla v_{h} \quad \forall v_{h} \in V_{h}
$$

The error estimator can then be defined as

$$
\begin{equation*}
\hat{\eta}^{2}=\left(\eta_{1}^{2}+\hat{\eta}_{2}^{2}\right)^{\frac{1}{2}} \quad \text { with } \quad \hat{\eta}_{2}^{2}:=\varepsilon^{2} \sum_{j=1}^{L}\left\|\nabla W_{j, h}\right\|_{L^{2}(D)}^{2} \tag{1.50}
\end{equation*}
$$

Notice that the computation of $\hat{\eta}$ in (1.50) requires the solution of $L$ additional Poisson problems compared to the error estimator $\eta$ based on (1.42), and a strategy to reduce the computational cost could be to introduce auxiliary local problems defined on an element or a small subdomain, see e.g. $[15,107]$ and references therein. We mention that the extra computational effort to get $\hat{\eta}_{2}$ instead of $\eta_{2}$ is apparently not worth to pay in the present case, since the $a$ posteriori error estimator based on Proposition 1.3.5 is efficient, at least for all the numerical experiments we have performed.

## A posteriori error estimate in the $L_{P}^{2}\left(\Omega ; L^{2}(D)\right)$ norm

We now give an a posteriori error estimate of the error between $u$ and $u_{0, h}$ in the $L^{2}$ norm in space, which leads to a gain of one order in $h$. To do so, we use a duality argument (often called the Aubin-Nitsche trick). We thus consider the dual problem of problem (1.11) given by:
find $\phi: D \times \Omega \rightarrow \mathbb{R}$ such that $P$-almost everywhere:

$$
\left\{\begin{align*}
-\operatorname{div}(a(\mathbf{x}, \omega) \nabla \phi(\mathbf{x}, \omega)) & =u(\mathbf{x}, \omega)-u_{0, h}(\mathbf{x}) & & \mathbf{x} \in D  \tag{1.51}\\
\phi(\mathbf{x}, \omega) & =0 & & \mathbf{x} \in \partial D
\end{align*}\right.
$$

whose pointwise in $\mathbf{y} \in \Gamma$ weak form reads:
find $\phi(\cdot, \mathbf{y}) \in H_{0}^{1}(D)$ such that

$$
\begin{equation*}
\int_{D} a(\mathbf{x}, \mathbf{y}) \nabla \phi(\mathbf{x}, \mathbf{y}) \cdot \nabla v(\mathbf{x}) d \mathbf{x}=\int_{D}\left(u(\mathbf{x}, \mathbf{y})-u_{0, h}(\mathbf{x})\right) v(\mathbf{x}) d \mathbf{x} \quad \forall v \in H_{0}^{1}(D), \rho \text {-a.e. in } \Gamma . \tag{1.52}
\end{equation*}
$$

Under regularity conditions on $D$, we have the following a posteriori error upper bound, which implies that the convergence rate of the error is $\mathscr{O}\left(h^{2}+\varepsilon\right)$ in that case. That is that we gain one order in $h$ compared to the error in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm. However, the order of the statistical error is not improved.

Proposition 1.3.9. Let $u$ and $u_{0}$ be the (weak) solutions of problems (1.11) and (1.20), respectively, and let $u_{0, h}$ be the solution of problem (1.30). If $\phi(\cdot, \mathbf{Y}(\omega)) \in H^{2}(D)$ and $\|\phi\|_{H^{2}(D)} \leq$ $C\left\|u-u_{0, h}\right\|_{L^{2}(D)}$ a.s. in $\Omega$, then there exist constants $C_{1}, C_{2}>0$ independent of $u, h$ and $\varepsilon$ such that

$$
\begin{equation*}
\mathbb{E}\left[\left\|u-u_{0, h}\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \sqrt{2}\left[C_{1} \eta_{1}^{2}+C_{2} \eta_{2}^{2}\right]^{\frac{1}{2}} \tag{1.53}
\end{equation*}
$$

with

$$
\begin{align*}
\eta_{1}^{2} & :=\sum_{K \in \mathscr{T}_{h}} h_{K}^{4} \int_{K}\left(f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e}^{3} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}^{2}  \tag{1.54}\\
\eta_{2}^{2} & :=\varepsilon^{2} \int_{D} \sum_{j=1}^{L} a_{j}^{2}\left|\nabla u_{0, h}\right|^{2} . \tag{1.55}
\end{align*}
$$

Remark 1.3.10. Since we assumed $a_{j} \in W^{1, \infty}(D), j=0, \ldots, L$, the assumptions of Proposition 1.3.9 on the regularity of the dual solution $\phi$ are satisfied if, for instance, $D$ is a convex polygon (see [84]). The constant $C$ in $\|\phi\|_{H^{2}(D)} \leq C\left\|u-u_{0, h}\right\|_{L^{2}(D)}$ may depend on the uniform bounds of $Y_{j}, a_{j}$ and $\nabla a_{j}$ and on $\varepsilon_{\max }$ but is independent of $\varepsilon$.

Proof. First note that if we take $v=u(\cdot, \mathbf{y})-u_{0, h}, \rho$-a.e. in $\Gamma$, in (1.52), we directly get the $L^{2}$ norm in space of the error at the right-hand side. We thus only need to estimate the lefthand side by a quantity which does not depend on the exact solutions $u=u(\mathbf{x}, \mathbf{Y}(\omega))$ and $\phi=\phi(\mathbf{x}, \mathbf{Y}(\omega))$ of respectively the primal and dual problems. In what follows, all equations hold a.s. in $\Omega$ without specifically mentioning it. Since

$$
\int_{D} a \nabla\left(u-u_{0, h}\right) \cdot \nabla v_{h}+\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla v_{h}=0 \quad \forall v_{h} \in V_{h}
$$

we have for any $v_{h} \in V_{h}$

$$
\begin{align*}
\left\|u-u_{0, h}\right\|_{L^{2}(D)}^{2} & =\int_{D} a \nabla\left(u-u_{0, h}\right) \cdot \nabla \phi \\
& =\int_{D} a \nabla\left(u-u_{0, h}\right) \cdot \nabla\left(\phi-v_{h}\right)-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla v_{h} \\
& =\underbrace{\int_{D} f\left(\phi-v_{h}\right)-\int_{D} a_{0} \nabla u_{0, h} \nabla\left(\phi-v_{h}\right)}_{=: A_{1}} \underbrace{-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla \phi}_{=: A_{2}} . \tag{1.56}
\end{align*}
$$

We now treat each term separately. For the first one, we follow the usual procedure. For any $v_{h} \in V_{h}$, we have

$$
\begin{aligned}
A_{1} & =\sum_{K \in \mathscr{T}_{h}} \int_{K} f\left(\phi-v_{h}\right)-\sum_{K \in \mathscr{\mathscr { T }}_{h}} \int_{K} a_{0} \nabla\left(\phi-v_{h}\right) \nabla u_{0, h} \\
& \leq \sum_{K \in \mathscr{T}_{h}} \| f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\left\|_{L^{2}(K)}\right\| \phi-v_{h}\left\|_{L^{2}(K)}+\sum_{e \in \mathscr{F}_{h}}\right\|\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\left\|_{L^{2}(e)}\right\| \phi-v_{h} \|_{L^{2}(e)} .\right.
\end{aligned}
$$

If we take $v_{h}=r_{h} \phi$, the Lagrange interpolant of $\phi$, thanks to the interpolation error estimate (1.25), the trace inequality and the standard elliptic regularity result $\|\phi\|_{H^{2}(D)} \leq C\left\|u-u_{0, h}\right\|_{L^{2}(D)}$ (see [31,49] for instance), we obtain

$$
\begin{align*}
A_{1} & \leq C_{1}\left[\left(\sum_{K \in \mathscr{F}_{h}} h_{K}^{4} \int_{K}\left(f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)^{2}\right)^{\frac{1}{2}}+\left(\sum_{e \in \mathscr{T}_{h}} h_{e}^{3} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}^{2}\right)^{\frac{1}{2}}\right]|\phi|_{H^{2}(D)} \\
& \leq \sqrt{2} C_{1}\left(\sum_{K \in \mathscr{T}_{h}} h_{K}^{4} \int_{K}\left(f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e}^{3} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}^{2}\right)^{\frac{1}{2}}\left\|u-u_{0, h}\right\|_{L^{2}(D)}, \tag{1.57}
\end{align*}
$$

where $C_{1}$ is a constant whose value might change from one line to another. Consider now the second term $A_{2}$ of (1.56). We have

$$
A_{2}=-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla \phi \leq\left(\int_{D}\left(a-a_{0}\right)^{2}\left|\nabla u_{0, h}\right|^{2}\right)^{\frac{1}{2}}\|\nabla \phi\|_{L^{2}(D)}
$$

and thus, it only remains to obtain an upper bound for $\|\nabla \phi\|_{L^{2}(D)}$. Taking $v=\phi$ in the weak form (1.52) of the dual problem yields

$$
\int_{D} a \nabla \phi \cdot \nabla \phi=\int_{D}\left(u-u_{0, h}\right) \phi \leq\left\|u-u_{0, h}\right\|_{L^{2}(D)}\|\phi\|_{L^{2}(D)} .
$$

Since $a$ is bounded from below by $a_{\text {min }}$, thanks to the Poincaré inequality we get

$$
a_{\min }\|\nabla \phi\|_{L^{2}(D)}^{2} \leq C_{P}\left\|u-u_{0, h}\right\|_{L^{2}(D)}\|\nabla \phi\|_{L^{2}(D)},
$$

and thus

$$
\|\nabla \phi\|_{L^{2}(D)} \leq \frac{C_{P}}{a_{m i n}}\left\|u-u_{0, h}\right\|_{L^{2}(D)} .
$$

Therefore, $A_{2}$ can be bounded by

$$
\begin{equation*}
A_{2} \leq \frac{C_{P}}{a_{m i n}}\left(\int_{D}\left(a-a_{0}\right)^{2}\left|\nabla u_{0, h}\right|^{2}\right)^{\frac{1}{2}}\left\|u-u_{0, h}\right\|_{L^{2}(D)} \tag{1.58}
\end{equation*}
$$

Inserting (1.57) and (1.58) into (1.56) yields

$$
\begin{aligned}
\left\|u-u_{0, h}\right\|_{L^{2}(D)} \leq & \sqrt{2} C_{1}\left(\sum_{K \in \mathscr{T}_{h}} h_{K}^{4} \int_{K}\left(f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e}^{3} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}^{2}\right)^{\frac{1}{2}} \\
& +\frac{C_{P}}{a_{\min }}\left(\int_{D}\left(a-a_{0}\right)^{2}\left|\nabla u_{0, h}\right|^{2}\right)^{\frac{1}{2}}
\end{aligned}
$$

and thus

$$
\begin{align*}
\left\|u-u_{0, h}\right\|_{L^{2}(D)}^{2} \leq & 2\left[2 C_{1}^{2}\left(\sum_{K \in \mathscr{T}_{h}} h_{K}^{4} \int_{K}\left(f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e}^{3} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}^{2}\right)\right. \\
& \left.+\frac{C_{P}^{2}}{a_{\min }^{2}} \int_{D}\left(a-a_{0}\right)^{2}\left|\nabla u_{0, h}\right|^{2}\right] . \tag{1.59}
\end{align*}
$$

Since $\mathbb{E}\left[\left(a-a_{0}\right)^{2}\right]=\varepsilon^{2} \sum_{j=1}^{L} a_{j}^{2}$, the result follows from taking first the expected value and then the square root on both sides of (1.59).

## Goal-oriented error estimate

The a posteriori error estimates obtained so far yield upper bounds on the error in global norms. In the case where we are interested in a particular quantity of interest, e.g. point values or contour integrals, these estimates may not be appropriate. Goal-oriented error estimation has thus been developed (see $[13,22,100]$ and $[4,33,35,92]$ and the references therein for the deterministic and stochastic framework, respectively) to bound a given functional using optimal control techniques (based on a duality-argument). In this section we only sketch the derivation of a goal-oriented error upper bound for the first-order FEM approximation $u_{0, h}$. Assume that we are interested in computing $Q(u)$ with $Q$ a linear functional on $H_{0}^{1}(D)$ representing a quantity of interest which depends on the random vector $\mathbf{Y}$ only through the random solution $u(\cdot, \mathbf{Y})$ itself. We introduce the dual problem:

$$
\begin{equation*}
\text { find } \varphi(\cdot, \mathbf{y}) \in H_{0}^{1}(D) \text { such that } \mathscr{A}(v, \varphi(\cdot, \mathbf{y}) ; \mathbf{y})=Q(\nu), \quad \forall v \in H_{0}^{1}(D), \rho \text {-a.e. in } \Gamma \text {, } \tag{1.60}
\end{equation*}
$$

where $\mathscr{A}$ is defined by (1.15). Let $\mathbf{y}_{0}=\mathbb{E}[\mathbf{Y}]=\mathbf{0}$ denotes the nominal value for $\mathbf{Y}$, for which $a\left(\mathbf{x}, \mathbf{y}_{0}\right)=a_{0}(\mathbf{x})$, and let $\varphi_{0}$ be the deterministic solution of (1.60) with $\mathbf{y}=\mathbf{y}_{0}$ and $\varphi_{0, h}$ its FE approximation. Using the fact that $Q$ does not depend on $\mathbf{Y}$ explicitly, we can easily show that
a.s. in $\Omega$

$$
\begin{aligned}
Q(u(\cdot, \mathbf{Y}(\omega)))-Q\left(u_{0, h}\right)= & \underbrace{\int_{D} f \varphi_{0}-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla \varphi_{0}}_{=: A_{1}} \underbrace{-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla \varphi_{0, h}}_{=: A_{2}} \\
& \underbrace{-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla\left(\varphi_{0}-\varphi_{0, h}\right)}_{=: A_{3}}-\underbrace{-\int_{D}\left(a-a_{0}\right) \nabla\left(u-u_{0, h}\right) \cdot \nabla \varphi_{0, h}}_{=: A_{4}} \\
& \underbrace{-\int_{D}\left(a-a_{0}\right) \nabla\left(u-u_{0, h}\right) \cdot \nabla\left(\varphi_{0}-\varphi_{0, h}\right)} .
\end{aligned}
$$

The first term $A_{1}$, which is deterministic and of order $h^{2}$, can be bounded using standard techniques such as the Dual-weighted residual (DWR) method (see e.g. [13, 22]) or using the parallelogram identity as proposed by Oden and Prudhomme in [100]. In the DWR method, the upper bound depends on the unknown influence function $\varphi_{0}$, either through $\left|\varphi_{0}\right|_{H^{2}(K)}$ or $\left\|\nabla\left(\varphi_{0}-\varphi_{0, h}\right)\right\|_{L^{2}(K)}, K$ being an element of the mesh. In the former case, the $H^{2}$ semi-norm can be estimated by a discrete analogue and in the latter case, the influence function might be replaced by a discrete solution computed on a space richer than $V_{h}$ or by post-processing. All the other terms can be bounded provided we can obtain an upper bound for $\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}$, which is given by (1.48), as well as an upper bound for $\left\|\nabla\left(\varphi_{0}-\varphi_{0, h}\right)\right\|_{L^{2}(D)}$ which can be done as in the previous sections. Moreover, based on the results obtained in the previous sections we have

$$
A_{1}=\mathscr{O}\left(h^{2}\right), A_{2}=\mathscr{O}(\varepsilon), A_{3}=\mathscr{O}(h \varepsilon), A_{4}=\mathscr{O}\left(h \varepsilon+\varepsilon^{2}\right) \text { and } A_{5}=\mathscr{O}\left(h^{2} \varepsilon+\varepsilon^{2} h\right)
$$

We might be interested in estimating the expectation or the variance of $Q(u(\cdot, \mathbf{Y}))-Q\left(u_{0, h}\right)$. In the former case, notice that $\mathbb{E}\left[A_{2}\right]=\mathbb{E}\left[A_{3}\right]=0$ and since $A_{1}$ is a deterministic quantity, we have

$$
\mathbb{E}\left[Q(u)-Q\left(u_{0, h}\right)\right]=A_{1}+\mathbb{E}\left[A_{4}\right]+\mathbb{E}\left[A_{5}\right] .
$$

Moreover, the term $\mathbb{E}\left[A_{5}\right]$ is of higher order than $\mathbb{E}\left[A_{4}\right]$ and can thus be neglected, so that we have $\mathbb{E}\left[Q(u)-Q\left(u_{0, h}\right)\right]=\mathscr{O}\left(h^{2}+h \varepsilon+\varepsilon^{2}\right)$. In the latter case, we have

$$
\mathbb{E}\left[\left|Q(u)-Q\left(u_{0, h}\right)\right|^{2}\right] \leq 5\left(A_{1}^{2}+\mathbb{E}\left[A_{2}^{2}\right]+\mathbb{E}\left[A_{3}^{2}\right]+\mathbb{E}\left[A_{4}^{2}\right]+\mathbb{E}\left[A_{5}^{2}\right]\right)
$$

As before, the term $\mathbb{E}\left[A_{5}^{2}\right]$ can be neglected and we have $\mathbb{E}\left[\left|Q(u)-Q\left(u_{0, h}\right)\right|^{2}\right]^{\frac{1}{2}}=\mathscr{O}\left(h^{2}+\varepsilon+h \varepsilon\right)$. Moreover, if the mesh space $h$ is chosen such that $h^{2} \sim \varepsilon$, then both terms $\mathbb{E}\left[A_{3}^{2}\right]$ and $\mathbb{E}\left[A_{4}^{2}\right]$ can also be omitted in the estimation of the variance and $\mathbb{E}\left[\mid Q(u)-Q\left(u_{0, h}\right)^{2}\right]^{\frac{1}{2}}=\mathscr{O}\left(h^{2}+\varepsilon\right)$.

Finally, we mention that the estimate on the variance of $Q(u)-Q\left(u_{0, h}\right)$ can be used to have a rough estimate on the failure probability $P\left(Q(u)>Q_{c r i t}\right)$ with some critical value $Q_{c r i t}$
sufficiently far from $Q\left(u_{0, h}\right)$. Indeed, using the Bienaymé-Tchebychev inequality we have

$$
P\left(Q(u)>Q_{c r i t}\right) \leq \frac{\mathbb{E}\left[\left(Q(u)-Q\left(u_{0, h}\right)\right)^{2}\right]}{\left(Q\left(u_{0, h}\right)-Q_{c r i t}\right)^{2}} .
$$

### 1.4 Error analysis for higher order approximations

In this section, we generalize the a posteriori error estimate of Proposition 1.3.5 to higher order approximation, that is when more terms in the expansion (1.19) of $u$ are taken into account. We start by giving the result for the second order approximation before generalizing to any order of approximation.

### 1.4.1 Second order approximation

In this section, instead of considering the error between $u$ and $u_{0, h}$, we will give an estimation of the error between $u$ and $u_{h}^{1}$, the FE approximation of $u^{1}:=u_{0}+\varepsilon u_{1}=u_{0}+\varepsilon \sum_{j=1}^{L} U_{j} Y_{j}$, where $U_{j}$ is the solution of problem (1.21). Since the random variables $Y_{j}, j=1, \ldots, L$, are assumed to be bounded, the error due to the stochastic approximation of $u$ is of order $\varepsilon^{2}$ in this case. Indeed, if we do not take the finite element approximation error into account, we have a.s. in $\Omega$ (see (1.38) for details)

$$
\begin{equation*}
\int_{D} a \nabla\left(u-u^{1}\right) \cdot \nabla v=-\varepsilon^{2} \int_{D} \sum_{i, j=1}^{L} Y_{i} Y_{j} a_{j} \nabla U_{i} \cdot \nabla v \tag{1.61}
\end{equation*}
$$

and only the term of order $\varepsilon^{2}$ remains. Let us now take the error due to the approximation of $u^{1}$ by $u_{h}^{1}:=u_{0, h}+\varepsilon u_{1, h}$ into account, where $u_{1, h}=\sum_{j=1}^{L} Y_{j} U_{j, h}$ and, for $j=1, \ldots, L, U_{j, h}$ is the solution of

$$
\begin{equation*}
\int_{D} a_{0} \nabla U_{j, h} \cdot \nabla v_{h}=-\int_{D} a_{j} \nabla u_{0, h} \cdot \nabla v_{h} \quad \forall v_{h} \in V_{h} \tag{1.62}
\end{equation*}
$$

To simplify the notation, we define

$$
w_{j, h}:=a_{0} \nabla U_{j, h}+a_{j} \nabla u_{0, h}
$$

We can show that, if the solution is regular enough in physical space, the convergence of the error is in $\mathscr{O}\left(h+\varepsilon h+\varepsilon^{2}\right)$, i.e., that for a mesh size $h$ of order $\varepsilon^{2}$, the error is divided by 4 when $\varepsilon$ is halved. The following proposition provides an a posteriori error estimate.

Proposition 1.4.1. Let $u$ be the weak solution of problem (1.11) and let $u_{0, h}$ and $U_{j, h}, j=$ $1, \ldots, L$, be the solutions of problems (1.30) and (1.62), respectively. There exist two constants $C_{1}, C_{2}>0$ depending only on the constants in (1.26) and (1.28) such that

$$
\begin{equation*}
\mathbb{E}\left[\left\|\nabla\left(u-u_{h}^{1}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{\sqrt{3}}{a_{\min }}\left[C_{1} \eta_{1}^{2}+C_{2} \eta_{2}^{2}+\eta_{3}^{2}\right]^{\frac{1}{2}}, \tag{1.63}
\end{equation*}
$$

with

$$
\begin{align*}
& \eta_{1}^{2}=\sum_{K} h_{K}^{2}\left\|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e} h_{e}\left\|\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{\left.L^{2}(e)\right)^{\prime}}^{2},  \tag{1.64}\\
& \eta_{2}^{2}=\varepsilon^{2}\left(\sum_{K} h_{K}^{2} \int_{K} \sum_{j=1}^{L}\left(\nabla \cdot w_{j, h}\right)^{2}+\sum_{e} h_{e} \int_{e} \sum_{j=1}^{L}\left[w_{j, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}^{2}\right),  \tag{1.65}\\
& \eta_{3}^{2}=\varepsilon^{4}\left(\int_{D} \sum_{i=1}^{L} a_{i}^{2}\left|\nabla U_{i, h}\right|^{2} \mathbb{E}\left[Y_{i}^{4}\right]+\int_{D_{D}} \sum_{\substack{i, j=1 \\
i \neq j}}^{L}\left[a_{i}^{2}\left|\nabla U_{j, h}\right|^{2}+2 a_{i} a_{j} \nabla U_{i, h} \cdot \nabla U_{j, h}\right]\right) . \tag{1.66}
\end{align*}
$$

From (1.63), we see that the error splits into three parts, namely the error due to the FE approximation of $u_{0}$, the FE approximation of the $U_{j}, j=1, \ldots, L$ and the truncation in the expansion of $u$ with respect to $\varepsilon$.

Proof. For any $v \in H_{0}^{1}(D)$ and a.s. in $\Omega$ we have

$$
\begin{align*}
\int_{D} a \nabla\left(u-u_{h}^{1}\right) \cdot \nabla v= & \underbrace{\int_{D} f v-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v}_{=: A_{1}} \underbrace{-\varepsilon \int_{D} \sum_{j=1}^{L} Y_{j}\left(a_{0} \nabla U_{j, h}+a_{j} \nabla u_{0, h}\right) \cdot \nabla v}_{=: A_{2}} \\
& \underbrace{-\varepsilon \int_{D}\left(a-a_{0}\right) \nabla u_{1, h} \cdot \nabla v}_{=: A_{3}} \tag{1.67}
\end{align*}
$$

where $A_{1}$ and $A_{2}$ are respectively the residual for $u_{0, h}$ and for $U_{j, h}$, for $j=1, \ldots, L$, while $A_{3}$ is due to the truncation in the expansion (1.19) of $u$. Let us treat each term separately. The first term $A_{1}$ is bounded by (see Section 1.3)

$$
\begin{equation*}
A_{1} \leq C_{1}\left[\sum_{K \in \mathscr{F}_{h}} h_{K}^{2}\left\|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e}\left\|\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2}\right]^{\frac{1}{2}}\|\nabla v\|_{L^{2}(D)} . \tag{1.68}
\end{equation*}
$$

Let us consider now the term $A_{2}$. Since $\int_{D} w_{j, h} \cdot \nabla v_{h}=0$ for all $v_{h} \in V_{h}$, we have

$$
\begin{align*}
A_{2} & =-\varepsilon \int_{D} \sum_{j=1}^{L} Y_{j} w_{j, h} \cdot \nabla\left(\nu-\mathscr{I}_{h} v\right) \\
& =\varepsilon \sum_{K \in \mathscr{T}_{h}} \int_{K}\left(\sum_{j=1}^{L} Y_{j} \nabla \cdot w_{j, h}\right)\left(\nu-\mathscr{I}_{h} \nu\right)+\varepsilon \sum_{e \in \mathscr{F}_{h}} \int_{e}\left[\sum_{j=1}^{L} Y_{j} w_{j, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\left(\nu-\mathscr{I}_{h} v\right) \\
& \leq C_{2}\left(\sum_{K \in \mathscr{F}_{h}} \varepsilon^{2} h_{K}^{2}\left\|\sum_{j=1}^{L} Y_{j} \nabla \cdot w_{j, h}\right\|_{L^{2}(K)}^{2}+\sum_{e \in \mathscr{F}_{h}} \varepsilon^{2} h_{e}\left\|\left[\sum_{j=1}^{L} Y_{j} w_{j, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2}\right)^{\frac{1}{2}}\|\nabla v\|_{L^{2}(D)}, \tag{1.69}
\end{align*}
$$

where $C_{2}$ depends only on the interpolation constants that appear in (1.26) and (1.28). Finally, we estimate the last term $A_{3}$. We have

$$
\begin{align*}
A_{3} & =-\varepsilon \int_{D}\left(\varepsilon \sum_{j=1}^{L} Y_{j} a_{j}\right) \nabla\left(\sum_{i=1}^{L} Y_{i} U_{i, h}\right) \cdot \nabla v=-\varepsilon^{2} \int_{D} \sum_{i, j=1}^{L} Y_{i} Y_{j} a_{j} \nabla U_{i, h} \cdot \nabla v \\
& \leq \varepsilon^{2}\left\|\sum_{i, j=1}^{L} Y_{i} Y_{j} a_{j} \nabla U_{i, h}\right\|_{L^{2}(D)}\|\nabla v\|_{L^{2}(D)} \tag{1.70}
\end{align*}
$$

Since $a$ is bounded from below by $a_{\text {min }}$, combining (1.67) with (1.68), (1.69) and (1.70) with $v=u(\cdot, \mathbf{Y}(\omega))-u_{h}^{1}(\cdot, \mathbf{Y}(\omega)) \in H_{0}^{1}(D)$ yields a.s. in $\Omega$

$$
\begin{aligned}
\left\|\nabla\left(u-u_{h}^{1}\right)\right\|_{L^{2}(D)} \leq & \frac{\sqrt{3}}{a_{\min }}\left[C_{1}^{2}\left(\sum_{K \in \mathscr{T}_{h}} h_{K}^{2}\left\|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e}\left\|\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2}\right)\right. \\
& +C_{2}^{2}\left(\sum_{K \in \mathscr{T}_{h}} \varepsilon^{2} h_{K}^{2}\left\|\sum_{j=1}^{L} Y_{j} \nabla \cdot w_{j, h}\right\|_{L^{2}(K)}^{2}+\sum_{e \in \mathscr{T}_{h}} \varepsilon^{2} h_{e}\left\|\left[\sum_{j=1}^{L} Y_{j} w_{j, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2}\right) \\
& \left.+\varepsilon^{4}\left\|\sum_{i, j=1}^{L} Y_{i} Y_{j} a_{j} \nabla U_{i, h}\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}},
\end{aligned}
$$

using the inequality $(a+b+c) \leq \sqrt{3}\left(a^{2}+b^{2}+c^{2}\right)^{\frac{1}{2}}$. To conclude the proof, it only remains to take the expected value on both sides of the square of this last inequality. By linearity of the expected value, we can consider the three terms of the right-hand side separately. The first term is a deterministic quantity and thus, taking the expected value on it has no effect. For the two other terms, we just have to evaluate $\mathbb{E}\left[Y_{i} Y_{j}\right]$ for $1 \leq i, j \leq L$ and $\mathbb{E}\left[Y_{i} Y_{j} Y_{k} Y_{l}\right]$ for $1 \leq i, j, k, l \leq L$. Since the random variables are assumed to be independent, with zero mean and unit variance, we have $\mathbb{E}\left[Y_{i} Y_{j}\right]=\delta_{i j}$ and

$$
\mathbb{E}\left[Y_{i} Y_{j} Y_{k} Y_{l}\right]=\left\{\begin{array}{cl}
\mathbb{E}\left[Y_{j}^{4}\right] & \text { if } i=j=k=l \\
1 & \text { if the indices are pairwise equal } \\
0 & \text { otherwise } .
\end{array}\right.
$$

Let us write

$$
B:=\sum_{i, j, k, l=1}^{L} Y_{i} Y_{j} Y_{k} Y_{l} a_{j} a_{k} \nabla U_{i, h} \cdot \nabla U_{l, h},
$$

which we split into three parts $B_{1}$ (all indices are equal), $B_{2}$ (two pairs of indices) and $B_{3}$ (remaining indices). Thanks to the linearity of expectation, we have $\mathbb{E}[B]=\mathbb{E}\left[B_{1}\right]+\mathbb{E}\left[B_{2}\right]+\mathbb{E}\left[B_{3}\right]$. First, we can notice that $\mathbb{E}\left[B_{3}\right]=0$. Moreover, the contribution to $\mathbb{E}[B]$ when $i=j=k=l$ is

$$
\mathbb{E}\left[B_{1}\right]=\sum_{i=1}^{L} a_{i}^{2}\left|\nabla U_{i, h}\right|^{2} \mathbb{E}\left[Y_{i}^{4}\right] .
$$

Let us consider now all the cases when we have pairwise equal pairs of indices. Out of 4 indices, there are three different ways to form two pairs of indices, namely ( $j=k, i=l),(j=i, k=l)$
and ( $j=l, k=i$ ). Since the two last cases lead to the same result, we get

$$
\mathbb{E}\left[B_{2}\right]=\sum_{\substack{i, j=1 \\ i \neq j}}^{L} a_{j}^{2}\left|\nabla U_{i, h}\right|^{2}+2 \sum_{\substack{i, j=1 \\ i \neq j}}^{L} a_{i} a_{j} \nabla U_{i, h} \cdot \nabla U_{j, h}
$$

Altogether, we finally get

$$
\mathbb{E}[B]=\sum_{i=1}^{L} a_{i}^{2}\left|\nabla U_{i, h}\right|^{2} \mathbb{E}\left[Y_{i}^{4}\right]+\sum_{\substack{i, j=1 \\ i \neq j}}^{L}\left[a_{i}^{2}\left|\nabla U_{j, h}\right|^{2}+2 a_{i} a_{j} \nabla U_{i, h} \cdot \nabla U_{j, h}\right],
$$

which concludes the proof.

### 1.4.2 Generalization

Suppose now that the random solution $u$ of problem (1.11) is expanded with respect to $\varepsilon$ up to order $N \in \mathbb{N}$, see (1.19). For $1 \leq n \leq N$, let us write

$$
\begin{equation*}
u_{n}(\mathbf{x}, \mathbf{Y}(\omega))=\sum_{j_{1}, j_{2}, \ldots, j_{n}=1}^{L} U_{j_{1} j_{2} \cdots j_{n}}(\mathbf{x}) Y_{j_{1}}(\omega) Y_{j_{2}}(\omega) \cdots Y_{j_{n}}(\omega) \tag{1.71}
\end{equation*}
$$

the $n^{\text {th }}$ term in the expansion. The $L^{n}$ functions $U_{j_{1} j_{2} \cdots j_{n}}$ are obtained by solving for $j_{1}, j_{2}, \ldots, j_{n}=$ $1, \ldots, L$ the deterministic problem

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{j_{1}}(\mathbf{x}) \nabla U_{j_{2} \cdots j_{n}}(\mathbf{x})+a_{0}(\mathbf{x}) \nabla U_{j_{1} \cdots j_{n}}(\mathbf{x})\right) & =0 & \mathbf{x} \in D  \tag{1.72}\\
U_{j_{1} \cdots j_{n}}(\mathbf{x}) & =0 & \mathbf{x} \in \partial D
\end{array}\right.
$$

using the solutions $U_{j_{2} \cdots j_{n}}, j_{2}, \ldots, j_{n}=1, \ldots, L$, obtained for the $(n-1)^{t h}$ order term. Proceeding as in Sections 1.3 and 1.4.1, it is easy to show that the error due to the truncation in the expansion of $u$ is of order $\varepsilon^{N+1}$. More precisely, we have for any $v \in H_{0}^{1}(D)$ and almost surely

$$
\begin{equation*}
\int_{D} a \nabla\left(u-\sum_{n=0}^{N} \varepsilon^{n} u_{n}\right) \cdot \nabla v=-\varepsilon^{N+1} \sum_{j_{0}, j_{1}, \ldots, j_{N}=1}^{L} Y_{j_{0}} Y_{j_{1}} \cdots Y_{j_{N}} \int_{D} a_{j_{0}} \nabla U_{j_{1} j_{2} \cdots j_{N}} \cdot \nabla v \tag{1.73}
\end{equation*}
$$

Since $Y_{j}, j=1, \ldots, L$ are bounded, in particular they have bounded $2(N+1)^{t h}$ moment. When the various deterministic functions are approximated using finite elements, if the solution is regular enough in physical space then the error $u-\sum_{n=0}^{N} \varepsilon^{n} u_{n, h}$ in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm is of order

$$
h+\varepsilon h+\varepsilon^{2} h+\ldots+\varepsilon^{N} h+\varepsilon^{N+1} .
$$

The error in $\mathscr{O}\left(\varepsilon^{n} h\right), 0 \leq n \leq N$, corresponds to the error made when the functions $U_{j_{1} \cdots j_{n}}\left(u_{0}\right.$ for $n=0$ ) are replaced by their FE approximation $U_{j_{1} \cdots j_{n}, h}$ (resp. $u_{0, h}$ ). An a posteriori error estimate can thus easily be obtained as follows. First, the term in $\mathscr{O}(h)$, which corresponds to the residual for $u_{0, h}$, is obtained by estimating $\int_{D}\left(f v-a_{0} \nabla u_{0, h} \cdot \nabla v\right)$, see (1.47). For the term in $\mathscr{O}\left(h \varepsilon^{n}\right), n=1, \ldots, N$, it suffices to estimate for $j_{1}, \ldots, j_{n}=1, \ldots, L$ the residual defined for any
$v \in H_{0}^{1}(D)$ by

$$
\left\langle\mathscr{R}\left(U_{j_{1} \cdots j_{n}, h}\right), v\right\rangle:=\int_{D}\left(a_{j_{1}} \nabla U_{j_{2} \cdots j_{n}, h}+a_{0} \nabla U_{j_{1} \cdots j_{n}, h}\right) \cdot \nabla v,
$$

where $\langle\cdot, \cdot\rangle$ denotes the duality pairing bracket. For an explicit error estimate, computable up to multiplicative interpolation constants, we finally need to express the expectation of the product of $n$ random variables $\mathbb{E}\left[Y_{j_{1}} \cdots Y_{j_{n}}\right]$ for all combinations of indices and for $n=$ $1, \ldots, 2(N+1)$. More precisely, we can show the following result.

Proposition 1.4.2. Let $u$ be the weak solution of problem (1.11) and $u_{h}^{N}=\sum_{n=0}^{N} \varepsilon^{n} u_{n, h}$, where $u_{n, h}$ is the FE approximation of $u_{n}$ given by (1.71). There exist $N+1$ constants $C_{n}>0, n=$ $0,1, \ldots, N$, depending only on the constants in (1.26) and (1.28) such that

$$
\begin{equation*}
\mathbb{E}\left[\left\|\nabla\left(u-u_{h}^{N}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{\sqrt{N+2}}{a_{\min }}\left[C_{0} \eta_{0}^{2}+\sum_{n=1}^{N} C_{n} \eta_{n}^{2}+\eta_{N+1}^{2}\right]^{\frac{1}{2}} \tag{1.74}
\end{equation*}
$$

with

$$
\begin{aligned}
\eta_{0}^{2}= & \sum_{K} h_{K}^{2}\left\|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e} h_{e}\left\|\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2}, \\
\eta_{n}^{2}= & \varepsilon^{2 n} \mathbb{E}\left[\sum_{K} h_{K}^{2}\left\|\sum_{j_{1}, \ldots, j_{n}=1}^{L} Y_{j_{1}} \cdots Y_{j_{n}} \nabla \cdot w_{j_{1} \cdots j_{n}, h}\right\|_{L^{2}(K)}^{2}\right. \\
& \left.+\sum_{e} h_{e}\left\|\left[\sum_{j_{1}, \ldots, j_{n}=1}^{L} Y_{j_{1}} \cdots Y_{j_{n}} w_{j_{1} \cdots j_{n}, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2}\right] \\
\eta_{N+1}^{2}= & \varepsilon^{2(N+1)} \mathbb{E}\left[\left\|\sum_{j_{0}, j_{1}, \ldots, j_{N}=1}^{L} Y_{j_{0}} Y_{j_{1}} \cdots Y_{j_{N}} a_{j_{0}} \nabla U_{j_{1} \cdots j_{N}, h}\right\|_{L^{2}(D)}^{2}\right]
\end{aligned}
$$

where

$$
w_{j_{1} \cdots j_{n}, h}:=a_{j_{1}} \nabla U_{j_{2} \cdots j_{n}, h}+a_{0} \nabla U_{j_{1} \cdots j_{n}, h} \quad j_{1}, \ldots, j_{n}=1, \ldots, L
$$

Proceeding similarly, this generalization can also be applied to the other error estimates we obtained in Section 1.3. Finally, notice that the constant $\sqrt{N+2}$ that appears in (1.74) can be avoided thanks to the triangle inequality for the $L_{P}^{2}(\Omega)$ norm, yielding an upper bound of the form $a_{\text {min }}^{-1}\left(C_{0} \eta_{0}+\ldots+C_{N} \eta_{N}+\eta_{N+1}\right)$. The same holds for all the error estimates obtained in Sections 1.3 and 1.4.1.

### 1.5 Extension to nonlinear problems

Keeping the same notations as in the previous sections, we are now interested in solving problems of the form:
find $u: D \times \Omega \rightarrow \mathbb{R}$ such that almost surely:

$$
\left\{\begin{array}{rll}
F(a, u) & =0 & \text { in } D  \tag{1.75}\\
u & =0 & \text { on } \partial D
\end{array}\right.
$$

where $F$ is a smooth nonlinear mapping that depends on the uncertain input $a$ given by (1.12). Again, the random solution $u$ is expanded with respect to $\varepsilon$ up to a certain order

$$
u(\mathbf{x}, \mathbf{Y}(\omega))=u_{0}(\mathbf{x})+\varepsilon u_{1}(\mathbf{x}, \mathbf{Y}(\omega))+\mathscr{O}\left(\varepsilon^{2}\right)
$$

Formally, we have

$$
F(a, u)=F\left(a_{0}, u_{0}\right)+D_{a} F\left(a_{0}, u_{0}\right)\left(a-a_{0}\right)+D_{u} F\left(a_{0}, u_{0}\right)\left(u-u_{0}\right)+\mathscr{O}\left(\varepsilon^{2}\right)
$$

where $D_{a}$ and $D_{u}$ denote the Fréchet derivatives with respect to $a$ and $u$ respectively, the deterministic part $u_{0}$ of $u$ is the solution of the (nonlinear) problem

$$
\left\{\begin{align*}
F\left(a_{0}, u_{0}\right) & =0 \quad \text { in } D  \tag{1.76}\\
u_{0} & =0 \quad \text { on } \partial D
\end{align*}\right.
$$

while the $U_{j}$ in $u_{1}=\sum_{j=1}^{L} Y_{j} U_{j}$ can be found by solving the (linear) problems

$$
\left\{\begin{align*}
D_{a} F\left(a_{0}, u_{0}\right)\left(a_{j}\right)+D_{u} F\left(a_{0}, u_{0}\right)\left(U_{j}\right) & =0 \quad \text { in } D  \tag{1.77}\\
U_{j} & =0 \quad \text { on } \partial D, \quad j=1, \ldots, L
\end{align*}\right.
$$

We can directly see one of the advantages of expanding the solution as proposed here, namely that a single nonlinear problem must be solved to find $u_{0}$, the other problems being linear. A new FE solver corresponding to (1.77) has to be implemented to approximate the $U_{j}, j=$ $1, \ldots, L$.

In the case of quasi-linear problems, the error analysis is very similar to the linear case considered in Section 1.1. Indeed, under certain conditions such as well-posedness of the problem, only the part of the estimate corresponding to the residual error in the physical space has to be changed in the a posteriori estimate of the error between $u$ and $u_{0, h}$ in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm. For instance, let us consider problem (1.75) with

$$
\begin{equation*}
F(a(\mathbf{x}, \omega), u(\mathbf{x}, \omega)):=-\operatorname{div}(a(\mathbf{x}, \omega) \nabla u(\mathbf{x}, \omega))+u^{3}(\mathbf{x}, \omega)-f(\mathbf{x}) \tag{1.78}
\end{equation*}
$$

This well-posed problem has a unique solution in $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ and we can show the following a posteriori error estimate for $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$, where $u_{0, h} \in V_{h}$ is the deterministic solution of

$$
\begin{equation*}
\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v_{h}+\int_{D} u_{0, h}^{3} v_{h}=\int_{D} f v_{h} \quad \forall v_{h} \in V_{h} \tag{1.79}
\end{equation*}
$$

Proposition 1.5.1. Let $u$ be the weak solution of problem (1.75) with $F$ given by (1.78), and let $u_{0, h}$ be the solution of (1.79). There exists a constant $C>0$ depending only on the constants in
(1.26) and (1.28) such that

$$
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{C}{a_{\min }}\left[\eta_{1}^{2}+\eta_{2}^{2}\right]^{\frac{1}{2}}
$$

with

$$
\begin{aligned}
\eta_{1}^{2} & :=\sum_{K \in \mathscr{T}_{h}} h_{K}^{2} \int_{K}\left(f-u_{0, h}^{3}+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{\mathbf{e}}\right]_{\mathbf{n}_{\mathbf{e}}}^{2} \\
\eta_{2}^{2} & :=\varepsilon^{2} \int_{D} \sum_{j=1}^{L} a_{j}^{2}\left|\nabla u_{0, h}\right|^{2}
\end{aligned}
$$

Proof. Since the proof is very similar to the one of Proposition 1.3.5, we only give the key ingredients here. First, for any $v \in V$ we have almost surely

$$
\int_{D} a \nabla\left(u-u_{0, h}\right) \cdot \nabla v=\int_{D}\left(f-u_{0, h}^{3}\right) v-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla v-\int_{D}\left(u^{3}-u_{0, h}^{3}\right) v
$$

Then, for $v=u-u_{0, h}$ the last term in the above equality is non-positive. Indeed, using that

$$
u^{3}-u_{0, h}^{3}=\int_{0}^{1} 3\left(u_{0, h}+t\left(u-u_{0, h}\right)\right)^{2}\left(u-u_{0, h}\right) d t
$$

we get

$$
-\int_{D}\left(u^{3}-u_{0, h}^{3}\right)\left(u-u_{0, h}\right)=-\int_{D} \int_{0}^{1} 3\left(u_{0, h}+t\left(u-u_{0, h}\right)\right)^{2}\left(u-u_{0, h}\right)^{2} \leq 0
$$

Therefore, this term can be omitted since we are looking for an upper bound of the error.

Another example is the following. Let $k>0$ be such that $\frac{k C_{P}^{2}}{a_{\min }}<1$, or in other words $\frac{k C_{p}^{2}}{a_{m i n}} \leq 1-\delta$ for any $\delta \in(0,1)$. If we take

$$
\begin{equation*}
F(a(\mathbf{x}, \omega), u(\mathbf{x}, \omega)):=-\operatorname{div}(a(\mathbf{x}, \omega) \nabla u(\mathbf{x}, \omega))-g(u(\mathbf{x}, \omega)) \tag{1.80}
\end{equation*}
$$

in problem (1.75), where $g$ is a Lipschitz function with Lipschitz constant $k$, then we can show the well-posedness of the problem and the following a posteriori error estimate for the error $u-u_{0, h}$, where $u_{0, h} \in V_{h}$ is the deterministic solution of

$$
\begin{equation*}
\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v_{h}=\int_{D} g\left(u_{0, h}\right) v_{h} \quad \forall v_{h} \in V_{h} . \tag{1.81}
\end{equation*}
$$

Proposition 1.5.2. Let $u$ be the weak solution of problem (1.75) with $F$ given by (1.80), and let $u_{0, h}$ be the solution of (1.81). There exists a constant $C>0$, depending only on $\delta$ and the constants in (1.26) and (1.28), such that

$$
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{C}{a_{\min }}\left[\eta_{1}^{2}+\eta_{2}^{2}\right]^{\frac{1}{2}}
$$

with

$$
\begin{aligned}
& \eta_{1}^{2}:=\sum_{K \in \mathscr{T}_{h}} h_{K}^{2} \int_{K}\left(g\left(u_{0, h}\right)+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{\mathrm{e}}\right]_{\mathbf{n}_{\mathbf{e}}}^{2} \\
& \eta_{2}^{2}:=\varepsilon^{2} \int_{D} \sum_{j=1}^{L} a_{j}^{2}\left|\nabla u_{0, h}\right|^{2} .
\end{aligned}
$$

Proof. Again, we only give the key ingredients of the proof. First, for any $v \in V$ we have almost surely

$$
\begin{align*}
\int_{D} a \nabla\left(u-u_{0, h}\right) \cdot \nabla v= & \underbrace{\int_{D} g\left(u_{0, h}\right) v-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla v}_{=: A(v)}  \tag{1.82}\\
& -\int_{D}\left(g(u)-g\left(u_{0, h}\right)\right) v .
\end{align*}
$$

With $v=u-u_{0, h}$, the last term is bounded by

$$
\begin{equation*}
-\int_{D}\left(g(u)-g\left(u_{0, h}\right)\right)\left(u-u_{0, h}\right) \leq k C_{P}^{2}\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2} \tag{1.83}
\end{equation*}
$$

Since

$$
a_{m i n}\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2} \leq \int_{D} a\left|\nabla\left(u-u_{0, h}\right)\right|^{2},
$$

taking (1.83) to the left-hand side of (1.82) and using $k C_{P}^{2} \leq a_{\text {min }}(1-\delta)$ yield

$$
a_{\text {min }} \delta\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2} \leq A\left(u-u_{0, h}\right) .
$$

A bound on $A\left(u-u_{0, h}\right)$, which contains the residual for $u_{0}$ and a term of order $\varepsilon$, is found proceeding exactly as in the proof of Proposition 1.3.5.

The constant $C$ that appears in the error estimate of Proposition 1.5.2 is of order $\delta^{-1}$, and thus explodes when $\delta$ tends to zero, i.e. when $\frac{k C_{p}^{2}}{a_{\text {min }}}$ is close to one. In practise, it is usual to restrict the analysis to Lipschitz function with Lipschitz constant $k$ such that $k \leq \frac{a_{\text {min }}}{2 C_{P}^{2}}$, so that $\delta \geq \frac{1}{2}$.
Finally, let us consider an example where the uncertain coefficient is associated to the nonlinear term, namely the problem (1.75) with

$$
\begin{equation*}
F(a(\mathbf{x}, \omega), u(\mathbf{x}, \omega))=-\Delta u(\mathbf{x}, \omega)+a(\mathbf{x}, \omega) u^{3}(\mathbf{x}, \omega)-f(\mathbf{x}) \tag{1.84}
\end{equation*}
$$

In this case, we can show the well-posedness of the problem and the following a posteriori error estimate in $H_{0}^{1}(D)$-norm in physical space for the first order approximation $u \approx u_{0, h}$, where $u_{0, h}$ is the solution of

$$
\begin{equation*}
\int_{D} \nabla u_{0, h} \cdot \nabla v_{h}+\int_{D} a_{0} u_{0, h}^{3} v_{h}=\int_{D} f v_{h} \quad \forall v_{h} \in V_{h} . \tag{1.85}
\end{equation*}
$$

Proposition 1.5.3. Let $u$ be the weak solution of problem (1.75) with $F$ given by (1.84), and let $u_{0, h}$ be the solution of (1.85). There exists a constant $C>0$ depending only on the constants in (1.26) and (1.28) such that

$$
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq C\left[\eta_{1}^{2}+\eta_{2}^{2}\right]^{\frac{1}{2}}
$$

with

$$
\begin{aligned}
\eta_{1}^{2} & :=\sum_{K \in \mathscr{T}_{h}} h_{K}^{2} \int_{K}\left(f+\Delta u_{0, h}-a_{0} u_{0, h}^{3}\right)^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e} \int_{e}\left[\nabla u_{0, h} \cdot \mathbf{n}_{\mathbf{e}}\right]_{\mathbf{n}_{\mathbf{e}}}^{2} \\
\eta_{2}^{2} & :=\varepsilon^{2} \int_{D} \sum_{j=1}^{L} a_{j}^{2} u_{0, h}^{6} .
\end{aligned}
$$

Proof. The proof is based on the relations

$$
\int_{D} \nabla\left(u-u_{0, h}\right) \cdot \nabla v=\int_{D} f v-\int_{D} a_{0} u_{0, h}^{3} v-\int_{D} \nabla u_{0, h} \cdot \nabla v-\int_{D}\left(a u^{3}-a_{0} u_{0, h}^{3}\right) v
$$

and

$$
-\int_{D}\left(a u^{3}-a_{0} u_{0, h}^{3}\right) v=-\int_{D} a \int_{0}^{1} 3\left(u_{0, h}+t\left(u-u_{0, h}\right)\right)^{2}\left(u-u_{0, h}\right) d t v-\int_{D}\left(a-a_{0}\right) u_{0, h}^{3} v
$$

Since $a$ is positive, the first term of the right-hand side of the last equality is less or equal to zero for $v=u-u_{0, h}$.

### 1.6 Computational costs

We perform here a comparison of the computational costs between the SC-FEM method [7, 124] and the one presented here, called perturbation method in the sequel, when comparable accuracy is reached. Briefly, the SC-FEM applied to the model problem (1.11) consists, given a set of (collocation) points $\left\{\mathbf{y}_{k} \in \Gamma, k=1, \ldots, N_{c}\right\}$, in finding $u_{h}\left(\cdot, \mathbf{y}_{k}\right) \in V_{h}$ such that

$$
\int_{D} a\left(\mathbf{x}, \mathbf{y}_{k}\right) \nabla u_{h}\left(\mathbf{x}, \mathbf{y}_{k}\right) \cdot \nabla v_{h}(\mathbf{x}) d \mathbf{x}=\int_{D} f(\mathbf{x}) v_{h}(\mathbf{x}) d \mathbf{x} \quad \forall v_{h} \in V_{h}
$$

for $k=1, \ldots, N_{c}$ and building a global polynomial approximation

$$
u_{h, N_{c}}(\mathbf{x}, \mathbf{y})=\sum_{k=1}^{N_{c}} u_{h}\left(\mathbf{x}, \mathbf{y}_{k}\right) \psi_{k}(\mathbf{y})
$$

for appropriate multivariate polynomials $\left\{\psi_{k}\right\}_{k=1}^{N_{c}}$. Since the FEM is used to approximate the physical space in both methods (stochastic collocation and perturbation), we use the same mesh for the discretization of $D$. For a comparable statistical error, say an error with convergence rate of order $\varepsilon^{2}$, we take $N=1$ in the expansion (1.19) of $u$ for the perturbation
method and use a sparse grid of level 1 for the SC method, based either on Clenshaw-Curtis (see [51]) or Gaussian abscissas. The construction of the sparse grid interpolant of level 1 is briefly described in the following. We refer to [65, 97, 124] for more details and the general construction of sparse grid of arbitrary level. First, the sparse grid interpolant of level 0 of a function $f(\mathbf{y})$, denoted $S_{0} f$, is simply the evaluation of the function at $\left(y_{1}^{0}, \ldots, y_{L}^{0}\right)$, where $y_{j}^{0}$ is the unique interpolation point in direction $j$. Next, for each variable $y_{j}$, we define the sequence of interpolation points at level $i \geq 1$ by $\left\{y_{j, k}^{i}, k=1, \ldots, m(i)\right\}$, where the number of collocation points $m(i)$ can be taken for instance as

$$
m(i)=i+1 \quad \text { or } \quad m(i)=\left\{\begin{array}{cc}
1 & \text { if } i=0 \\
2^{i}+1 & \text { if } i \geq 1
\end{array}\right.
$$

The former choice for $m$ corresponds to a total degree (TD) approximation space while the latter corresponds to a Smolyak one (see [11]). Notice that compared to the articles mentioned above, the level index $i$ starts here at 0 instead of 1 . We define then the one dimensional (Lagrange) interpolation operator in direction $j$ at level $i=1$ by

$$
\mathscr{U}_{j}^{1} f\left(y_{1}, \ldots, y_{L}\right):=\sum_{k=1}^{m(1)} f\left(y_{1}^{0}, \ldots, y_{j-1}^{0}, y_{j, k}^{1}, y_{j+1}^{0}, \ldots, y_{L}^{0}\right)\left(\prod_{l=1, l \neq k}^{m(1)} \frac{y_{j}-y_{j, l}^{1}}{y_{j, k}^{1}-y_{j, l}^{1}}\right)
$$

which is a polynomial of degree $m(1)-1$ in the direction $j$ and constant in all other directions. Finally, the level 1 sparse grid interpolant is defined as

$$
S_{1} f:=S_{0} f+\sum_{j=1}^{L}\left(\mathscr{U}_{j}^{1} f-S_{0} f\right)=(1-L) S_{0} f+\sum_{j=1}^{L} \mathscr{U}_{j}^{1} f
$$

which is nothing else than the sum of the level 0 sparse grid interpolant and the details in each direction.

Remark 1.6.1. It can be proved that the SC approximation computed with a sparse grid of level 1 indeed yields an error of order $\varepsilon^{2}$, using for instance a scaling argument together with the fact that $S_{1}$ is exact for any polynomial of (total) degree at most 1 (see [18]). More generally, we can show that a sparse grid of level l yields an error of order $\varepsilon^{l+1}$ for the choice $m(i)=i$, while for the second choice of $m$ it is of order $\varepsilon^{l+k+1}$, where $k=0$ if $l<L$ and $k=l-L+1$ otherwise.

The type of points in each direction is chosen according to the distribution of the random variables. Note that the use of Clenshaw-Curtis points, which are the extrema of Chebyshev polynomials and which are suitable for uniformly distributed random variables, and Smolyak sparse grid leads to nested set of abscissas. However, since only sparse grids of level 1 are considered, there is no real advantage to consider hierarchical sparse grids. In both cases $m(1)=2$ and Gauss-Legendre abscissas and $m(1)=3$ and Clenshaw-Curtis abscissas, referred to as SC1 and SC2 in the following, the sparse grid of level 1 consists of $2 L+1$ collocation points (due to the use of nested set of abscissas in each direction for SC2).

Let $W_{l}$, respectively $W_{n l}$, denote the work to solve once a given linear, respectively nonlinear, problem. Moreover, let $W_{\tilde{l}}$ denote the work to solve the linear problem for $U_{j}$ associated with the nonlinear one, see (1.77). Table 1.1 contains the computational costs for the SC-FEM and the perturbation method. Notice that the work to construct the sparse grid is not taken into account.

|  | linear problem | nonlinear problem |
| :---: | :---: | :---: |
| SC-FEM | $(2 L+1) \cdot W_{l}$ | $(2 L+1) \cdot W_{n l}$ |
| perturbation method | $(L+1) \cdot W_{l}$ | $W_{n l}+L \cdot W_{\tilde{l}}$ |

Table 1.1: Computational costs for the SC-FEM and the perturbation method.

The perturbation method presents no real advantage for solving linear problems since the costs for both methods differ only by a factor 2 . The situation is different when a nonlinear problem is considered. Indeed, when using the SC method, we need to solve as many nonlinear problems as collocation points, i.e. $2 L+1$ problems, whereas only one nonlinear problem needs to be solved for the perturbation method. The $L$ remaining problems, to compute the $U_{j}, j=1, \ldots, L$, are linear and so usually much cheaper to solve. However, one should invest extra effort to derive by hand the Fréchet derivatives and implement the problems solved by the $U_{j}, j=1, \ldots, L$.

### 1.7 Numerical results

This section is devoted to illustration and validation of the theoretical results obtained in the previous sections. We start with the analysis of 1D problems, analysing first the convergence rate for various errors and norms and presenting, next, algorithms which adaptively refine the (physical) mesh to balance the two sources of error: the physical space discretization and the uncertainty. We present then two 2D examples and conclude this section with a comparison with the stochastic collocation method in term of computational costs when solving linear and nonlinear problems.

### 1.7.1 1D problems

Let $D=(0,1)$. In what follows, the true errors in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ and $L_{P}^{2}\left(\Omega ; L^{2}(D)\right)$ norms have been accurately approximated with the standard Monte Carlo method, with a sample of size $K=10000$, i.e. for $V=H_{0}^{1}(D)$ or $L^{2}(D)$ we approximate

$$
\|v\|_{L_{P}^{2}(\Omega ; V)} \approx\left(\frac{1}{K} \sum_{k=1}^{K}\left\|v\left(\cdot, \mathbf{y}_{k}\right)\right\|_{V}^{2}\right)^{\frac{1}{2}} \quad \forall v \in L_{P}^{2}(\Omega ; V)
$$

where $\left\{\mathbf{y}_{k}\right\} \in \Gamma$ are i.i.d realizations of the random vector $\mathbf{Y}$. With this choice for the sample size, the variance of the estimation of the error for all the considered values of $h$ and $\varepsilon$ is at most $10^{-5}$ the estimated error. In what follows, whenever we refer to error it should be understood that the true error has been accurately computed by the Monte Carlo procedure. Since the exact random solution of the problems considered below is not known, the error is computed with respect to a reference solution computed on a fine uniform mesh for $D$, namely with a mesh-grid of length $h_{r e f}=2^{-12}$. Notice that if we take a FE space of mesh size $h=h_{r e f}$, then only the statistical error is considered. Finally, all the involved integrals are evaluated numerically with sufficiently accurate quadrature formulas that permit to neglect the effect of quadrature.

Let us first consider $L=50$ random variables $Y_{j}, j=1, \ldots, L$, which can take the values $\pm 1$ with probability $\frac{1}{2}$. Such discrete random variables have zero mean, unit variance and unit fourth moment. Similarly to what is done in [124], we take a diffusion coefficient of the form

$$
\begin{equation*}
a(x, \mathbf{Y}(\omega))=1+\varepsilon \sum_{j=1}^{L} \frac{\cos (2 \pi j x)}{(\pi j)^{2}} Y_{j}(\omega) \tag{1.86}
\end{equation*}
$$

which is similar to a (truncated) Karhunen-Loève expansion with eigenvalues of order $\frac{1}{j^{4}}$. With this choice of random diffusion coefficient, we have $1-\frac{\varepsilon}{6} \leq a(x, y) \leq 1+\frac{\varepsilon}{6}$. We take $\varepsilon \in[0,4]$ which guarantee property (A1) with $a_{\min }=\frac{1}{3}$ and $a_{\max }=\frac{5}{3}$. Finally, we consider two different right-hand sides, namely

$$
\begin{equation*}
f_{1}(x)=1 \quad \text { and } \quad f_{2}(x)=72\left(1-72(x-0.5)^{2}\right) e^{-36(x-0.5)^{2}} \tag{1.87}
\end{equation*}
$$

The latter corresponds to the exact solution $u_{0}(x)=e^{-36(x-0.5)^{2}}-e^{-9}$ for problem (1.20) while it is $u_{0}(x)=0.5 x(1-x)$ for the case $f=f_{1}$.

Error in $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$-norm
We consider first the error measured in $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$-norm. We show in Figure 1.1 the convergence rate of the error $u-u_{0, h}$ with respect to $2^{-9} \leq h \leq 2^{-3}$ for $\varepsilon=32 h$, along with the a posteriori estimator based on (1.43) and (1.44). Based on this result, we can see that a division of $h$ and $\varepsilon$ by two halves the error, which is in agreement with the convergence of $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ in $\mathscr{O}(h+\varepsilon)$ predicted by the foregoing error analysis. Moreover, for the two cases $f_{1}$ and $f_{2}$, the gap between the error and the estimator is of about 1.6 and 2.8 , respectively, which is comprised between the effectivity index of the stochastic error estimator (1) and the spatial error estimator (3.46), see below for details. Concerning the convergence rate of the second order approximation, we present in Figure 1.2 the error between $u$ and $u_{h}^{1}$ with respect to $2^{-3} \leq \varepsilon \leq 2$ for $h=\varepsilon^{2} / 32$. This result confirms the convergence in $\mathscr{O}\left(\varepsilon^{2}\right)$ of the stochastic truncation predicted by (1.63), when the exact solution is approximated by $u_{0}+\varepsilon u_{1}$.

The error estimators depicted in Figures 1.1 and 1.2 do not take into account the unknown

Chapter 1. Elliptic model problems with random diffusion coefficient


Figure 1.1: Convergence orders for problem (1.11) with $f=f_{1}$ (left) and $f=f_{2}$ (right). Log log scale plot of the error between $u$ and $u_{0, h}$ in $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$-norm w.r.t $h$ with $\varepsilon=32 h$.


Figure 1.2: Convergence orders for problem (1.11) with $f=f_{1}$ (left) and $f=f_{2}$ (right). Log $\log$ scale plot of the error between $u$ and $u_{h}^{1}$ in $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$-norm w.r.t $\varepsilon$ with $h=\varepsilon^{2} / 32$.
constants due to interpolation error that appear in (1.42) and (1.63). These constants can be estimated numerically as follows to obtain a sharp error estimator: consider the problem $-u_{0}^{\prime \prime}=f$ with $f$ such that the exact solution is known, for instance $f=f_{1}$ or $f=f_{2}$, and define $1 / C_{H_{0}^{1}}:=3.46 \approx \eta_{1} /\left\|u_{0}-u_{0, h}\right\|_{H_{0}^{1}(D)}$ for $h$ small enough. This estimation can be done once for all since $C_{H_{0}^{1}}$ does not depend on the input data. We define then $\tilde{\eta}:=\left(C_{H_{0}^{1}}^{2} \eta_{1}^{2}+\eta_{2}^{2}\right)^{\frac{1}{2}}$ as an estimator for the error $\left\|u-u_{0, h}\right\|_{L_{p}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$. We will say that $\tilde{\eta}$ is a good approximation of the error if the ratio $\tilde{\eta} /\left\|u-u_{0, h}\right\|_{L_{p}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ remains between $a_{\text {min }}$ and $a_{\text {max }}$. Since in the considered case the ratio $a_{\text {max }} / a_{\text {min }}$ tends to 1 as $\varepsilon$ goes to 0 , we expect the effectivity index of the estimator $\tilde{\eta}$ to approach 1 as $\varepsilon$ gets smaller. We give in Tables 1.2 and 1.3 the results obtained when the constant $C_{H_{0}^{1}}$ is considered. In Table 1.2, the mesh size is fixed to $h=2^{-7}$ while in Table 1.3 we fix $\varepsilon=0.25$. In both cases, the ratio of the estimator $\tilde{\eta}$, which contains the estimated constant $C_{H_{0}^{1}}$, over the error is close to one.

| $\varepsilon$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $\eta_{2}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 4 | $1.2167 \mathrm{e}-1$ | $2.2579 \mathrm{e}-3$ | $9.1996 \mathrm{e}-2$ | $9.2024 \mathrm{e}-2$ | 0.75632 |
| 4 | $4.9276 \mathrm{e}-2$ | $2.2579 \mathrm{e}-3$ | $4.5998 \mathrm{e}-2$ | $4.6054 \mathrm{e}-2$ | 0.93461 |
|  | 1 | $2.3460 \mathrm{e}-2$ | $2.2579 \mathrm{e}-3$ | $2.2999 \mathrm{e}-2$ | $2.3110 \mathrm{e}-2$ |
|  |  |  |  |  |  |
| 0.5 | $1.1760 \mathrm{e}-2$ | $2.2579 \mathrm{e}-3$ | $1.1500 \mathrm{e}-2$ | $1.1719 \mathrm{e}-2$ | 0.99505 |
| 0.25 | $6.1805 \mathrm{e}-3$ | $2.2579 \mathrm{e}-3$ | $5.7498 \mathrm{e}-3$ | $6.1772 \mathrm{e}-3$ | 0.99947 |
| 0.125 | $3.6545 \mathrm{e}-3$ | $2.2579 \mathrm{e}-3$ | $2.8749 \mathrm{e}-3$ | $3.6556 \mathrm{e}-3$ | 1.00031 |


|  | $\varepsilon$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $\eta_{2}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\begin{aligned} & \mathbb{N} \\ & \stackrel{N}{I I} \\ & \underset{4}{2} \end{aligned}$ | 4 | $9.5591 \mathrm{e}-1$ | $6.4347 \mathrm{e}-2$ | 7.8646e-1 | $7.8909 \mathrm{e}-1$ | 0.82548 |
|  | 2 | $4.1806 \mathrm{e}-1$ | $6.4347 \mathrm{e}-2$ | 3.9323e-1 | $3.9846 \mathrm{e}-1$ | 0.95312 |
|  | 1 | $2.0916 \mathrm{e}-1$ | $6.4347 \mathrm{e}-2$ | 1.9661e-1 | $2.0688 \mathrm{e}-1$ | 0.98910 |
|  | 0.5 | $1.1782 \mathrm{e}-1$ | $6.4347 \mathrm{e}-2$ | $9.8307 \mathrm{e}-2$ | $1.1749 \mathrm{e}-1$ | 0.99720 |
|  | 0.25 | $8.0974 \mathrm{e}-2$ | $6.4347 \mathrm{e}-2$ | $4.9154 \mathrm{e}-2$ | 8.0973e-2 | 0.99999 |
|  | 0.125 | $6.8769 \mathrm{e}-2$ | $6.4347 \mathrm{e}-2$ | $2.4577 \mathrm{e}-2$ | $6.8881 \mathrm{e}-2$ | 1.00163 |

Table 1.2: Error $\left\|u-u_{0, h}\right\|_{L_{p}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$, estimators $\eta_{1}, \eta_{2}$ and $\tilde{\eta}$ and ratio $\tilde{\eta} /\left\|u-u_{0, h}\right\|_{L_{p}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ for $h=2^{-7}$ and various $\varepsilon$ for both cases $f_{1}$ and $f_{2}$.

The same observation holds for the approximation $u \approx u_{0, h}+\varepsilon u_{1, h}$ taking $C_{1}=C_{2}=C_{H_{0}^{1}}^{2}$ in (1.63) and for the generalization (1.74) with $C_{i}=C_{H_{0}^{1}}^{2}$ for $i=0, \ldots, N$, see Table 1.4 where the case $u \approx u_{0, h}+\varepsilon u_{1, h}$ is presented for the case $f=f_{2}$. Recall that $\eta_{1}, \eta_{2}$ and $\eta_{3}$ are given in (1.64), (1.65) and (1.66), respectively, and here $\tilde{\eta}:=\left(C_{H_{0}^{2}}^{2} \eta_{1}^{2}+C_{H_{0}^{1}}^{2} \eta_{2}^{2}+\eta_{2}^{3}\right)^{\frac{1}{2}}$.

## Error in $L_{P}^{2}\left(\Omega ; L^{2}(D)\right)$-norm

We consider now the error $u-u_{0, h}$ in $L_{P}^{2}\left(\Omega ; L^{2}(D)\right)$-norm. According to the theoretical result, we should get a convergence of order $h^{2}$ for $\varepsilon=C h^{2}$. Figure 1.3, which contains the plot of the error and estimator based on (1.54) and (1.55) for $C=32$ and $2^{-6} \leq h \leq 2^{-2}$, confirms that this

|  | $N$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $\eta_{2}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\stackrel{41}{4}$ | 8 | $3.6575 \mathrm{e}-2$ | $3.6127 \mathrm{e}-2$ | 5.5801e-3 | 3.6556e-2 | 0.99946 |
|  | 16 | $1.8951 \mathrm{e}-2$ | $1.8064 \mathrm{e}-2$ | $5.7030 \mathrm{e}-3$ | $1.8942 \mathrm{e}-2$ | 0.99955 |
|  | 32 | $1.0712 \mathrm{e}-2$ | $9.0318 \mathrm{e}-3$ | $5.7384 \mathrm{e}-3$ | $1.0701 \mathrm{e}-2$ | 0.99890 |
|  | 64 | 7.3076e-3 | $4.5159 \mathrm{e}-3$ | $5.7475 \mathrm{e}-3$ | $7.3094 \mathrm{e}-3$ | 1.00024 |
|  | 128 | $6.1765 \mathrm{e}-3$ | $2.2580 \mathrm{e}-3$ | $5.7498 \mathrm{e}-3$ | $6.1772 \mathrm{e}-3$ | 1.00011 |
|  | 256 | $5.8822 \mathrm{e}-3$ | 1.1290e-3 | $5.7503 \mathrm{e}-3$ | $5.8601 \mathrm{e}-3$ | 0.99625 |


|  | $N$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $\eta_{2}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\underset{4}{\mathbb{N}_{11}}$ | 8 | $9.7697 \mathrm{e}-1$ | $1.0441 \mathrm{e}-0$ | 4.6189e-2 | $1.0451 \mathrm{e}-0$ | 1.06977 |
|  | 16 | $5.1089 \mathrm{e}-1$ | $5.1478 \mathrm{e}-1$ | 4.8261e-2 | $5.1704 \mathrm{e}-1$ | 1.01204 |
|  | 32 | 2.6109e-1 | $2.5739 \mathrm{e}-1$ | $4.8942 \mathrm{e}-2$ | $2.6200 \mathrm{e}-1$ | 1.00349 |
|  | 64 | $1.3766 \mathrm{e}-1$ | $1.2869 \mathrm{e}-1$ | $4.9112 \mathrm{e}-2$ | $1.3775 \mathrm{e}-1$ | 1.00066 |
|  | 128 | 8.0919e-2 | $6.4347 \mathrm{e}-2$ | $4.9154 \mathrm{e}-2$ | $8.0973 \mathrm{e}-2$ | 1.00066 |
|  | 256 | $5.8787 \mathrm{e}-2$ | $3.2174 \mathrm{e}-2$ | $4.9164 \mathrm{e}-2$ | $5.8756 \mathrm{e}-2$ | 0.99946 |

Table 1.3: Error $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$, estimators $\eta_{1}, \eta_{2}$ and $\tilde{\eta}$ and ratio $\tilde{\eta} /\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ for $\varepsilon=0.25$ and various $h=1 / N$ for both cases $f_{1}$ and $f_{2}$.

|  | $\varepsilon$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $C_{H_{0}^{1}} \eta_{2}$ | $\eta_{3}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\begin{aligned} & 0 \\ & 1 \\ & N \\ & 11 \\ & \\ & \end{aligned}$ | 4 | 3.5236e-1 | $8.0434 \mathrm{e}-3$ | $1.8607 \mathrm{e}-3$ | $2.7488 \mathrm{e}-1$ | $2.7500 \mathrm{e}-1$ | 0.78044 |
|  | 2 | $7.3380 \mathrm{e}-2$ | $8.0434 \mathrm{e}-3$ | $9.3037 \mathrm{e}-4$ | $6.8719 \mathrm{e}-2$ | $6.9194 \mathrm{e}-2$ | 0.94295 |
|  | 1 | $1.9054 \mathrm{e}-2$ | $8.0434 \mathrm{e}-3$ | $4.6519 \mathrm{e}-4$ | $1.7180 \mathrm{e}-2$ | $1.8975 \mathrm{e}-2$ | 0.99586 |
|  | 0.5 | $8.9126 \mathrm{e}-3$ | $8.0434 \mathrm{e}-3$ | $2.3259 \mathrm{e}-4$ | $4.2949 \mathrm{e}-3$ | $9.1213 \mathrm{e}-3$ | 1.02341 |
|  | 0.25 | $7.8616 \mathrm{e}-3$ | $8.0434 \mathrm{e}-3$ | 1.1630e-4 | $1.0737 \mathrm{e}-3$ | 8.1156e-3 | 1.03231 |
|  | 0.125 | $7.7840 \mathrm{e}-3$ | $8.0434 \mathrm{e}-3$ | 5.8148e-5 | 2.6843e-4 | 8.0481e-3 | 1.03393 |


|  | $N$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $C_{H_{0}^{1}} \eta_{2}$ | $\eta_{3}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\stackrel{\rightharpoonup}{11}$ | 32 | $2.6010 \mathrm{e}-1$ | $2.5739 \mathrm{e}-1$ | $1.4792 \mathrm{e}-2$ | 1.6950e-2 | $2.5837 \mathrm{e}-1$ | 0.99337 |
|  | 64 | $1.3107 \mathrm{e}-1$ | $1.2869 \mathrm{e}-1$ | $7.4313 \mathrm{e}-3$ | $1.7121 \mathrm{e}-2$ | $1.3004 \mathrm{e}-1$ | 0.99219 |
|  | 128 | $6.7384 \mathrm{e}-2$ | $6.4347 \mathrm{e}-2$ | $3.7201 \mathrm{e}-3$ | $1.7165 \mathrm{e}-2$ | $6.6701 \mathrm{e}-2$ | 0.98987 |
|  | 256 | $3.6796 \mathrm{e}-2$ | $3.2174 \mathrm{e}-2$ | 1.8606e-3 | $1.7176 \mathrm{e}-2$ | $3.6519 \mathrm{e}-2$ | 0.99246 |
|  | 512 | $2.3761 \mathrm{e}-2$ | $1.6087 \mathrm{e}-2$ | 9.3036e-4 | $1.7179 \mathrm{e}-2$ | $2.3554 \mathrm{e}-2$ | 0.99128 |
|  | 1024 | $1.9131 \mathrm{e}-2$ | $8.0434 \mathrm{e}-3$ | $4.6519 \mathrm{e}-4$ | $1.7180 \mathrm{e}-2$ | $1.8975 \mathrm{e}-2$ | 0.99188 |

Table 1.4: Error and estimators for the approximation $u \approx u_{0, h}+\varepsilon u_{1, h}$ with $h$ fixed (top) and $\varepsilon$ fixed (bottom) for the case $f=f_{2}$.
is the case. Similarly to the error in $H_{0}^{1}(D)$-norm, the constant due to interpolation error could


Figure 1.3: Convergence orders for problem (1.11) with $f=f_{1}$ (left) and $f=f_{2}$ (right). Log log scale plot of the error between $u$ and $u_{0, h}$ in $L_{P}^{2}\left(\Omega ; L^{2}(D)\right)$-norm w.r.t $h$ with $\varepsilon$ fixed to $32 h^{2}$.
be estimated numerically once for all following the same procedure as above. However, even with a sharp estimation of such constant, there is no guarantee that the estimator is efficient though it has the correct convergence rate. We see two reasons for that. First of all, there are no proofs, to our knowledge, that the part of the estimator due to the uncertainty $\left(\eta_{2}\right)$ is a lower bound for the error in $L^{2}(D)$-norm, mainly due to the use of the Poincaré inequality. Considering $h=h_{r e f}$, the estimator over estimates the error by a factor of about 4.2 for $f=f_{1}$ and 9 for $f=f_{2}$, showing that the constant multiplying $\eta_{2}$ does depend on $f$. Moreover, the constant $C_{1}$ in (1.53) depends in an implicit way on the uniform bound for $a$ and $\nabla a$ (see Remark 1.3.10).

## Different setup

Similar results are obtained when other input data are considered. For instance, let us consider independent uniformly distributed random variables in $[-\sqrt{3}, \sqrt{3}]$. In this case, the random variables still have zero mean and unit variance but $\mathbb{E}\left[Y_{j}^{4}\right]=\frac{9}{5}$. This only modifies the part $\eta_{3}$ in the a posteriori error estimate (1.63) for $\left\|u-u_{h}^{1}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$. Moreover, we also modify the functions $a_{j}$ considering here

$$
\begin{equation*}
a(x, \mathbf{Y}(\omega))=1+\varepsilon \sum_{j=1}^{50} \frac{\cos (8 \pi j x) \sin (2 \pi j x)}{(\pi j)^{2}} Y_{j}(\omega) \tag{1.88}
\end{equation*}
$$

for the random diffusion coefficient. Notice that this choice satisfies $1-\frac{\sqrt{3} \varepsilon}{6} \leq a(x, y) \leq 1+\frac{\sqrt{3} \varepsilon}{6}$. We give in Figure 1.4 some realizations of $a$ and the corresponding solution for the case $\varepsilon=1$ and $f=f_{2}$ defined in (1.87).

The results obtained when the constant $C_{H_{0}^{1}}=1 / 3.46$ is taken into account are given in Table 1.5. First, the mesh size is fixed to $h=1 / N=2^{-8}$ and $\varepsilon$ varies and then, we set $\varepsilon=0.5$ and consider various partitions of $[0,1]$. When $h$ is fixed, the error decreases linearly with respect


Figure 1.4: Five realizations of the random diffusion coefficient $a$ given in (1.88) with $\varepsilon=1$ (left) and the corresponding solution for $f=f_{2}$ (right).
to $\varepsilon$ until the FE error is no longer negligible. The same observation holds when $\varepsilon$ is fixed and $h$ varies. In both cases, the effectivity index of the error estimator $\tilde{\eta}=\left(C_{H_{0}^{1}}^{2} \eta_{1}^{2}+\eta_{2}^{2}\right)^{\frac{1}{2}}$ is close to one.

| $\varepsilon$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $\eta_{2}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 2 | $3.2152 \mathrm{e}-1$ | $3.2174 \mathrm{e}-2$ | $3.0331 \mathrm{e}-1$ | $3.0501 \mathrm{e}-1$ | 0.94866 |
| 1 | $1.5541 \mathrm{e}-1$ | $3.2174 \mathrm{e}-2$ | $1.5165 \mathrm{e}-1$ | $1.5503 \mathrm{e}-1$ | 0.99754 |
| 0.5 | $8.1168 \mathrm{e}-2$ | $3.2174 \mathrm{e}-2$ | $7.5827 \mathrm{e}-2$ | $8.2371 \mathrm{e}-2$ | 1.01482 |
| 0.25 | $4.9399 \mathrm{e}-2$ | $3.2174 \mathrm{e}-2$ | $3.7914 \mathrm{e}-2$ | $4.9725 \mathrm{e}-2$ | 1.00659 |
| 0.125 | $3.7192 \mathrm{e}-2$ | $3.2174 \mathrm{e}-2$ | $1.8957 \mathrm{e}-2$ | $3.7343 \mathrm{e}-2$ | 1.00406 |
| 0.0625 | $3.3432 \mathrm{e}-2$ | $3.2174 \mathrm{e}-2$ | $9.4784 \mathrm{e}-3$ | $3.3541 \mathrm{e}-2$ | 1.00325 |


| $N$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $\eta_{2}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 8 | $9.7920 \mathrm{e}-1$ | $1.0441 \mathrm{e}-0$ | $9.7528 \mathrm{e}-2$ | $1.0487 \mathrm{e}-0$ | 1.07093 |
| 16 | $5.1403 \mathrm{e}-1$ | $5.1478 \mathrm{e}-1$ | $7.6937 \mathrm{e}-2$ | $5.2050 \mathrm{e}-1$ | 1.01258 |
| 32 | $2.6726 \mathrm{e}-1$ | $2.5739 \mathrm{e}-1$ | $7.5506 \mathrm{e}-2$ | $2.6824 \mathrm{e}-1$ | 1.00365 |
| 64 | $1.4900 \mathrm{e}-1$ | $1.2869 \mathrm{e}-1$ | $7.5726 \mathrm{e}-2$ | $1.4932 \mathrm{e}-1$ | 1.00217 |
| 128 | $9.8399 \mathrm{e}-2$ | $6.4347 \mathrm{e}-2$ | $7.5805 \mathrm{e}-2$ | $9.9433 \mathrm{e}-2$ | 1.01051 |
| 256 | $8.1817 \mathrm{e}-2$ | $3.2174 \mathrm{e}-2$ | $7.5827 \mathrm{e}-2$ | $8.2371 \mathrm{e}-2$ | 1.00676 |

Table 1.5: Error $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$, estimators $\eta_{1}, \eta_{2}$ and $\tilde{\eta}$ and ratio $\tilde{\eta} /\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ for $h=2^{-8}$ (top) and $\varepsilon=0.5$ (bottom).

## Adaptive algorithm

We propose here adaptive algorithms to determine, for a given $\varepsilon$, a mesh for $D$ that balances the two sources of error. The convergence rate of the error in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm with respect to $h$ for uniform refinements and for the first, second and third order approximation for several given (fixed) values of $\varepsilon$ is depicted in Figure 1.5 in the case $f=1$ and $a$ given in
(1.86). First, we can notice that a better accuracy is reached when $u$ is approximated by $u_{h}^{2}$ than


Figure 1.5: Convergence rate for problem (1.11) with $f=f_{1}$ for $\varepsilon=4,1,0.25,0.0625$. $\log \log$ scale plot of the error in $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$-norm w.r.t $h$.
with $u_{h}^{1}$, which in turn provides a better approximation than only the deterministic part $u_{0, h}$. This observation holds except for coarse meshes where the FE error is dominating yielding comparable accuracy in all cases. Moreover, the global approximation error remains constant for mesh sizes smaller than a critical value $h_{0}$ of the mesh-size. Any further refinement of the mesh below this value should thus be avoided since it would not improve the global approximation error, being dominated by the stochastic error.

Based on this observation, it is interesting to determine how fine the mesh should be to get a comparable error in $h$ and $\varepsilon$. More precisely, for a given $\varepsilon$ and for the approximation $u \approx u_{0, h}$, we would like to find a mesh for $D$ such that

$$
\begin{equation*}
\frac{T-1}{T} \eta_{2} \leq \eta_{1} \leq \frac{T+1}{T} \eta_{2} \tag{1.89}
\end{equation*}
$$

for a given preset tolerance $T>1$, where $\eta_{1}$ and $\eta_{2}$ are given by (1.43) and (1.44), respectively. Notice that in all what follows, $\eta_{1}$ can be replaced by $C_{H_{0}^{1}} \eta_{1}$ if the estimated constant $C_{H_{0}^{1}}$ is at disposal, so that the correct balance of the two sources of error is considered. Moreover, we mention that the choice of the law of the $Y_{j}, j=1, \ldots, L$, is irrelevant here as long as $\mathbb{E}\left[Y_{j}\right]=0$ and $\operatorname{Var}\left(Y_{j}\right)=1$. Indeed, the error estimator $\eta_{2}$ given in (1.44) is valid under these conditions irrespectively of the law of $Y_{j}$ and only the solution $u_{0, h}$ is computed.

## Uniform refinement

The adaptation can be done in 1D using Algorithm 1 given below, where $N_{h}+1$ denotes the number of discretization points in $[0,1]$.

```
Algorithm 1 find \(h=N_{h}^{-1}\) such that (1.89) holds
Require: \(N_{\text {init }}\) and \(T\)
Ensure: mesh-size \(h\) which yield comparable accuracy in \(h\) and \(\varepsilon\)
    \(N_{h}=N_{\text {init }}\)
    Compute \(u_{0, h}\) on the uniform partition \(x_{i}=i h, h=N_{h}^{-1}, i=0,1, \ldots, N_{h}\)
    Compute \(\eta_{1}\) and \(\eta_{2}\) according to (1.43) and (1.44)
    if \(\frac{T-1}{T} \leq \frac{\eta_{1}}{\eta_{2}} \leq \frac{T+1}{T}\) then
        stop
    else
        if \(\frac{\eta_{1}}{\eta_{2}}<\frac{T-1}{T}\) then
            \(N_{h} \leftarrow\left\lfloor\frac{N_{h}}{2}\right\rfloor\) (mesh too fine)
        else
            \(N_{h} \leftarrow 2 N_{h}\) (mesh too coarse)
        end if
        go to 2 .
    end if
```

Applying Algorithm 1 to our problem for $T=2$ and various given $\varepsilon$, we get the results presented in Table 1.6.

|  | $f_{1}$ |  |  | $f_{2}$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\varepsilon$ | $N_{h}$ | $\eta_{1}$ | $\eta_{2}$ | $N_{h}$ | $\eta_{1}$ | $\eta_{2}$ |
| 1 | 32 | 0.03125 | 0.02295 | 128 | 0.22264 | 0.19661 |
| 0.5 | 64 | 0.01563 | 0.01149 | 256 | 0.11132 | 0.09833 |
| 0.25 | 128 | 0.00781 | 0.00575 | 512 | 0.05566 | 0.04917 |
| 0.125 | 256 | 0.00391 | 0.00288 | 1024 | 0.02783 | 0.02458 |
| 0.0625 | 512 | 0.00195 | 0.00144 | 2048 | 0.01392 | 0.01229 |

Table 1.6: Value of $h=N_{h}^{-1}$ with respect to $\varepsilon$ such that (1.89) holds with $T=2$.

We mention that if T is large, i.e. $\frac{T-1}{T}$ is close to $\frac{T+1}{T}$, the algorithm might not converge due to an oscillation of the ratio $\frac{\eta_{1}}{\eta_{2}}$ below the lower bound $\frac{T-1}{T}$ and above the upper bound $\frac{T+1}{T}$ in two consecutive steps. Such behaviour will be observed if no uniform partition of $D$ satisfies (1.89). Moreover, notice that with Algorithm 1, only refinement or only coarsening is performed, depending on the initial number $N_{\text {init }}$ of subintervals.

## Non-uniform refinement

Algorithm 1 given above only uses uniform refinement or coarsening. Of course, adaptive refinements can be considered as well exploiting the local nature of the estimator $\eta_{1}$, which can indeed be written as

$$
\begin{equation*}
\eta_{1}^{2}=\sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2} \quad \text { with } \quad \eta_{K}^{2}=h_{K}^{2} \int_{K}\left(f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)^{2}+\frac{1}{2} \sum_{e \subset \partial K} h_{e} \int_{e}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{\mathbf{e}}\right]_{\mathbf{n}_{\mathbf{e}}}^{2} \tag{1.90}
\end{equation*}
$$

taking into account that each edge is then counted twice.
Remark 1.7.1. The factor $\frac{1}{2}$ could in fact be replaced by $\frac{1}{4}$ if we do not split the summation over the elements and the edges in the derivation of the error estimate in (1.47), namely if we consider an element point of view. Indeed, we can use the fact that for any $v \in H_{0}^{1}(D)$ and any $v_{h} \in V_{h}$ we have

$$
\int_{D} f v-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v=\sum_{K \in \mathscr{T}_{h}}\left[\int_{K}\left(f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)\left(v-v_{h}\right)+\int_{\partial K} \frac{1}{2}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\left(v-v_{h}\right)\right] .
$$

Recall that in 1D, for a partition $0=x_{0}<x_{1}<\ldots<x_{N_{h}}=1$, the error estimator $\eta_{1}$ reads

$$
\eta_{1}^{2}=\sum_{i=0}^{N_{h}-1} \eta_{1, i}^{2} \quad \text { with } \quad \eta_{1, i}^{2}=h_{i}^{2}\left\|f+\left(a_{0} u_{0, h}^{\prime}\right)^{\prime}\right\|_{L^{2}\left(x_{i}, x_{i+1}\right)}^{2}
$$

The goal being to satisfy (1.89), a first possibility is to require that

$$
\begin{equation*}
B_{\mathrm{inf}}:=\left(\frac{T-1}{T}\right)^{2} \eta_{2}^{2} \frac{1}{N_{h}} \leq \eta_{1, i}^{2} \leq\left(\frac{T+1}{T}\right)^{2} \eta_{2}^{2} \frac{1}{N_{h}}=: B_{\text {sup }} \quad \forall i=0, \ldots, N_{h}-1 \tag{1.91}
\end{equation*}
$$

Another sufficient condition for (1.89) to hold is to impose that

$$
\begin{equation*}
B_{\mathrm{inf}}:=\left(\frac{T-1}{T}\right)^{2} \eta_{2}^{2} \frac{h_{i}}{|D|} \leq \eta_{1, i}^{2} \leq\left(\frac{T+1}{T}\right)^{2} \eta_{2}^{2} \frac{h_{i}}{|D|}=: B_{\text {sup }} \quad \forall i=0, \ldots, N_{h}-1 \tag{1.92}
\end{equation*}
$$

using the fact that $\sum_{i=0}^{N_{h}-1} h_{i}=|D|$. The criterion (1.91) imposes an equidistribution of the error, enforcing a comparable value of the local error estimator on each subinterval regardless of its length. In the second strategy (1.92), the repartition of the error is weighted by $h_{i}$. This is commonly used in a time-adaptivity framework so that the solution does not need to be computed until the final time before adapting the time step.

We give in Algorithm 2 an adaptive procedure which find a (non-uniform) partition of $D$ for which (1.89) holds. The idea is to check for each subinterval $\left[x_{i}, x_{i+1}\right], i=0, \ldots, N_{h}-1$, of the current partition of $D$ if the local error estimator $\eta_{1, i}$ satisfies the criterion (1.91) or (1.92). If it is too large, then we should refine the interval $\left[x_{i}, x_{i+1}\right]$, for instance by adding its midpoint, while a coarsening should be done if it is too small.

To better appreciate the behaviour of the non-uniform adaptation, we test Algorithm 2 with a

```
Algorithm 2 adaptive algorithm with non-uniform partition
Require: \(T\) and initial partition \(\mathscr{T}_{h}=\left\{x_{i}, i=0, \ldots, N_{h}-1\right\}\)
Ensure: partition of \(D\) such that (1.89) holds
    Compute \(u_{0, h}\) on \(\mathscr{T}_{h}\)
    Compute \(\eta_{1}\) and \(\eta_{2}\) according to (1.43) and (1.44)
    if \(\frac{T-1}{T} \leq \frac{\eta_{1}}{\eta_{2}} \leq \frac{T+1}{T}\) then
        stop
    else
        for \(i=0, \ldots, N_{h}-1\) do
            if \(\eta_{1, i}^{2}>B_{\text {sup }}\) then
                add the midpoint \(\frac{x_{i}+x_{i+1}}{2}\) to \(\mathscr{T}_{h}\)
            else if \(\eta_{1, i}^{2}<B_{\text {inf }}\) then
                remove the endpoint \(x_{i+1}\) from \(\mathscr{T}_{h}\left(x_{i}\right.\) if \(\left.i=N_{h}-1\right)\)
            end if
        end for
    end if
    go to 1.
```

different forcing term than in the previous sections, keeping the diffusion coefficient $a$ as in (1.86) and all other input data being unchanged. We consider the source term $f$ for which the corresponding solution $u_{0}$ of problem (1.20) is given by ${ }^{1}$

$$
\begin{equation*}
u_{0}(x)=x-\frac{1-e^{x \tau^{-1}}}{1-e^{\tau^{-1}}} \tag{1.93}
\end{equation*}
$$

The solution presents a boundary layer near $x=1$ of width proportional to $\tau$, see Figure 1.6. It is linear on the remaining part of the interval, where only few points are thus sufficient to obtain a good approximation. In the numerical results below, we choose $\tau=0.05$.

We give in Tables 1.7 and 1.8 the results obtained for various values of $\varepsilon$ when using the two adaptive criterion (1.91) and (1.92), respectively. We have denoted by $N_{h}$ the number of subintervals of $D$ (i.e. $N_{h}+1$ is the number of nodes), $h_{\min }=\min _{i} h_{i}$ and $h_{\max }=\max _{i} h_{i}$ are the minimum and maximum mesh sizes, respectively, and iter stands for the number of iterations of the adaptive algorithm. In all cases, we have started the adaptation with the initial partition $\{0,0.5,1\}$.

First, we can see that the number of iterations is similar in both cases and the same holds for the values of the error estimators $\eta_{1}$ and $\eta_{2}$. Moreover, the number of nodes is smaller when criterion (1.91) is used while the maximum subinterval length $h_{\text {max }}$ is in general larger with (1.92). The latter strategy indeed allows to have large subintervals if the corresponding local error estimator is small. This can be seen in Figure 1.6 where the repartition of the nodes is given for various values of $\varepsilon$ and for both criteria (1.91) and (1.91). The continuous line

[^0]| $\varepsilon$ | $N_{h}$ | $h_{\min }$ | $h_{\max }$ | $\eta_{1}$ | $\eta_{2}$ | iter |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 28 | $3.91 \mathrm{e}-3$ | $6.56 \mathrm{e}-1$ | $3.6424 \mathrm{e}-1$ | $3.0474 \mathrm{e}-1$ | 8 |
| 0.5 | 53 | $1.95 \mathrm{e}-3$ | $6.25 \mathrm{e}-1$ | $1.9848 \mathrm{e}-1$ | $1.5272 \mathrm{e}-1$ | 9 |
| 0.1 | 231 | $4.88 \mathrm{e}-4$ | $4.69 \mathrm{e}-1$ | $3.8504 \mathrm{e}-2$ | $3.0727 \mathrm{e}-2$ | 11 |
| 0.05 | 461 | $2.44 \mathrm{e}-4$ | $2.50 \mathrm{e}-1$ | $1.9233 \mathrm{e}-2$ | $1.5164 \mathrm{e}-2$ | 12 |
| 0.01 | 2056 | $6.10 \mathrm{e}-5$ | $1.88 \mathrm{e}-1$ | $4.4334 \mathrm{e}-3$ | $3.0339 \mathrm{e}-3$ | 14 |
| 0.005 | 4119 | $3.05 \mathrm{e}-5$ | $2.81 \mathrm{e}-1$ | $2.2138 \mathrm{e}-3$ | $1.5178 \mathrm{e}-3$ | 15 |
| 0.001 | 25646 | $3.81 \mathrm{e}-6$ | $1.05 \mathrm{e}-1$ | $3.3770 \mathrm{e}-4$ | $3.0304 \mathrm{e}-4$ | 18 |
| 0.0005 | 51292 | $1.91 \mathrm{e}-6$ | $1.05 \mathrm{e}-1$ | $1.6884 \mathrm{e}-4$ | $1.5150 \mathrm{e}-4$ | 19 |
| 0.0001 | 233216 | $4.77 \mathrm{e}-7$ | $5.27 \mathrm{e}-2$ | $3.7686 \mathrm{e}-5$ | $3.0301 \mathrm{e}-5$ | 21 |

Table 1.7: Adaptive partition of $D$ such that (1.89) holds with $T=2$ when criterion (1.91) is used.

| $\varepsilon$ | $N_{h}$ | $h_{\min }$ | $h_{\max }$ | $\eta_{1}$ | $\eta_{2}$ | iter |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 64 | $9.77 \mathrm{e}-04$ | $5.63 \mathrm{e}-01$ | $2.4702 \mathrm{e}-1$ | $3.0720 \mathrm{e}-1$ | 10 |
| 0.5 | 70 | $1.95 \mathrm{e}-03$ | $6.25 \mathrm{e}-01$ | $1.7704 \mathrm{e}-1$ | $1.5273 \mathrm{e}-1$ | 9 |
| 0.1 | 293 | $4.88 \mathrm{e}-04$ | $4.69 \mathrm{e}-01$ | $3.6756 \mathrm{e}-2$ | $3.0704 \mathrm{e}-2$ | 11 |
| 0.05 | 581 | $2.44 \mathrm{e}-04$ | $5.47 \mathrm{e}-01$ | $2.2340 \mathrm{e}-2$ | $1.5356 \mathrm{e}-2$ | 12 |
| 0.01 | 3880 | $3.05 \mathrm{e}-05$ | $2.50 \mathrm{e}-01$ | $3.2449 \mathrm{e}-3$ | $3.0329 \mathrm{e}-3$ | 15 |
| 0.005 | 7741 | $1.53 \mathrm{e}-05$ | $4.38 \mathrm{e}-01$ | $1.7937 \mathrm{e}-3$ | $1.5338 \mathrm{e}-3$ | 16 |
| 0.001 | 33949 | $3.81 \mathrm{e}-06$ | $3.75 \mathrm{e}-01$ | $4.0887 \mathrm{e}-4$ | $3.0531 \mathrm{e}-4$ | 18 |
| 0.0005 | 99606 | $9.54 \mathrm{e}-07$ | $1.88 \mathrm{e}-01$ | $1.6707 \mathrm{e}-4$ | $1.5170 \mathrm{e}-4$ | 20 |
| 0.0001 | 295692 | $4.77 \mathrm{e}-07$ | $2.50 \mathrm{e}-01$ | $4.0904 \mathrm{e}-5$ | $3.0320 \mathrm{e}-5$ | 21 |

Table 1.8: Adaptive partition of $D$ such that (1.89) holds with $T=2$ when criterion (1.92) is used.


Figure 1.6: Repartition of the nodes for $\varepsilon=1$ (top), $\varepsilon=0.1$ (middle) and $\varepsilon=0.01$ (bottom) in the case $T=2$. Left: strategy (1.91), right: strategy (1.92).
represents the exact solution $u_{0}$ given in (1.93).
As we have seen in Tables 1.7 and 1.8, the two methods yield comparable results. The number of nodes for criterion (1.92) is larger but it allows, in general, larger maximum mesh size $h_{\text {max }}$.

Finally, we compare the results of Tables 1.7 and 1.8 with those obtained using a Dörfler [57] bulk-chasing marking commonly used in adaptive finite element method (AFEM), see for instance [42, 114]. To reach the target $\frac{\eta_{1}}{\eta_{2}} \leq \frac{T+1}{T}$, a suitable fraction of the subintervals with highest local error estimator is selected for refinement at each iteration. More precisely, for a given parameter $\theta \in(0,1]$, we select an index set $J \subseteq\left\{0,1, \ldots, N_{h}-1\right\}$ of minimal cardinality such that

$$
\left(\sum_{j \in J} \eta_{1, j}^{2}\right)^{\frac{1}{2}} \geq \theta\left(\sum_{i=0}^{N_{h}-1} \eta_{1, i}^{2}\right)^{\frac{1}{2}}=\theta \eta_{1}
$$

This marking strategy is often referred to as equilibration strategy and yields comparable
results than the so-called maximum strategy, see [119]. Notice that if $\theta$ is closed to 0 , then only few subintervals will be refined at each iteration while choosing $\theta$ close to 1 will generate a set $J$ of large cardinality. In particular, the case $\theta=1$ gives similar results than Algorithm 1 without coarsening, namely all the subintervals are refined at each iteration, except those for which ${ }^{2} \eta_{1, i}=0$. The procedure based on Dörfler marking is described in Algorithm 3. The search for the index $i \in\left\{0, \ldots, N_{h}-1\right\} \backslash J$ with largest $\eta_{1, i}^{2}$ (see line 8) can be achieved by sorting the local estimators $\eta_{1, i}$ in decreasing order before the while loop.

```
Algorithm 3 adaptive algorithm with Dörfler marking
Require: \(T, \theta\) and initial partition \(\mathscr{T}_{h}=\left\{x_{i}, i=0, \ldots, N_{h}-1\right\}\)
Ensure: partition of \(D\) such that \(\frac{\eta_{1}}{\eta_{2}} \leq \frac{T+1}{T}\)
    Compute \(u_{0, h}\) on \(\mathscr{T}_{h}\)
    Compute \(\eta_{1}\) and \(\eta_{2}\) according to (1.43) and (1.44)
    if \(\frac{\eta_{1}}{\eta_{2}} \leq \frac{T+1}{T}\) then
        stop
    else
        \(J=\varnothing\) and \(\vartheta=0\)
        while \(\vartheta<\theta \eta_{1}\) do
            \(J \leftarrow J \cup\{j\}\) with \(j=\arg \max _{i \in\left\{0, \ldots, N_{h}-1\right\} \backslash J} \eta_{1, i}^{2}\)
            \(\vartheta \leftarrow \vartheta+\eta_{1, j}^{2}\)
            add the midpoint \(\frac{x_{j}+x_{j+1}}{2}\) to \(\mathscr{T}_{h}\)
        end while
    end if
    go to 1 .
```

We give in Table 1.9 the results obtained using the Dörfler strategy of Algorithm 3 for the same values of $\varepsilon$ than in Tables 1.7 and 1.8.

| $\varepsilon$ | $N_{h}$ | $h_{\min }$ | $h_{\max }$ | $\eta_{1}$ | $\eta_{2}$ | iter |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 23 | $3.91 \mathrm{e}-03$ | $5.00 \mathrm{e}-01$ | $4.3240 \mathrm{e}-1$ | $3.0736 \mathrm{e}-1$ | 16 |
| 0.5 | 41 | $1.95 \mathrm{e}-03$ | $5.00 \mathrm{e}-01$ | $2.2862 \mathrm{e}-1$ | $1.5371 \mathrm{e}-1$ | 21 |
| 0.1 | 201 | $4.88 \mathrm{e}-04$ | $5.00 \mathrm{e}-01$ | $4.3679 \mathrm{e}-2$ | $3.0756 \mathrm{e}-2$ | 36 |
| 0.05 | 419 | $2.44 \mathrm{e}-04$ | $2.50 \mathrm{e}-01$ | $2.0873 \mathrm{e}-2$ | $1.5164 \mathrm{e}-2$ | 43 |
| 0.01 | 2017 | $3.05 \mathrm{e}-05$ | $2.50 \mathrm{e}-01$ | $4.3602 \mathrm{e}-3$ | $3.0330 \mathrm{e}-3$ | 58 |
| 0.005 | 4177 | $1.53 \mathrm{e}-05$ | $2.50 \mathrm{e}-01$ | $2.1044 \mathrm{e}-3$ | $1.5165 \mathrm{e}-3$ | 65 |
| 0.001 | 19715 | $3.81 \mathrm{e}-06$ | $1.25 \mathrm{e}-01$ | $4.4296 \mathrm{e}-4$ | $3.0306 \mathrm{e}-4$ | 80 |
| 0.0005 | 40705 | $1.91 \mathrm{e}-06$ | $1.25 \mathrm{e}-01$ | $2.1412 \mathrm{e}-4$ | $1.5147 \mathrm{e}-4$ | 87 |
| 0.0001 | 191790 | $4.77 \mathrm{e}-07$ | $6.25 \mathrm{e}-02$ | $4.5111 \mathrm{e}-5$ | $3.0300 \mathrm{e}-5$ | 102 |

Table 1.9: Dörfler strategy such that $\frac{\eta_{1}}{\eta_{2}} \leq \frac{T+1}{T}$ holds with $T=2$ and $\theta=0.5$.

Compared to the results obtained with the two previous adaptive strategies, the Dörfler

[^1]marking procedure requires more iterations but produces a partition of $D$ satisfying $\frac{\eta_{1}}{\eta_{2}} \leq \frac{T+1}{T}$ with fewer nodes. Moreover, this last inequality is tight here which is an expected feature for moderate $\theta$, or when few local error estimators are large compared to the others, since only few subintervals are refined at each step. It is therefore more likely to stop the refinement process when the tolerance is just satisfied. We give in Table 1.10 the results obtained when changing the value of $\theta$.

| $\theta$ | $N_{h}$ | $h_{\min }$ | $h_{\max }$ | $\eta_{1}$ | $\eta_{2}$ | iter |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 0.1 | 1934 | $3.05 \mathrm{e}-5$ | $2.50 \mathrm{e}-1$ | $4.5434 \mathrm{e}-3$ | $3.0330 \mathrm{e}-3$ | 704 |
| 0.4 | 2034 | $3.05 \mathrm{e}-5$ | $2.50 \mathrm{e}-1$ | $4.3241 \mathrm{e}-3$ | $3.0330 \mathrm{e}-3$ | 86 |
| 0.7 | 2202 | $3.05 \mathrm{e}-5$ | $2.50 \mathrm{e}-1$ | $3.9900 \mathrm{e}-3$ | $3.0330 \mathrm{e}-3$ | 31 |
| 0.95 | 2356 | $3.05 \mathrm{e}-5$ | $2.50 \mathrm{e}-1$ | $3.7186 \mathrm{e}-3$ | $3.0330 \mathrm{e}-3$ | 16 |
| 1 | 15872 | $6.10 \mathrm{e}-5$ | $1.22 \mathrm{e}-4$ | $3.8602 \mathrm{e}-3$ | $3.0303 \mathrm{e}-3$ | 14 |

Table 1.10: Dörfler strategy such that $\frac{\eta_{1}}{\eta_{2}} \leq \frac{T+1}{T}$ holds with $T=2$ in the case $\varepsilon=0.01$.

We see that when $\theta$ is small, the number of nodes is small but it requires many iterations of the adaptive process. On the contrary, a large value of $\theta$ yields a partition of $D$ with many nodes obtained with few iterations. Notice that here, all cases but $\theta=1$ yield comparable results in terms of number of nodes, minimal and maximal mesh sizes and estimators. As mentioned above, the case $\theta=1$ yields similar results to those obtained with uniform refinement of the mesh. The only difference lies in the fact that here, the midpoint of a subinterval $\left[x_{i}, x_{i+1}\right]$ is not added if $\eta_{1, i}$ is (numerically) zero. This explain why in Table 1.10 we get $h_{\min } \neq h_{\max }$. If we consider $f_{1}$ or $f_{2}$ as forcing term and $\varepsilon=0.0625$, in which cases no local error estimator $\eta_{1, i}$ vanishes, we get $N_{h}=512$ and $N_{h}=2048$ for $f=f_{1}$ and $f=f_{2}$, respectively, as in Table 1.6.

## Adaptation for higher-order approximation in $\varepsilon$

Here, we give only a sketch of a possible adaptive scheme to achieve an approximate solution with a prescribed accuracy, but we do not provide numerical experiments. As mentioned previously, further mesh refinement should be avoided once the two error estimators $\eta_{1}$ and $\eta_{2}$ are balanced since it would not decrease the total error. The latter can be decreased only by adding more terms in the expansion of $u$. Based on this observation, we can think of a strategy to adaptively increase the degree $N$ in the expansion (1.19) of $u$ together with adaptive mesh refinements for each deterministic term in this expansion. Recall that the estimator for the error $u-u_{h}^{N}=u-\sum_{n=0}^{N} \varepsilon^{n} u_{h, n}$ in the $L_{p}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm reads $a_{m i n}^{-1}\left(C_{H_{0}^{1}} \sum_{n=0}^{N} \eta_{n}+\eta_{N+1}\right)$, see Section 1.4.2. Starting with $N=0$, we find a mesh of $D$ (using Algorithm 2 for instance) such that $C_{H_{0}^{1}} \eta_{0} \approx \eta_{1}$. If the error estimate does not reach the given tolerance, we increase $N$ by one and find a mesh such that $C_{H_{0}^{1}}\left(\eta_{0}+\eta_{1}\right) \approx \eta_{2}$ and proceed then iteratively. Notice that different meshes could be used for the FE approximation of each deterministic part of the solution ( $u_{0}$, $\left.U_{1}, U_{2}, \ldots\right)$.

### 1.7.2 2D problems

The numerical results obtained for the one-dimensional case generalize to problems of higher dimensions. To motivate this statement, we present two numerical examples in 2D. In both cases, the physical domain is $D=(0,1)^{2}$ that we partition using uniform meshes of size $h \sim 1 / n$ for different values of $n$. The true error in the norm $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ is computed via the MonteCarlo method with sample size $K=1000$ and a reference solution computed on the finest mesh considered which corresponds here to $n_{r e f}=2^{8}$.

## First example

We consider first the problem (1.11) with $f(\mathbf{x})=32\left(x_{1}\left(1-x_{1}\right)+x_{2}\left(1-x_{2}\right)\right)$ and

$$
a(\mathbf{x}, \mathbf{Y}(\omega))=1+\varepsilon \sum_{j=1}^{5} \frac{\cos \left(2 \pi j x_{1}\right)+\cos \left(2 \pi j x_{2}\right)}{(\pi j)^{2}} Y_{j}(\omega)
$$

for $\mathbf{x}=\left(x_{1}, x_{2}\right) \in D$, where $Y_{j}, j=1, \ldots, 5$, are uniform random variables in $[-\sqrt{3}, \sqrt{3}]$. In this setting, the exact solution $u_{0}$ for the deterministic case $\varepsilon=0$ is given by $u_{0}(\mathbf{x})=x_{1} x_{2}\left(1-x_{1}\right)(1-$ $x_{2}$ ). The expected value and the standard deviation of $u$ for the case $\varepsilon=0.5$ is given in Figure 1.7.


Figure 1.7: Expected value (left) and standard deviation (right) of the solution with $\varepsilon=0.5$ for the first example.

Similarly to the 1D case, the constant due to interpolation can be estimated numerically, yielding ${ }^{3} C_{H_{0}^{1}}:=1 / 5.7$. We define then $\tilde{\eta}=\left(C_{H_{0}^{1}}^{2} \eta_{1}^{2}+\eta_{2}^{2}\right)^{\frac{1}{2}}$ with $\eta_{1}$ and $\eta_{2}$ given by (1.43) and (1.44), respectively. We report in Table 1.11 the results obtained for $\varepsilon=0.5$ fixed and uniform meshes of various sizes $h \sim 1 / n$ while in Table 1.12, we fix $n=64$ and vary $\varepsilon$.

In Table 1.12, where $\varepsilon$ is fixed and $n$ varies, the error decreases linearly with respect to $h \sim 1 / n$

[^2]| $\varepsilon$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $\eta_{2}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 0.1749 | 0.0604 | 0.1842 | 0.1939 | 1.108 |
| 0.5 | 0.0974 | 0.0604 | 0.0921 | 0.1101 | 1.131 |
| 0.25 | 0.0703 | 0.0604 | 0.0461 | 0.0759 | 1.081 |
| 0.125 | 0.0622 | 0.0604 | 0.0230 | 0.0646 | 1.039 |
| 0.0625 | 0.0597 | 0.0604 | 0.0115 | 0.0615 | 1.029 |

Table 1.11: Error $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$, estimators $\eta_{1}, \eta_{2}$ and $\tilde{\eta}$ and ratio $\tilde{\eta} /\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ with $n=64$ for the first example.

| $n$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $\eta_{2}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 8 | 0.4891 | 0.4649 | 0.0927 | 0.4741 | 0.969 |
| 16 | 0.2551 | 0.2381 | 0.0923 | 0.2554 | 1.001 |
| 32 | 0.1439 | 0.1202 | 0.0922 | 0.1515 | 1.053 |
| 64 | 0.0974 | 0.0604 | 0.0921 | 0.1101 | 1.131 |
| 128 | 0.0833 | 0.0303 | 0.0921 | 0.0969 | 1.164 |

Table 1.12: Error $\left\|u-u_{0, h}\right\|_{L_{p}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$, estimators $\eta_{1}, \eta_{2}$ and $\tilde{\eta}$ and ratio $\tilde{\eta} /\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ with $\varepsilon=0.5$ for the first example.
when $\eta_{2}$ is negligible compared to $C_{H_{0}^{1}} \eta_{1}$. When it is no longer the case, the error continues diminishing with refinement of the mesh but with a smaller rate. The same observation holds for the results of Table 1.11 switching the role of $h$ and $\varepsilon$. Finally, we observe in both cases that the effectivity index of the error estimator $\tilde{\eta}$ that contains the estimated constant $C_{H_{0}^{1}}$ is close to 1 .

## Second example

Let $\left\{\lambda_{i}, \varphi_{i}\right\}$ be the eigenpairs of the Karhunen-Loève expansion of a (1D) Gaussian random field with exponential covariance function $C: D \times D \rightarrow \mathbb{R}$ given by

$$
C\left(x, x^{\prime}\right)=\sigma^{2} e^{\frac{\left|x-x^{\prime}\right|}{L_{c}}}
$$

for which the analytical expression is known, see for instance [67] or [90]. We set the variance $\sigma^{2}$ and the correlation length $L_{c}$ to $\sigma^{2}=L_{c}=1$ and we consider the random diffusion coefficient $a$ obtained by tensorization

$$
a(\mathbf{x}, \mathbf{Y}(\omega))=1+\varepsilon \sum_{i=1}^{3} \sum_{k=1}^{3} \sqrt{\lambda_{i} \lambda_{k}} \varphi_{i}\left(x_{1}\right) \varphi_{k}\left(x_{2}\right) Y_{i k}(\omega)=1+\varepsilon \sum_{j=1}^{9} a_{j}(\mathbf{x}) Y_{j}(\omega)
$$

where $Y_{j}, j=1, \ldots, 9$, are uniform random variables in $[-\sqrt{3}, \sqrt{3}]$. Finally, we choose here $f(\mathbf{x})=10 \sin \left(2 \pi\left(x_{1}+x_{2}\right)\right)$ for the forcing term. We give in Figure 1.8 the functions $\sqrt{\lambda_{i} \lambda_{k}} \varphi_{i}\left(x_{1}\right) \varphi_{k}\left(x_{2}\right)$
for $i, k=1,2,3$. Notice that we can choose the global index $j$ so that $\lambda_{j}=\lambda_{i} \lambda_{k}$ is nondecreasing but it is irrelevant here. Indeed, we do not perform a truncation on $j$ and so an ordering to keep the more relevant functions is not required.


Figure 1.8: Plot of the functions $a_{j}, j=1, \ldots, 9$, constructed by tenzorization of onedimensional KL functions.

The expected value and the standard deviation of $u$ for the case $\varepsilon=0.5$ is given in Figure 1.9.
Finally, the results for a fixed $n=128$ and a fixed $\varepsilon=0.05$ are given in Tables 1.13 and 1.14 , respectively.

The conclusions for this second example are the same as in the previous example.

### 1.7.3 Comparison with the stochastic collocation method

We finally illustrate the findings of Section 1.6 concerning the computation costs for the SCFEM and the perturbation method. We consider the linear problem (1.11) and the nonlinear problem (1.75) with $F$ given by (1.78). In both cases, homogeneous Dirichlet boundary condition are considered and we assume that the random variables $Y_{j}, j=1, \ldots, L$, that appear in the characterization (1.86) of $a$ are uniform random variables in $[-\sqrt{3}, \sqrt{3}]$. We compare the computation time to solve the two problems with accuracy of order 2 in $\varepsilon$. Such accuracy is reached when we consider a sparse grid of level 1 for the SC-FEM method and the second


Figure 1.9: Expected value (left) and standard deviation (right) of the solution with $\varepsilon=0.1$ for the second example.

| $\varepsilon$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $\eta_{2}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 0.1 | 0.0623 | 0.0201 | 0.0605 | 0.0637 | 1.0227 |
| 0.05 | 0.0356 | 0.0201 | 0.0302 | 0.0363 | 1.0195 |
| 0.025 | 0.0245 | 0.0201 | 0.0151 | 0.0252 | 1.0269 |
| 0.0125 | 0.0210 | 0.0201 | 0.0076 | 0.0215 | 1.0263 |
| 0.00625 | 0.0200 | 0.0201 | 0.0038 | 0.0205 | 1.0274 |

Table 1.13: Error $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$, estimators $\eta_{1}, \eta_{2}$ and $\tilde{\eta}$ and ratio $\tilde{\eta} /\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ with $n=128$ for the second example.

| $n$ | error | $C_{H_{0}^{1}} \eta_{1}$ | $\eta_{2}$ | $\tilde{\eta}$ | $\tilde{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 8 | 0.3397 | 0.2762 | 0.0260 | 0.2774 | 0.8167 |
| 16 | 0.1804 | 0.1527 | 0.0291 | 0.1555 | 0.8616 |
| 32 | 0.0941 | 0.0791 | 0.0300 | 0.0848 | 0.9007 |
| 64 | 0.0527 | 0.0401 | 0.0302 | 0.0505 | 0.9577 |
| 128 | 0.0358 | 0.0201 | 0.0302 | 0.0367 | 1.0259 |

Table 1.14: Error $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$, estimators $\eta_{1}, \eta_{2}$ and $\tilde{\eta}$ and ratio $\tilde{\eta} /\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ with $\varepsilon=0.05$ for the second example.
order approximation $u \approx u_{0, h}+\varepsilon u_{1, h}$ for the perturbation method. Note that $u_{1, h}=\sum_{j=1}^{L} U_{j, h} Y_{j}$ where $U_{j, h}$ for $j=1, \ldots, L$ is the solution of

$$
\int_{D} a_{0} \nabla U_{j, h} \cdot \nabla v_{h}+\int_{D} 3 u_{0, h}^{2} U_{j, h} v_{h}=-\int_{D} a_{j} \nabla u_{0, h} \cdot v_{h} \quad \forall v_{h} \in V_{h} .
$$

when problem (1.75) is considered. Finally, we use the same physical space discretization for both methods, namely a uniform partition with $h=2^{-12}$. With this choice of mesh size, the work to solve the $(2 L+1)$ problems dominates the one needed to construct the grid. The computational time to solve both problems with respect to the number of random variables $L$ is given in Figure 1.10.


Figure 1.10: Time to solve the linear problem (1.11) and the nonlinear problem (1.78) with accuracy of order 2 in $\varepsilon$ using the SC-FEM and the perturbation method.

As predicted in Section 1.6, the perturbation method presents no real advantage in terms of computation time over the stochastic collocation one, since it is only twice faster. This factor 2 comes from the fact that the perturbation method requires the solution of $L+1$ problems, while $2 L+1$ problems need to be solve in the stochastic collocation method. The situation is different for nonlinear problems. In this case, the perturbation method is significantly faster than the stochastic collocation one. Indeed, only one nonlinear problem and $L$ linear problems need to be solve for the former, to obtain respectively the deterministic part $u_{0}$ of $u$ and the $U_{j}, j=1, \ldots, L$. For the SC method, we need to solve as many nonlinear problems as collocation points. Even for the nonlinear problem considered here, where the nonlinearity comes from the term $u^{3}$ and which is quite cheap to solve, the perturbation method is about 8 times faster.

To conclude, we can mention that for $h=h_{\text {ref }}$, i.e. without error due to FE approximation and a convergence of the error in $\mathscr{O}\left(\varepsilon^{2}\right)$, the error for the perturbation method is about 1.4 and 3.5 times larger than the error obtained using respectively SC1 and SC2. Again, the error for the perturbation method and the SC method has been accurately computed using the Monte Carlo method. However, for a given problem, that is for fixed value of $\varepsilon$ and $L$, the perturbation method perform better than the SC method in terms of CPU time versus error for $h>h_{r e f}$, especially for nonlinear problems. We plot in Figure 1.11 the computation time with respect to the error for problems (1.11) and (1.78) with $f=f_{2}, \varepsilon=0.5, L=10$ and $2^{-10} \leq h \leq 2^{-3}$. Notice
that the results for $\mathrm{SC1}$ are not depicted on this figure since they are indistinguishable from those of SC2. Finally, we mention that it would be better, in terms of computational costs, to adapt the level $l$ of the sparse grid for the SC-FEM, respectively the order in the approximation $u \approx \sum_{n=0}^{l} \varepsilon^{n} u_{n, h}$ for the perturbation method, with respect to $h$. Indeed, for the value of $h$ for which the total error is not too small, namely of order $\varepsilon$ or larger, it is more suitable to take $l=0$ than $l=1$ since comparable accuracy is reached at lower computational costs. However, the error due to the uncertainty, which is of order $\varepsilon$ for $l=0$, will be dominating at some point (see also Figure 1.5) and the value of $l$ must be increased to be able to further reduce the error.


Figure 1.11: $\log \log$ scale plot of the computational time w.r.t. the error in $L_{p}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$-norm using the SC-FEM with Smolyak and Clenshaw-Curtis abscissas and the perturbation method.

## Conclusions

In this chapter, we have performed error analyses for elliptic PDEs with coefficients affected by small uncertainties, characterized through random variables. The exact random solution has been approximated using a perturbation approach combined with the finite element method for the physical space discretization.

For the first order approximation $u \approx u_{0, h}$, we derived strong and weak a priori error estimates as well as a posteriori error estimates in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ and $L_{P}^{2}\left(\Omega ; L^{2}(D)\right)$ norms. These estimates naturally split into two parts, namely the error in $h$ due to the physical discretization and the error in $\varepsilon$ due to the model. In the a priori error estimation, we have shown that the order of the weak error in the model is twice the order of the strong error, the order of the error due to FE approximation being the same in both cases. The a posteriori error estimator in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm that we have obtained is a computable quantity of order $h+\varepsilon$ if the solution is regular enough in physical space. Given $u_{0, h}$, this estimator is cheap to compute and does not require any other FE solution. It can be used for mesh adaptation so that comparable accuracy in $h$ and $\varepsilon$ is reached. We have shown that taking the $L^{2}$ norm in physical space leads to a gain of one order in $h$ but no improvement in the error due to the model. Finally, we gave a sketch of the derivation of a goal-oriented error estimate, which is more suitable than an estimate in global norm when a particular quantity of interest is
considered.
The a posteriori error estimation procedure for the error in the $L^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm has been applied to the second-order approximation $u \approx u_{0, h}+\varepsilon u_{1, h}$, before giving a generalization for approximations of any order. This reliable error upper bound can be used to adaptively determine the order of approximation and partitions of $D$ such that the total error is below a given tolerance.

A posteriori error estimates have then been derived for a class of nonlinear problems through three different examples. A comparison in terms of computational costs with the stochastic collocation method has been performed, considering an error of order 2 in the model. The perturbation method presents only mild advantages for solving linear problems, the computational cost being halved with respect to the SC method. The situation is different for nonlinear problems. Indeed, the SC method requires the resolution of as many nonlinear problems as collocation points while for the perturbation method, only one nonlinear problem has to be solved for $u_{0, h}$, the remaining problems being linear.

## 1.A Derivation of problems (1.20), (1.21) and (1.22)

We make here some remarks about the derivation of the problems (1.20), (1.21) and (1.22) that we need to solve to build the approximate solution $u_{0}+\varepsilon u_{1}+\varepsilon^{2} u_{2}$. In particular, we will see that the deterministic problems for the terms $u_{0}$ and $u_{1}$ are uniquely determined while those for $u_{2}$ are not. We thus discuss the various ways to build the term $u_{2}$. Moreover, we will make a more precise link between each term and the derivatives of $u=u(\mathbf{x}, \mathbf{y})$ with respect to the $y_{j}$, $j=1, \ldots, L$.

Let us first give some details about the derivation of the problems. Recall that we assume that the diffusion coefficient $a$ has the form

$$
a(\mathbf{x}, \omega)=a(\mathbf{x}, \mathbf{Y}(\omega))=a_{0}(\mathbf{x})+\varepsilon \sum_{j=1}^{L} a_{j}(\mathbf{x}) Y_{j}(\omega) .
$$

Moreover, the random solution $u$ is expanded as

$$
u(\mathbf{x}, \mathbf{Y}(\omega))=u_{0}(\mathbf{x})+\varepsilon u_{1}(\mathbf{x}, \mathbf{Y}(\omega))+\varepsilon^{2} u_{2}(\mathbf{x}, \mathbf{Y}(\omega))+\ldots
$$

with $u_{1}=\sum_{j=1}^{L} U_{j} Y_{j}$ and $u_{2}=\sum_{j, k=1}^{L} U_{j k} Y_{j} Y_{k}$. Similar expansion can be used for the higher order terms, see (1.71) where the general case is treated or [126, 127]. If we substitute the expansions of $a$ and $u$ in the first equation of problem (1.13), we get

$$
-\nabla \cdot\left(\left(a_{0}+\varepsilon \sum_{j=1}^{L} a_{j} Y_{j}\right) \nabla\left(u_{0}+\varepsilon \sum_{j=1}^{L} U_{j} Y_{j}+\varepsilon^{2} \sum_{j, k=1}^{L} U_{j k} Y_{j} Y_{k}+\cdots\right)\right)=f
$$

After recalling that $f$ is deterministic by assumption, we separate then the terms of different
order in $\varepsilon$. The equation for the $\mathscr{O}(1)$ term is

$$
-\nabla \cdot\left(a_{0} \nabla u_{0}\right)=f
$$

which yields problem (1.20) after adding suitable boundary conditions. Next, the equation for the $\mathscr{O}(\varepsilon)$ term is

$$
\begin{equation*}
-\varepsilon \sum_{j=1}^{L} Y_{j} \nabla \cdot\left(a_{0} \nabla U_{j}+a_{j} \nabla u_{0}\right)=0 \tag{1.94}
\end{equation*}
$$

Since the set $\left\{Y_{j}: j=1, \ldots, L\right\}$ is orthonormal, it is in particular linearly independent. Therefore, equation (1.94) holds if and only if each term is zero, i.e.

$$
\begin{equation*}
\nabla \cdot\left(a_{0} \nabla U_{j}+a_{j} \nabla u_{0}\right)=0 \quad \forall j=1, \ldots, L, \tag{1.95}
\end{equation*}
$$

which is nothing else than the first equation of problem (1.21). Notice that the relation (1.95) can also be obtained by multiplying (1.94) by $Y_{k}$ and taking the ensemble mean, see [127], thanks again to the fact that $\mathbb{E}\left[Y_{j} Y_{k}\right]=\delta_{j k}$. Finally, we collect the terms in $\mathscr{O}\left(\varepsilon^{2}\right)$ to obtain

$$
\begin{equation*}
-\varepsilon^{2} \sum_{j, k=1}^{L} Y_{j} Y_{k} \nabla \cdot\left(a_{0} \nabla U_{j k}+a_{j} \nabla U_{k}\right)=0 \tag{1.96}
\end{equation*}
$$

A sufficient condition for (1.96) to hold is that

$$
\begin{equation*}
\nabla \cdot\left(a_{0} \nabla U_{j k}+a_{j} \nabla U_{k}\right)=0 \quad \forall j, k=1, \ldots, L, \tag{1.97}
\end{equation*}
$$

which corresponds to the set of PDEs in (1.22). However, it is not necessary to have (1.97) to verify (1.96) since the set $\left\{Y_{j} Y_{k}: j, k=1, \ldots, L\right\}$ is not linearly independent. Using the fact that $Y_{j} Y_{k}=Y_{k} Y_{j}$, we can rewrite (1.96) as

$$
\begin{equation*}
-\varepsilon^{2} \sum_{1 \leq j \leq k \leq L} Y_{j} Y_{k} \nabla \cdot\left(a_{0} \nabla\left(U_{j k}+U_{k j}\right)+a_{j} \nabla U_{k}+a_{k} \nabla U_{j}\right) \beta_{j k}=0 \tag{1.98}
\end{equation*}
$$

where $\beta_{j k}=1-\frac{1}{2} \delta_{j k}$ is introduced to allow to keep the cases $j<k$ and $j=k$ under the same summation sign. Now, the set $\left\{Y_{j} Y_{k}: 1 \leq j \leq k \leq L\right\}$ is linearly independent [127] and thus (1.98) holds if and only if

$$
\nabla \cdot\left(a_{0} \nabla\left(U_{j k}+U_{k j}\right)+a_{j} \nabla U_{k}+a_{k} \nabla U_{j}\right)=0 \quad \forall 1 \leq j \leq k \leq L .
$$

If we write $\tilde{U}_{j k}:=\frac{U_{j k}+U_{k j}}{2}$ for $j, k=1, \ldots, L$ we have then

$$
u_{2}=\sum_{j, k=1}^{L} U_{j k} Y_{j} Y_{k}=\sum_{1 \leq j \leq k \leq L} \beta_{j k}\left(U_{j k}+U_{k j}\right) Y_{j} Y_{k}=\sum_{j, k=1}^{L} \tilde{U}_{j k} Y_{j} Y_{k} .
$$

Notice that $\tilde{U}_{j k}$ solves

$$
-\nabla \cdot\left(a_{0} \nabla \tilde{U}_{j k}+\frac{a_{j} \nabla U_{k}+a_{k} \nabla U_{j}}{2}\right) \quad \forall j, k=1, \ldots, L
$$

and $\tilde{U}_{j k}+\tilde{U}_{k j}=U_{j k}+U_{k j}$. The advantage of building $u_{2}$ with the $\tilde{U}_{j k}$ instead of the $U_{j k}$ relies in the fact that $\tilde{U}_{j k}=\tilde{U}_{k j}$ while $U_{j k}$ is not necessarily equal to $U_{k j}$. Therefore, the construction of $u_{2}$ with the $\tilde{U}_{j k}$ requires the resolution of $\frac{L(L+1)}{2}$ whereas $L^{2}$ problems need to be solved when the $U_{j k}$ are used.

Notice that the problems we obtain for $u_{0}, U_{j}, U_{j k}$ and $U_{j_{1} j_{2} \ldots j_{n}}$, given by (1.20), (1.21), (1.22) and (1.71), respectively, are equivalent to those derived in [6]. In that paper, the authors apply what they called the method of successive approximations which uses the Karhunen-Loève expansion to represent the stochastic diffusion coefficient combined with the Neumann series expansion method. In fact, applied to the specific linear elliptic diffusion model problem (1.11), the (generalized or standard) Neumann expansion method and the perturbation method are equivalent [121].

In the remaining part of this section, we clarify the link between the various terms $u_{0}, U_{j}, U_{j k}$ and $\tilde{U}_{j k}$ defined above and the derivatives of $u=u(\mathbf{x}, \mathbf{y})$ with respect to the $y_{j}$. In other words, we compare the expansion (1.19) of $u$ with its Taylor expansion around $\mathbf{y}_{0}=\mathbb{E}[\mathbf{Y}]=\mathbf{0}$. Recall that it has been proved (see for instance [7]) that the weak solution $u=u(\mathbf{x}, \mathbf{y})$ of problem (1.13), i.e. the solution of (1.14), is analytic with respect to each variable $y_{j}, j=1, \ldots, L$. First of all, we have

$$
a\left(\mathbf{x}, \mathbf{y}_{0}\right)=a_{0}(\mathbf{x}), \quad \frac{\partial a}{\partial y_{j}}\left(\mathbf{x}, \mathbf{y}_{0}\right)=\varepsilon a_{j}(\mathbf{x}) \quad \text { and } \quad \frac{\partial^{2} a}{\partial y_{k} \partial y_{j}}\left(\mathbf{x}, \mathbf{y}_{0}\right)=0 \quad \forall j, k=1, \ldots, L
$$

Then, we recall that for each $\mathbf{y} \in \Gamma$ the solution $u(\cdot, \mathbf{y}) \in H_{0}^{1}(D)$ of problem (1.14) satisfies

$$
\begin{equation*}
\int_{D} a(\mathbf{x}, \mathbf{y}) \nabla u(\mathbf{x}, \mathbf{y}) \cdot \nabla v(\mathbf{x}) d \mathbf{x}=\int_{D} f(\mathbf{x}) v(\mathbf{x}) d \mathbf{x} \quad \forall v \in H_{0}^{1}(D), \rho-\text { a.e. in } \Gamma . \tag{1.99}
\end{equation*}
$$

The evaluation of equation (1.99) at $\mathbf{y}_{0}$ yields

$$
\begin{equation*}
\int_{D} a_{0}(\mathbf{x}) \nabla u\left(\mathbf{x}, \mathbf{y}_{0}\right) \cdot \nabla v(\mathbf{x}) d \mathbf{x}=\int_{D} f(\mathbf{x}) v(\mathbf{x}) d \mathbf{x} \tag{1.100}
\end{equation*}
$$

We can formally differentiate equation (1.99) with respect to $y_{j}$ to get

$$
\begin{equation*}
\int_{D}\left(\frac{\partial a}{\partial y_{j}} \nabla u+a \nabla \frac{\partial u}{\partial y_{j}}\right)(\mathbf{x}, \mathbf{y}) \cdot \nabla v(\mathbf{x}) d \mathbf{x}=0, \quad j=1, \ldots, L \tag{1.101}
\end{equation*}
$$

and thus for $\mathbf{y}=\mathbf{y}_{0}$ we have

$$
\begin{equation*}
\int_{D}\left(\varepsilon a_{j}(\mathbf{x}) \nabla u\left(\mathbf{x}, \mathbf{y}_{0}\right)+a_{0}(\mathbf{x}) \nabla \frac{\partial u}{\partial y_{j}}\left(\mathbf{x}, \mathbf{y}_{0}\right)\right) \cdot \nabla v(\mathbf{x}) d \mathbf{x}=0, \quad j=1, \ldots, L \tag{1.102}
\end{equation*}
$$

Taking then the derivative of (1.102) with respect to $y_{k}$, or equivalently the second derivative of (1.99), we obtain for $j, k=1, \ldots, L$ the relation

$$
\int_{D}\left(\frac{\partial^{2} a}{\partial y_{k} \partial y_{j}} \nabla u+\frac{\partial a}{\partial y_{j}} \nabla \frac{\partial u}{\partial y_{k}}+\frac{\partial a}{\partial y_{k}} \nabla \frac{\partial u}{\partial y_{j}}+a \nabla \frac{\partial^{2} u}{\partial y_{k} \partial y_{j}}\right)(\mathbf{x}, \mathbf{y}) \cdot \nabla v(\mathbf{x}) d \mathbf{x}=0 .
$$

Since $\frac{\partial^{2} a}{\partial y_{k} \partial y_{j}}=0$, the evaluation of last relation at $\mathbf{y}_{0}$ gives us
$\int_{D}\left(\varepsilon a_{j}(\mathbf{x}) \nabla \frac{\partial u}{\partial y_{k}}\left(\mathbf{x}, \mathbf{y}_{0}\right)+\varepsilon a_{k}(\mathbf{x}) \nabla \frac{\partial u}{\partial y_{j}}\left(\mathbf{x}, \mathbf{y}_{0}\right)+a_{0}(\mathbf{x}) \nabla \frac{\partial^{2} u}{\partial y_{k} \partial y_{j}}\left(\mathbf{x}, \mathbf{y}_{0}\right)\right) \cdot \nabla v(\mathbf{x}) d \mathbf{x}=0, \quad j, k=1, \ldots, L$.

Finally, based on equations (1.100), (1.102) and (1.103) we conclude that $u_{0}=u\left(\cdot, \mathbf{y}_{0}\right), \quad \varepsilon U_{j}=\frac{\partial u}{\partial y_{j}}\left(\cdot, \mathbf{y}_{0}\right), \quad \varepsilon^{2}\left(U_{j k}+U_{k j}\right)=\frac{\partial^{2} u}{\partial y_{k} \partial y_{j}}\left(\cdot, \mathbf{y}_{0}\right) \quad$ and $\quad \varepsilon^{2} \tilde{U}_{j k}=\frac{1}{2} \frac{\partial^{2} u}{\partial y_{k} \partial y_{j}}\left(\cdot, \mathbf{y}_{0}\right)$ for $j, k=1, \ldots, L$.

## 1.B Upper and lower bounds for the error $u-u_{0, h}$ in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm

The goal here is to prove that the error estimator introduced in (1.50) provides both lower and upper bounds for the error $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$. We assume here that $D \subset \mathbb{R}^{d}$ with $d=2$, mentioning that the case $d=1$ can be treated easily since no jump terms occur while the extension to the case $d=3$ is straightforward. We first introduce the estimator in more details, starting from the relation

$$
\begin{aligned}
\mathscr{A}\left(u-u_{0, h}, v ; \mathbf{y}\right) & =\int_{D} f v-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla v \\
& =\mathscr{R}\left(v ; \mathbf{y}_{0}\right)+\left[\mathscr{R}(v ; \mathbf{y})-\mathscr{R}\left(v ; \mathbf{y}_{0}\right)\right]
\end{aligned}
$$

for all $v \in H_{0}^{1}(D)$ and $\rho$-a.e. in $\Gamma$, where $\mathbf{y}_{0}=\mathbb{E}[\mathbf{Y}]=\mathbf{0}$ and

$$
\mathscr{R}(v ; \mathbf{y}):=F(v)-\mathscr{A}\left(u_{0, h}, v ; \mathbf{y}\right)=\int_{D} f v-\int_{D} a(\cdot, \mathbf{y}) \nabla u_{0, h} \cdot \nabla v .
$$

For any $\mathbf{y} \in \Gamma$, let $r(\cdot ; \mathbf{y}): H_{0}^{1}(D) \rightarrow \mathbb{R}$ be defined by

$$
\begin{equation*}
r(\nu ; \mathbf{y}):=\mathscr{R}(\nu ; \mathbf{y})-\mathscr{R}\left(\nu ; \mathbf{y}_{0}\right)=-\int_{D}\left(a(\cdot, \mathbf{y})-a_{0}\right) \nabla u_{0, h} \cdot \nabla v \tag{1.104}
\end{equation*}
$$

The dual norm of $r$ is then given by $\|r(\cdot ; \mathbf{y})\|_{H^{-1}(D)}=\|\nabla w(\cdot, \mathbf{y})\|_{L^{2}(D)}$ with $w(\cdot, \mathbf{y})$ the solution of

$$
\begin{equation*}
\int_{D} \nabla w(\cdot, \mathbf{y}) \cdot \nabla v=r(v ; \mathbf{y}) \quad \forall v \in H_{0}^{1}(D), \rho \text {-a.e. in } \Gamma \text {. } \tag{1.105}
\end{equation*}
$$

We write then $w(\mathbf{x}, \mathbf{Y}(\omega))=\varepsilon \sum_{j=1}^{L} W_{j}(\mathbf{x}) Y_{j}(\omega)$ with $W_{j} \in H_{0}^{1}(D)$ such that

$$
\int_{D} \nabla W_{j} \cdot \nabla v=-\int_{D} a_{j} \nabla u_{0, h} \cdot \nabla v \quad \forall v \in H_{0}^{1}(D) .
$$

Let $w_{h}(\mathbf{x}, \mathbf{Y}(\omega))=\varepsilon \sum_{j=1}^{L} W_{j, h}(\mathbf{x}) Y_{j}(\omega)$, where $W_{j, h} \in V_{h}$ is the FE approximation of $W_{j}$, and let $R$ and $J$ denote the interior element residual and the jump defined on an element $K$ and an internal edge $e$ by respectively

$$
\left.R\right|_{K}=\left.\left(f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right)\right|_{K} \quad \text { and } \quad J_{e}=\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}} .
$$

The spatial and stochastic a posteriori error estimators $\eta_{1}$ and $\hat{\eta}_{2}$ are given by (1.43) and (1.50), respectively, definitions that we recall here for clarity

$$
\begin{align*}
\eta_{1}^{2} & :=\sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2} \quad \text { with } \quad \eta_{K}^{2}=h_{K}^{2}\|R\|_{L^{2}(K)}^{2}+\frac{1}{2} \sum_{e \subset \partial K} h_{e}\|J\|_{L^{2}(e)}^{2},  \tag{1.106}\\
\hat{\eta}_{2}^{2} & :=\varepsilon^{2} \sum_{j=1}^{L}\left\|\nabla W_{j, h}\right\|_{L^{2}(D)}^{2} \tag{1.107}
\end{align*}
$$

To prove the spatial lower bound, see (1.112), we will need some definitions and notation that we introduce now.

For any element $K \in \mathscr{T}_{h}$, using the notation given in Figure 1.12-left, we define the so-called element bubble function $\psi_{K}$ and edge bubble function $\psi_{e_{i}}$, see for instance [118], by

$$
\psi_{K}=27 \lambda_{1} \lambda_{2} \lambda_{3} \quad \text { and } \quad \psi_{e_{i}}=4 \lambda_{i+1} \lambda_{i+2}, \quad i=1,2,3
$$

where the indices are taken modulo 3 and $\lambda_{1}, \lambda_{2}, \lambda_{3}$ are the (linear) barycentric coordinates on $K$. Using the notation used in [118], we denote by $w_{K}$ the union of all the elements sharing an edge with $K$ and, for an internal edge $e$, we write $w_{e}$ the union of the two elements sharing $e$ as an edge, see Figure 1.12 for an illustration.


Figure 1.12: Notation for an element $K$ in $\mathscr{T}_{h}$ (left) and illustration of the domains $w_{K}$ (middle) and $w_{e}$ (right).

The bubble functions satisfy the following properties: for any polynomial $\varphi$ of degree less or equal to $k$ we have

$$
\begin{equation*}
\|\varphi\|_{L^{2}(K)} \leq c_{1}\left\|\psi_{K}^{\frac{1}{2}} \varphi\right\|_{L^{2}(K)}, \quad\left\|\nabla\left(\psi_{K} \varphi\right)\right\|_{L^{2}(K)} \leq c_{2} h_{K}^{-1}\|\varphi\|_{L^{2}(K)} \tag{1.108}
\end{equation*}
$$

and

$$
\begin{equation*}
\|\varphi\|_{L^{2}(e)} \leq c_{3}\left\|\psi_{e}^{\frac{1}{2}} \varphi\right\|_{L^{2}(e)}, \quad\left\|\nabla\left(\psi_{e} \varphi\right)\right\|_{L^{2}\left(w_{e}\right)} \leq c_{4} h_{e}^{-\frac{1}{2}}\|\varphi\|_{L^{2}(e)}, \quad\left\|\psi_{e} \varphi\right\|_{L^{2}\left(w_{e}\right)} \leq c_{5} h_{e}^{\frac{1}{2}}\|\varphi\|_{L^{2}(e)}, \tag{1.109}
\end{equation*}
$$

where the constants $C_{i}, i=1, \ldots, 5$, depend only on $k$ and on the shape regularity parameter of $\mathscr{T}_{h}$ given in (1.23). Moreover, we have

$$
0 \leq \psi_{K}(\mathbf{x}) \leq 1 \quad \forall \mathbf{x} \in K, \quad \psi_{K}(\mathbf{x})=0 \forall \mathbf{x} \notin K, \quad \max _{\mathbf{x} \in K} \psi_{K}(\mathbf{x})=1
$$

and

$$
0 \leq \psi_{e}(\mathbf{x}) \leq 1 \quad \forall \mathbf{x} \in w_{e}, \quad \psi_{e}(\mathbf{x})=0 \quad \forall \mathbf{x} \notin w_{e}, \quad \max _{\mathbf{x} \in w_{e}} \psi_{e}(\mathbf{x})=1
$$

For any element $K$, we denote by $\bar{g}_{K}$ the mean value of $g$ on $K$ and similarly we denote by $\bar{g}_{e}$ the mean value of $g$ on any internal edge $e$, i.e.

$$
\bar{g}_{K}=\frac{1}{|K|} \int_{K} g \quad \text { and } \quad \bar{g}_{e}=\frac{1}{|e|} \int_{e} g .
$$

Finally, we introduce the oscillation term $\theta_{K}$ defined by

$$
\begin{equation*}
\theta_{K}^{2}:=\sum_{T \subset w_{K}} h_{T}^{2}\left\|R-\bar{R}_{T}\right\|_{L^{2}(T)}^{2}+\sum_{e \subset \partial K} h_{e}\left\|J-\bar{J}_{e}\right\|_{L^{2}(e)}^{2} . \tag{1.110}
\end{equation*}
$$

We can now state the upper and lower bounds, given in the following proposition.
Proposition 1.B.1. Let $u$ be the weak solution of problem (1.11) and let $u_{0, h}$ be the solution of problem (1.30), respectively. There exist two constants $C_{1}, C_{2}>0$ depending only on the mesh aspect ratio and $s \in(0,1]$ such that

$$
\begin{gather*}
\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)} \leq \frac{1}{a_{\text {min }}}\left(C_{1} \eta_{1}+\hat{\eta}_{2}\right)+\mathscr{O}\left(\varepsilon h^{s}\right),  \tag{1.111}\\
\eta_{1} \leq C_{2}\left[a_{\text {max }}\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}+\hat{\eta}_{2}+\left(\sum_{K \in \mathscr{T}_{h}} \theta_{K}^{2}\right)^{\frac{1}{2}}\right]+\mathscr{O}\left(\varepsilon h^{s}\right) \tag{1.112}
\end{gather*}
$$

and

$$
\begin{equation*}
\hat{\eta}_{2} \leq a_{\max }\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}+C_{1} \eta_{1}+\mathscr{O}\left(\varepsilon h^{s}\right) . \tag{1.113}
\end{equation*}
$$

Proof. We first derive a bound for the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm of $w$ (resp. $w_{h}$ ) in term of the norm of $w_{h}$ (resp. $w$ ) and higher order terms, where $w$ is the solution (1.105) and $w_{h}$ its FE approximation. Let us introduce $\psi(\mathbf{x}, \mathbf{Y}(\omega))=\varepsilon \sum_{j=1}^{L} \psi_{j}(\mathbf{x}) Y_{j}(\omega)$, where $\psi_{j} \in H_{0}^{1}(D)$ is the solution of

$$
\int_{D} \nabla \psi_{j} \cdot \nabla v=-\int_{D} a_{j} \nabla u_{0} \cdot \nabla v \quad \forall v \in H_{0}^{1}(D),
$$

and let $\psi_{h}$ denotes its FE approximation. Notice that $\psi(\cdot, \mathbf{Y}(\omega))$ solves

$$
\int_{D} \nabla \psi \cdot \nabla v=-\int_{D}\left(a-a_{0}\right) \nabla u_{0} \cdot \nabla v \quad \forall v \in H_{0}^{1}(D), \text { a.s. in } \Omega,
$$

which is similar to the problem (1.105) for $w$, except that $u_{0, h}$ is replaced by $u_{0}$ in the righthand side. Thanks to the triangle inequality, we obtain

$$
\|\nabla w\|_{L^{2}(D)} \leq\left\|\nabla w_{h}\right\|_{L^{2}(D)}+\|\nabla(w-\psi)\|_{L^{2}(D)}+\left\|\nabla\left(\psi-\psi_{h}\right)\right\|_{L^{2}(D)}+\left\|\nabla\left(\psi_{h}-w_{h}\right)\right\|_{L^{2}(D)}
$$

from which we can deduce

$$
\|\nabla w\|_{L_{P}^{2}\left(\Omega ; L^{2}(D)\right)} \leq\left\|\nabla w_{h}\right\|_{L_{P}^{2}\left(\Omega ; L^{2}(D)\right)}+C \varepsilon h^{s}
$$

where $s \in(0,1]$ depends only on the regularity of $u_{0}, \psi_{j}, j=1, \ldots, L$, and the domain $D$ and $C$ is a (deterministic) positive constant independent of $h$ and $\varepsilon$ but dependent on the mesh aspect ratio, $\left|u_{0}\right|_{H^{1+s}(D)}$ and $\left|\psi_{j}\right|_{H^{1+s}(D)}, j=1, \ldots, L$. Therefore, recalling that $w_{h}=\varepsilon \sum_{j=1}^{L} W_{j, h} Y_{j}$ and using $\mathbb{E}\left[Y_{i} Y_{j}\right]=\delta_{i j}$ we get

$$
\begin{equation*}
\|\nabla w\|_{L_{P}^{2}\left(\Omega ; L^{2}(D)\right)} \leq \hat{\eta}_{2}+C \varepsilon h^{s} \tag{1.114}
\end{equation*}
$$

with $\hat{\eta}_{2}$ given in (1.107). Finally, proceeding in the same way we can obtain the relation

$$
\begin{equation*}
\hat{\eta}_{2}=\left\|\nabla w_{h}\right\|_{L_{P}^{2}\left(\Omega ; L^{2}(D)\right)} \leq\|\nabla w\|_{L_{P}^{2}\left(\Omega ; L^{2}(D)\right)}+C \varepsilon h^{s} \tag{1.115}
\end{equation*}
$$

We now prove the three bounds (1.111), (1.112) and (1.113) separately. The proof of (1.112) is inspired by what is done in $[99,118]$, while the idea for the proof of $(1.113)$ is based on the proof of efficiency of the error estimator proposed in [102] for the Reduced Basis method. In the sequel, all the equations hold a.s. in $\Omega$ without specifically mentioning it.

Upper bound The proof is similar to the one of Proposition 1.3.5, only the bound of term controlling the stochastic error is different. For any $v \in H_{0}^{1}(D)$, taking $v_{h}=I_{h}$ the Clément interpolant of $v$ we have

$$
\begin{aligned}
\int_{D} a \nabla\left(u-u_{0, h}\right) \cdot \nabla v & =\int_{D} f v-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla v \\
& =\int_{D} f\left(v-v_{h}\right)-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla\left(v-v_{h}\right)+\int_{D} \nabla w \cdot \nabla v \\
& \leq\left[C_{1}\left(\sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2}\right)^{\frac{1}{2}}+\|\nabla w\|_{L^{2}(D)}\right]\|\nabla v\|_{L^{2}(D)}
\end{aligned}
$$

where $C_{1}$ depends only on the constants in (1.26) and (1.28). Since $a_{\min }$ is a lower bound for $a$, taking $v=u-u_{0, h}$ we get

$$
a_{m i n}\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)} \leq C_{1} \eta_{1}+\|\nabla w\|_{L^{2}(D)}
$$

and thus, taking the $L_{P}^{2}(\Omega)$ norm on both sides of the last inequality we have

$$
a_{m i n}\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)} \leq C_{1} \eta_{1}+\|\nabla w\|_{L_{P}^{2}\left(\Omega ; L^{2}(D)\right)} .
$$

Finally, we obtain (1.111) using (1.114).
$h$-lower bound First of all, notice that for any $v \in H_{0}^{1}(D)$ we have

$$
\begin{align*}
\int_{D} a \nabla\left(u-u_{0, h}\right) \cdot \nabla v & =\sum_{K \in \mathscr{T}_{h}} \int_{K} R v+\sum_{e \in \mathscr{T}_{h}} \int_{e} J v-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla v \\
& =\sum_{K \in \mathscr{T}_{h}} \int_{K} R v+\sum_{e \in \mathscr{T}_{h}} \int_{e} J v+\int_{D} \nabla w \cdot \nabla v \tag{1.116}
\end{align*}
$$

The proof is then divided into three steps.

1. Let $K$ be any element in $\mathscr{T}_{h}$ and let $v_{K}=\bar{R}_{K} \psi_{K}$. We take $v=v_{K}$ in (1.116). Since $\operatorname{supp} \psi_{K} \subset K$, we have

$$
\int_{K} a \nabla\left(u-u_{0, h}\right) \cdot \nabla v_{K}=\int_{K} \bar{R}_{K} v_{K}+\int_{K}\left(R-\bar{R}_{K}\right) v_{K}+\int_{K} \nabla w \cdot \nabla v_{K}
$$

and thus, using the properties of the element bubble function given in (1.108), we obtain

$$
h_{K}\left\|\bar{R}_{K}\right\|_{L^{2}(K)} \leq c_{1}^{2} c_{2} a_{\max }\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(K)}+c_{1}^{2} c_{2}\|\nabla w\|_{L^{2}(K)}+c_{1}^{2} h_{K}\left\|R-\bar{R}_{K}\right\|_{L^{2}(K)} .
$$

Thanks to triangle's inequality, we finally obtain

$$
\begin{equation*}
h_{K}\|R\|_{L^{2}(K)} \leq c_{1}^{2} c_{2} a_{\max }\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(K)}+c_{1}^{2} c_{2}\|\nabla w\|_{L^{2}(K)}+\left(1+c_{1}^{2}\right) h_{K}\left\|R-\bar{R}_{K}\right\|_{L^{2}(K)} \tag{1.117}
\end{equation*}
$$

2. Let $e$ be any interior edge of $\mathscr{T}_{h}$, let $v_{e}=\bar{J}_{e} \psi_{e}$ and let $K_{1}$ and $K_{2}$ be the two elements that share $e$ as an edge. We take $v=v_{e}$ in (1.116) to get

$$
\int_{w_{e}} a \nabla\left(u-u_{0, h}\right) \cdot \nabla v_{e}=\sum_{K \in w_{e}} \int_{K} R v_{e}+\int_{e} \bar{J}_{e} v_{e}+\int_{e}\left(J-\bar{J}_{e}\right) v_{e}+\int_{w_{e}} \nabla w \cdot \nabla v_{e} .
$$

Therefore, using the properties of the edge bubble function given in (1.109), we obtain

$$
\begin{aligned}
h_{e}^{\frac{1}{2}}\left\|\bar{J}_{e}\right\|_{L^{2}(e)} \leq & c_{3}^{2} c_{4} a_{\max }\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}\left(w_{e}\right)}+c_{3}^{2} c_{5} h_{e}\|R\|_{L^{2}\left(w_{e}\right)} \\
& +c_{3}^{2} h_{e}^{\frac{1}{2}}\left\|J-\bar{J}_{e}\right\|_{L^{2}(e)}+c_{3}^{2} c_{4}\|\nabla w\|_{L^{2}\left(w_{e}\right)}
\end{aligned}
$$

and thus

$$
\begin{aligned}
h_{e}^{\frac{1}{2}}\|J\|_{L^{2}(e)} \leq & c_{3}^{2} c_{4} a_{m a x}\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}\left(w_{e}\right)}+c_{3}^{2} c_{5} h_{e}\|R\|_{L^{2}\left(w_{e}\right)}+\left(1+c_{3}^{2}\right) h_{e}^{\frac{1}{2}}\left\|J-\bar{J}_{e}\right\|_{L^{2}(e)} \\
& +c_{3}^{2} c_{4}\|\nabla w\|_{L^{2}\left(w_{e}\right)} \\
\leq & \sum_{i=1}^{2}\left[a_{m a x} c_{3}^{2}\left(c_{4}+c_{1}^{2} c_{2} c_{5}\right)\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}\left(K_{i}\right)}+\left(1+c_{1}^{2}\right) c_{3}^{2} c_{5} h_{K_{i}}\left\|R-\bar{R}_{K_{i}}\right\|_{L^{2}\left(K_{i}\right)}\right. \\
& \left.c_{3}^{2}\left(c_{4}+c_{1}^{2} c_{2} c_{5}\right)\|\nabla w\|_{L^{2}\left(K_{i}\right)}\right]+\left(1+c_{3}^{2}\right) h_{e}^{\frac{1}{2}}\left\|J-\bar{J}_{e}\right\|_{L^{2}(e)}
\end{aligned}
$$

using relation (1.117).
3. Putting everything together, we obtain for any element $K \in \mathscr{T}_{h}$

$$
\begin{aligned}
\eta_{K}^{2}= & h_{K}^{2}\|R\|_{L^{2}(K)}^{2}+\frac{1}{2} \sum_{e \subset \partial K} h_{e}\|J\|_{L^{2}(e)}^{2} \\
\leq & C_{2}\left(a_{\max }^{2}\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}\left(w_{K}\right)}^{2}+\|\nabla w\|_{L^{2}\left(w_{K}\right)}^{2}\right. \\
& \left.+\sum_{T \subset w_{K}} h_{T}^{2}\left\|R-\bar{R}_{T}\right\|_{L^{2}(T)}^{2}+\sum_{e \subset \partial K} h_{e}\left\|J-\bar{J}_{e}\right\|_{L^{2}(e)}^{2}\right)
\end{aligned}
$$

where $C_{2}$ depends only on the regularity of the mesh (through the constants $c_{i}, i=$ $1, \ldots, 5)$. Recalling the definition of $\theta_{K}$ in (1.110), if we sum over all $K \in \mathscr{T}_{h}$ and use the relation $\left(a^{2}+b^{2}+c^{2}\right) \leq(a+b+c)^{2}$ valid for any non-negative numbers $a, b, c$, we get

$$
\eta_{1} \leq C_{2}\left[a_{m a x}\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}+\|\nabla w\|_{L^{2}(D)}+\left(\sum_{K \in \mathscr{T}_{h}} \theta_{K}^{2}\right)^{\frac{1}{2}}\right]
$$

where $C_{2}$ has changed but still only depends on the mesh aspect ratio. Finally, we obtain (1.112) taking the $L_{P}^{2}(\Omega)$ norm and using (1.114).
$\varepsilon$-lower bound For any $v \in H_{0}^{1}(D)$ we have

$$
\begin{equation*}
\int_{D} \nabla w \cdot \nabla v=-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla v=\int_{D} a \nabla\left(u-u_{0, h}\right) \cdot \nabla v-\int_{D} a_{0} \nabla\left(u_{0}-u_{0, h}\right) \cdot \nabla v \tag{1.118}
\end{equation*}
$$

Taking $v=w$ in (1.118) and noticing that the last term of (1.118) is nothing else than (minus) the residual for $u_{0, h}$, we can easily derive the bound

$$
\|\nabla w\|_{L^{2}(D)}^{2} \leq\left[a_{\max }\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}+C_{1}\left(\sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2}\right)^{\frac{1}{2}}\right]\|\nabla w\|_{L^{2}(D)}
$$

where $C_{1}$ depends only on the constants in (1.26) and (1.28). From the last relation, we deduce taking the $L_{P}^{2}(\Omega)$ that

$$
\|\nabla w\|_{L_{P}^{2}\left(\Omega ; L^{2}(D)\right)} \leq a_{\max }\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}+C_{1} \eta_{1}
$$

which conclude the proof thanks to (1.115).
Remark 1.B.2. Since $u_{0, h}$ is piecewise affine, if $a_{0}$ is piecewise constant then we have $R=f$ and $J=\bar{J}_{e}$. Therefore, in this case $\theta_{K}$ reduces to $\sum_{T \subset w_{K}} h_{T}^{2}\left\|f-\bar{f}_{T}\right\|_{L^{2}(T)}^{2}$ which does no longer depend on $u_{0, h}$. It is often refereed to as data oscillation.

Remark 1.B.3. We deduce from the three relations (1.111), (1.112) and (1.113) that

$$
a_{\min } \leq \frac{\hat{\eta}_{2}}{\left\|u-u_{0, h}\right\|} \leq a_{\max } \quad \text { as } h \rightarrow 0
$$

and

$$
C_{1}^{-1} a_{\min } \leq \frac{\eta_{1}}{\left\|u-u_{0, h}\right\|} \leq C_{2} a_{\max } \quad \text { as } \varepsilon \rightarrow 0
$$

where $\|\cdot\|$ denotes the $L^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm and $C_{1}$ and $C_{2}$ are two positive constants depending only on the mesh aspect ratio.

## 1.C Estimation of the interpolation constant

In this section, we briefly present the value of the interpolation constant $C_{H_{0}^{1}}$ that can be included in the error estimator to get a sharp spatial error estimator. This value depends on the degree of the finite element space as well as if we are in 1D, 2D or 3D.

In the one-dimensional case, we have already mentioned that the constant for $\mathbb{P}_{1}$ finite element can be set to $C_{H_{0}^{1}}=\frac{1}{3.46} \approx \frac{1}{2 \sqrt{3}}$. The latter corresponds to the theoretical value $\left(\frac{1}{p+1}\right)^{1 / p} \frac{1}{2}$ with $p=2$ given in [9].

For the 2D case, we consider the (deterministic) Poisson problem $-\Delta u_{0}=f$ with homogeneous Dirichlet boundary conditions. We set $D=(0,1)^{2}$ and $u_{0}\left(x_{1}, x_{2}\right)=\sin \left(2 \pi x_{1}\right) \sin \left(4 \pi x_{2}\right)$ and compute the corresponding right-hand side given by

$$
\begin{equation*}
f\left(x_{1}, x_{2}\right)=20 \pi^{2} \sin \left(2 \pi x_{2}\right) \sin \left(4 \pi x_{2}\right) . \tag{1.119}
\end{equation*}
$$

We give in Table 1.15 the error $\left\|\nabla\left(u_{0}-u_{0, h}\right)\right\|_{L^{2}(D)}$ and the two estimators $\eta_{1}$ and $\hat{\eta}_{1}$ defined by

$$
\eta_{1}^{2}=\sum_{K \in \mathscr{T}_{h}} h_{K}^{2}\left\|f+\Delta u_{h}\right\|_{L^{2}(K)}^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e}\left\|\left[\nabla u_{h} \cdot \mathbf{n}_{\mathbf{e}}\right]_{\mathbf{n}_{\mathbf{e}}}\right\|_{L^{2}(e)}^{2}
$$

and

$$
\hat{\eta}_{1}^{2}=\sum_{K \in \mathscr{T}}^{h} \text { }\left[h_{K}^{2}\left\|f+\Delta u_{h}\right\|_{L^{2}(K)}^{2}+\frac{1}{4} \sum_{e \subset \partial K} h_{e}\left\|\left[\nabla u_{h} \cdot \mathbf{n}_{\mathbf{e}}\right]_{\mathbf{n}_{\mathbf{e}}}\right\|_{L^{2}(e)}^{2}\right]
$$

We consider both structured and Delaunay triangulations with $N=256$ equidistant vertices on each boundary of $D$, see Figure 1.13 where the meshes for the case $N=16$ are given.

The constant $1 / C_{H_{1}^{0}}$ can then be set to $\eta_{1} /\left\|\nabla\left(u_{0}-u_{0, h}\right)\right\|_{L^{2}(D)}$ or $\hat{\eta}_{1} /\left\|\nabla\left(u_{0}-u_{0, h}\right)\right\|_{L^{2}(D)}$ depending on the definition of the estimator.



Figure 1.13: Structured (left) and Delaunay (right) triangulations of $D$ with $N=16$.

|  | Structured mesh |  |  |  |  |  | Delaunay mesh |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | error | $\eta_{1}$ | e.i. | $\hat{\eta}_{1}$ | e.i. | error | $\eta_{1}$ | e.i. | $\hat{\eta}_{1}$ | e.i. |  |
| $\mathbb{P}_{1}$ | $1.279 \mathrm{e}-1$ | $7.352 \mathrm{e}-1$ | 5.75 | $6.472 \mathrm{e}-1$ | 5.06 | $1.037 \mathrm{e}-1$ | $5.934 \mathrm{e}-1$ | 5.72 | $5.296 \mathrm{e}-1$ | 5.11 |  |
| $\mathbb{P}_{1 b}$ | $1.204 \mathrm{e}-1$ | $5.225 \mathrm{e}-1$ | 4.34 | $3.939 \mathrm{e}-1$ | 3.27 | $9.450 \mathrm{e}-2$ | $3.590 \mathrm{e}-1$ | 3.80 | $2.668 \mathrm{e}-1$ | 2.82 |  |
| $\mathbb{P}_{2}$ | $9.592 \mathrm{e}-4$ | $8.464 \mathrm{e}-3$ | 8.82 | $8.195 \mathrm{e}-3$ | 8.54 | $6.905 \mathrm{e}-4$ | $6.473 \mathrm{e}-3$ | 9.37 | $6.385 \mathrm{e}-3$ | 9.25 |  |
| $\mathbb{P}_{3}$ | $3.130 \mathrm{e}-6$ | $7.136 \mathrm{e}-5$ | 22.80 | $6.924 \mathrm{e}-5$ | 22.12 | $2.017 \mathrm{e}-3$ | $4.865 \mathrm{e}-5$ | 24.12 | $4.749 \mathrm{e}-5$ | 23.55 |  |

Table 1.15: Error, estimator and effectivity index for the Poisson problem.

Notice that we get similar values when considering other cases than (1.119). We see from the results of Table 1.15 that, as expected, the interpolation constant depends on the polynomial degree of the finite elements. Moreover, we could go further by estimating separately the efficiency of the interior residual and the contribution of the jump terms, but we will not do it in this thesis.

## 2 Elliptic model problems with other sources of uncertainty

## Introduction

We extend here the results of Chapter 1 to include other sources of uncertainty. We first consider the case of random Neumann boundary conditions. The analysis is very similar to the one presented in Chapter 1. It is even easier in this case since the solution $u$ depends linearly on the random input, and thus only the first two terms in the expansion are nonzero. We consider then the case where two random input data are affected by uncertainty, namely we consider a random diffusion coefficient combined with a random forcing term. Two different sets of random variables are used to describe each uncertain input data. Finally, numerical results are given to illustrate the theoretical findings.

### 2.1 Neumann random boundary conditions

We consider the problem:
find $u: D \times \Omega \rightarrow \mathbb{R}$ such that a.s. in $\Omega$ :

$$
\left\{\begin{array}{rlrl}
-\operatorname{div}\left(a_{0}(\mathbf{x}) \nabla u(\mathbf{x}, \omega)\right) & =f(\mathbf{x}) & & \mathbf{x} \in D  \tag{2.1}\\
u(\mathbf{x}, \omega) & =0 & & \mathbf{x} \in \Gamma_{D} \\
a_{0}(\mathbf{x}) \frac{\partial u(\mathbf{x}, \omega)}{\partial \mathbf{n}} & =g(\mathbf{x}, \omega) & \mathbf{x} \in \Gamma_{N},
\end{array}\right.
$$

where $\Gamma_{D} \cup \Gamma_{N}=\partial D$ with $\Gamma_{D} \cap \Gamma_{N}=\varnothing$ and $\Gamma_{D} \neq \varnothing$. We assume that $a_{0}$ is bounded from below by $a_{0, \text { min }}$ and that $g$ is characterized by $L$ independent random variables $\left\{Y_{j}\right\}_{j=1}^{L}$ with zero mean and unit variance as

$$
\begin{equation*}
g(\mathbf{x}, \omega)=g\left(\mathbf{x}, Y_{1}(\omega), \ldots, Y_{L}(\omega)\right)=g_{0}(\mathbf{x})+\varepsilon \sum_{j=1}^{L} g_{j}(\mathbf{x}) Y_{j}(\omega) \tag{2.2}
\end{equation*}
$$

with $g_{j} \in L^{2}\left(\Gamma_{N}\right), j=0,1, \ldots, L$. Using the same notation as in the previous chapter, we can rewrite problem (2.1) in parametric form as:
find $u: D \times \Gamma \rightarrow \mathbb{R}$ such that $\rho$-a.e. in $\Gamma$ we have:

$$
\left\{\begin{align*}
-\operatorname{div}\left(a_{0}(\mathbf{x}) \nabla u(\mathbf{x}, \mathbf{y})\right) & =f(\mathbf{x}) & & \mathbf{x} \in D  \tag{2.3}\\
u(\mathbf{x}, \mathbf{y}) & =0 & & \mathbf{x} \in \Gamma_{D} \\
a_{0}(\mathbf{x}) \frac{\partial u(\mathbf{x}, \mathbf{y})}{\partial \mathbf{n}} & =g(\mathbf{x}, \mathbf{y}) & & \mathbf{x} \in \Gamma_{N}
\end{align*}\right.
$$

whose weak formulation reads:
find $u(\cdot, \mathbf{y}) \in W$ such that

$$
\begin{equation*}
\int_{D} a_{0} \nabla u(\cdot, \mathbf{y}) \cdot \nabla v=\int_{D} f v+\int_{\Gamma_{N}} g(\cdot, \mathbf{y}) v \quad \forall v \in W, \rho-\text { a.e. in } \Gamma \tag{2.4}
\end{equation*}
$$

with $W:=H_{\Gamma_{D}}^{1}(D)=\left\{v \in H^{1}(D): v=0\right.$ on $\left.\Gamma_{D}\right\}$ that we endow with the gradient norm $\|\cdot\|_{W}:=$ $\|\nabla \cdot\|_{L^{2}(D)}$. This can be done thanks to the Friedrich-Poincaré inequality

$$
\begin{equation*}
\|v\|_{L^{2}(D)} \leq C_{F}\|\nabla v\|_{L^{2}(D)} \quad \forall v \in W \tag{2.5}
\end{equation*}
$$

which holds as long as $\Gamma_{D} \neq \varnothing$. Using again a perturbation technique, we write

$$
u(\mathbf{x}, \mathbf{Y}(\omega))=u_{0}(\mathbf{x})+\varepsilon u_{1}(\mathbf{x}, \mathbf{Y}(\omega))+\varepsilon^{2} u_{2}(\mathbf{x}, \mathbf{Y}(\omega))+\ldots
$$

where $u_{0}: D \rightarrow \mathbb{R}$ is the solution of

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{0}(\mathbf{x}) \nabla u_{0}(\mathbf{x})\right) & =f(\mathbf{x}) & & \mathbf{x} \in D  \tag{2.6}\\
u_{0}(\mathbf{x}) & =0 & & \mathbf{x} \in \Gamma_{D} \\
a_{0}(\mathbf{x}) \frac{\partial u_{0}(\mathbf{x})}{\partial \mathbf{n}} & =g_{0}(\mathbf{x}) & \mathbf{x} \in \Gamma_{N}
\end{array}\right.
$$

and $u_{1}=\sum_{j=1}^{L} U_{j} Y_{j}$ with $U_{j}: D \rightarrow \mathbb{R}, j=1, \ldots, L$, the solution of

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{0}(\mathbf{x}) \nabla U_{j}(\mathbf{x})\right) & =0 & & \mathbf{x} \in D  \tag{2.7}\\
U_{j}(\mathbf{x}) & =0 & & \mathbf{x} \in \Gamma_{D} \\
a_{0}(\mathbf{x}) \frac{\partial U_{j}(\mathbf{x})}{\partial \mathbf{n}} & =g_{j}(\mathbf{x}) & & \mathbf{x} \in \Gamma_{N} .
\end{array}\right.
$$

Contrary to the problem with random diffusion coefficient $a$ of the previous chapter, we will show that we have here $u=u_{0}+\varepsilon u_{1}$, i.e. there is no term of order higher than one in $\varepsilon$. This is due to the linear dependence of $u$ with respect to the uncertain input data $g$. The same holds for instance when the forcing term $f$ is random, see also the next section. The weak formulation of problems (2.6) and (2.7) is given by, respectively,

$$
\begin{equation*}
\text { find } u_{0} \in W: \quad \int_{D} a_{0} \nabla u_{0} \cdot \nabla v=\int_{D} f v+\int_{\Gamma_{N}} g_{0} v \quad \forall v \in W \tag{2.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\text { find } U_{j} \in W: \quad \int_{D} a_{0} \nabla U_{j} \cdot \nabla v=\int_{\Gamma_{N}} g_{j} v \quad \forall v \in W . \tag{2.9}
\end{equation*}
$$

Notice that the problems for $u_{0}$ and the $U_{j}, j=1, \ldots, L$, are decoupled, that is the solution $u_{0}$ does not appear in the problem for $U_{j}$ as it is the case when dealing with random diffusion coefficient, see problem (1.21). We first show the following three properties.

Proposition 2.1.1. Let $u$ be the weak solution of problem (2.1) and let $u_{0}$ and $U_{j}, j=1, \ldots, L$, be the solutions of problems (2.8) and (2.9), respectively. Then for $u_{1}=\sum_{j=1}^{L} U_{j} Y_{j}$ we have

1. $\mathbb{E}[u]=u_{0}$
2. $u=u_{0}+\varepsilon u_{1}$
3. $\operatorname{Var}[u]=\varepsilon^{2} \sum_{j=1}^{L} U_{j}^{2}$.

Proof. First of all, if we take the expected value on both sides of equation (2.4) with $\mathbf{y}=\mathbf{Y}(\omega)$, we get

$$
\int_{D} a_{0} \nabla \mathbb{E}[u] \cdot \nabla v=\int_{D} f v+\int_{\Gamma_{N}} g_{0} v \quad \forall v \in W
$$

and thus, subtracting equation (2.8) we obtain

$$
\int_{D} a_{0} \nabla\left(\mathbb{E}[u]-u_{0}\right) \cdot \nabla v=0 \quad \forall v \in W
$$

If we take then $v=\mathbb{E}[u]-u_{0}$, we have

$$
0 \leq a_{0, \text { min }}\left\|\nabla\left(\mathbb{E}[u]-u_{0}\right)\right\|_{L^{2}(D)}^{2} \leq\left\|a_{0}^{\frac{1}{2}} \nabla\left(\mathbb{E}[u]-u_{0}\right)\right\|_{L^{2}(D)}^{2}=0
$$

which implies $\mathbb{E}[u]=u_{0}$ a.e. in $D$. We proceed similarly for the second relation. Indeed, without writing the dependence of each function, we have for any $v \in W$ and a.s. in $\Omega$

$$
\begin{aligned}
\int_{D} a_{0} \nabla\left(u-\left(u_{0}+\varepsilon u_{1}\right)\right) \cdot \nabla v & =\int_{D} a_{0} \nabla u \cdot \nabla v-\int_{D} a_{0} \nabla u_{0} \cdot \nabla v-\varepsilon \int_{D} a_{0} \nabla u_{1} \cdot \nabla v \\
& =\int_{\Gamma_{N}} g v-\int_{\Gamma_{N}} g_{0} v-\varepsilon \sum_{j=1}^{L} \int_{\Gamma_{N}} g_{j} v \\
& =0
\end{aligned}
$$

Taking then $v=u-u_{0}-\varepsilon u_{1} \in W$ a.s. in $\Omega$, we can easily show that $\left\|u-\left(u_{0}+\varepsilon u_{1}\right)\right\|_{L_{P}^{2}(\Omega ; W)}=0$ and thus $u=u_{0}+\varepsilon u_{1}$ a.e. in $D$ and a.s. in $\Omega$. Finally, we directly get

$$
\operatorname{Var}[u]=\mathbb{E}\left[(u-\mathbb{E}[u])^{2}\right]=\mathbb{E}\left[\varepsilon^{2} u_{1}^{2}\right]=\varepsilon^{2} \sum_{j=1}^{L} U_{j}^{2}
$$

using the fact that $\mathbb{E}\left[Y_{i} Y_{j}\right]=\delta_{i j}$.
Remark 2.1.2. Notice that we could also see that u does not contain any term of order $\mathscr{O}\left(\varepsilon^{k}\right)$ for any $k \geq 2$ by observing that the term $u_{k}$ in the expansion of $u$ would be the solution of the
problem

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{0}(\mathbf{x}) \nabla u_{k}(\mathbf{x}, \omega)\right) & = & 0 & \mathbf{x} \in D \\
u_{k}(\mathbf{x}, \omega) & = & 0 & \mathbf{x} \in \Gamma_{D} \\
a_{0}(\mathbf{x}) \frac{\partial u_{k}(\mathbf{x}, \omega)}{\partial \mathbf{n}} & = & 0 & \mathbf{x} \in \Gamma_{N}
\end{array}\right.
$$

for which $u_{k}=0$ is the obvious solution.

To simplify the notation in the a posteriori error estimates given below, we introduce the generalized jumps across an edge $e$ defined as

$$
J_{e, 0}\left(u_{0, h}\right):=\left\{\begin{array}{rll}
\frac{1}{2}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{e} & \text { if } & e \subset D \\
g_{0}-\lim _{t \rightarrow 0^{+}}\left(a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right)\left(\mathbf{x}-t \mathbf{n}_{e}\right) & \text { if } & e \subset \Gamma_{N} \\
0 & \text { if } & e \subset \Gamma_{D}
\end{array}\right.
$$

with $[\cdot]_{\mathbf{n}_{e}}$ the jump across an interior edge $e$ defined by

$$
[\varphi]_{\mathbf{n}_{e}}(\mathbf{x}):=\lim _{t \rightarrow 0^{+}}\left(\varphi\left(\mathbf{x}+t \mathbf{n}_{e}\right)-\varphi\left(\mathbf{x}-t \mathbf{n}_{e}\right)\right)
$$

For $j=1, \ldots, L$, the quantity $J_{e, j}\left(U_{j, h}\right)$ is defined analogously replacing $u_{0, h}$ and $g_{0}$ by $U_{j, h}$ and $g_{j}$, respectively. Moreover, we will need the following trace inequality (see for instance [109])

$$
\begin{equation*}
\|v\|_{L^{2}\left(\Gamma_{N}\right)} \leq C_{T}\|v\|_{H^{1}(D)} \quad \forall v \in H^{1}(D) \tag{2.10}
\end{equation*}
$$

## Error estimation for $u-u_{0, h}$

We consider the $\mathbb{P}_{1}$ finite element approximation of problem (2.8) given by

$$
\begin{equation*}
\text { find } u_{0, h} \in W_{h}: \quad \int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v_{h}=\int_{D} f v_{h}+\int_{\Gamma_{N}} g_{0} v_{h} \quad \forall v_{h} \in W_{h} \tag{2.11}
\end{equation*}
$$

with $W_{h}=\left\{\nu \in C^{0}(\bar{D}): \nu_{\mid K} \in \mathbb{P}_{1} \forall K \in \mathscr{T}_{h}\right\} \cap W$ and $\mathscr{T}_{h}$ a regular triangulation of $D$. We have the following a posteriori error estimate for the error $u-u_{0, h}$, yielding an error of order $\mathscr{O}\left(h^{s}+\varepsilon\right)$ with $s \in(0,1]$ depending on the regularity of the solution.

Proposition 2.1.3. Let $u$ be the weak solution of problem (2.1) and let $u_{0, h}$ be the solution of problem (2.11). Then, there exists a constant $C>0$ depending only on $C_{F}$ in (2.5), $C_{T}$ in (2.10) and the mesh aspect ratio such that

$$
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{C}{a_{0, \min }}\left(\eta_{h}^{2}+\eta_{\varepsilon}^{2}\right)^{\frac{1}{2}},
$$

with

$$
\begin{aligned}
& \eta_{h}^{2}:=\sum_{K \in \mathscr{T}_{h}} h_{K}^{2}\left\|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e}\left\|J_{e, 0}\left(u_{0, h}\right)\right\|_{L^{2}(e)}^{2} \\
& \eta_{\varepsilon}^{2}:=\varepsilon^{2} \sum_{j=1}^{L}\left\|g_{j}\right\|_{L^{2}\left(\Gamma_{N}\right)}^{2} .
\end{aligned}
$$

Proof. For any $v \in W$ and a.s. in $\Omega$ we have

$$
\int_{D} a_{0} \nabla\left(u-u_{0, h}\right) \cdot \nabla v=\underbrace{\int_{D} f v+\int_{\Gamma_{N}} g_{0} v-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v}_{=: \mathrm{I}}+\underbrace{\int_{\Gamma_{N}}\left(g-g_{0}\right) v}_{=: \mathrm{II}} .
$$

We bound each term separately. The term I, which is the residual for $u_{0, h}$, can be bounded as follows

$$
\mathrm{I} \leq C_{1}\left[\sum_{K \in \mathscr{T}_{h}} h_{K}^{2}\left\|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \in \mathscr{T}_{h}}\left\|J_{e, 0}\left(u_{0, h}\right)\right\|_{L^{2}(e)}^{2}\right]^{\frac{1}{2}}\|\nabla v\|_{L^{2}(D)}
$$

where $C_{1}$ depends only on the interpolation constants in (1.26) and (1.28). The second term is bounded by

$$
\mathrm{II}=\int_{\Gamma_{N}}\left(g-g_{0}\right) v \leq\left\|g-g_{0}\right\|_{L^{2}\left(\Gamma_{N}\right)}\|\nu\|_{L^{2}\left(\Gamma_{N}\right)} \leq C_{2}\left\|g-g_{0}\right\|_{L^{2}\left(\Gamma_{N}\right)}\|\nabla v\|_{L^{2}(D)}, \quad C_{2}=C_{T} \sqrt{1+C_{F}^{2}}
$$

Combining these two bounds with the fact that $a_{0}$ is larger than $a_{0, \min }$ we get

$$
\begin{aligned}
\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2} \leq & \frac{1}{a_{0, \text { min }}}\left[C_{1}\left(\sum_{K \in \mathscr{T}_{h}} h_{K}^{2}\left\|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \in \mathscr{T}_{h}}\left\|J_{e, 0}\left(u_{0, h}\right)\right\|_{L^{2}(e)}^{2}\right)^{\frac{1}{2}}\right. \\
& \left.+C_{2}\left\|g-g_{0}\right\|_{L^{2}\left(\Gamma_{N}\right)}\right]
\end{aligned}
$$

Taking the expected value of the square of last inequality and using the fact that $\mathbb{E}\left[Y_{i} Y_{j}\right]=\delta_{i j}$ allows us to conclude the proof.

## Error estimation for $u-\left(u_{0, h}+\varepsilon u_{1, h}\right)$

Let $U_{j, h}$ be the $\mathbb{P}_{1}$ finite element approximation of $U_{j}$ which solves

$$
\begin{equation*}
\text { find } U_{j, h} \in W_{h}: \quad \int_{D} a_{0} \nabla U_{j, h} \cdot \nabla v_{h}=\int_{\Gamma_{N}} g_{j} v_{h} \quad \forall v_{h} \in W_{h} . \tag{2.12}
\end{equation*}
$$

We have the following a posteriori error estimate for the error $u-\left(u_{0, h}+\varepsilon u_{1, h}\right)$, yielding an error of order $\mathscr{O}\left(h^{s}+\varepsilon h^{s}\right), s \in(0,1]$. In particular, there is no term of order $\mathscr{O}\left(\varepsilon^{k}\right), k \geq 2$, and thus no pure statistical error.

Proposition 2.1.4. Let $u$ be the weak solution of problem (2.1) and let $u_{0, h}$ be the solution of
problem (2.11). Moreover, let $u_{1, h}=\sum_{j=1}^{L} U_{j, h} Y_{j}$ with $U_{j, h}$ the solution of problem (2.12). Then, there exists a constant $C>0$ depending only on $C_{F}$ in (2.5), $C_{T}$ in (2.10) and the mesh aspect ratio such that

$$
\mathbb{E}\left[\left\|\nabla\left(u-\left(u_{0, h}+\varepsilon u_{1, h}\right)\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{C}{a_{0, \min }}\left(\eta_{h}^{2}+\eta_{\varepsilon h}^{2}\right)^{\frac{1}{2}},
$$

with

$$
\begin{aligned}
\eta_{h}^{2} & :=\sum_{K \in \mathscr{T}_{h}} h_{K}^{2}\left\|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e}\left\|J_{e, 0}\left(u_{0, h}\right)\right\|_{L^{2}(e)}^{2} \\
\eta_{\varepsilon h}^{2} & :=\varepsilon^{2} \sum_{j=1}^{L}\left[\sum_{K \in \mathscr{T}_{h}} h_{K}^{2}\left\|\nabla \cdot\left(a_{0} \nabla U_{j, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \in \mathscr{T}_{h}} h_{e}\left\|J_{e, j}\left(U_{j, h}\right)\right\|_{L^{2}(e)}^{2}\right] .
\end{aligned}
$$

Proof. The proof can easily be deduced from the relation

$$
\int_{D} a_{0} \nabla\left(u-\left(u_{0, h}+\varepsilon u_{1, h}\right)\right) \cdot \nabla v=\underbrace{\int_{D} f v+\int_{\Gamma_{N}} g_{0} v-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v}_{=: \mathrm{I}}+\underbrace{\varepsilon \sum_{j=1}^{L}\left(\int_{\Gamma_{N}} g_{j} v-\int_{D} a_{0} \nabla U_{j, h} \cdot \nabla v\right)}_{=: \mathrm{II}}
$$

a.s. in $\Omega$, where I and II are nothing else than the residual for $u_{0, h}$ and $u_{1, h}$, respectively. Each of these terms can be bounded in a standard way to conclude.

### 2.2 Two sources of uncertainty

We consider again the diffusion model problem but with two input data affected by uncertainty, namely the diffusion coefficient and the source term:
find $u: D \times \Omega \rightarrow \mathbb{R}$ such that a.s. in $\Omega$ it holds:

$$
\left\{\begin{array}{rll}
-\operatorname{div}(a(\mathbf{x}, \omega) \nabla u(\mathbf{x}, \omega)) & =f(\mathbf{x}, \omega) & \mathbf{x} \in D  \tag{2.13}\\
u(\mathbf{x}, \omega) & =0 & \\
\mathbf{x} \in \partial D
\end{array}\right.
$$

where $f(\cdot, \omega) \in L^{2}(D)$ a.s. in $\Omega$ and $a$ is uniformly bounded from below and above by $a_{\text {min }}$ and $a_{\max }$, respectively. We prescribe homogeneous Dirichlet boundary conditions for simplicity but we could easily extend the following results to other kinds of boundary conditions, including random boundary conditions as treated in the previous section. We assume that the two random inputs $a$ and $f$ are characterized through a finite number of random variables

$$
a(\mathbf{x}, \omega)=a\left(\mathbf{x}, Y_{1}(\omega), \ldots, Y_{L}(\omega)\right) \quad \text { and } \quad f(\mathbf{x}, \omega)=f\left(\mathbf{x}, Z_{1}(\omega), \ldots, Z_{M}(\omega)\right) .
$$

More precisely, we assume an affine dependence of $a$ and $f$ with respect to the random
variables as follows

$$
\begin{align*}
& a(\mathbf{x}, \omega)=a_{0}(\mathbf{x})+\varepsilon \sum_{j=1}^{L} a_{j}(\mathbf{x}) Y_{j}(\omega),  \tag{2.14}\\
& f(\mathbf{x}, \omega)=f_{0}(\mathbf{x})+\delta \sum_{j=1}^{M} f_{j}(\mathbf{x}) Z_{j}(\omega), \tag{2.15}
\end{align*}
$$

where $\left\{Y_{j}\right\}_{j=1}^{L}$ and $\left\{Z_{j}\right\}_{j=1}^{M}$ are two families of independent random variables with zero mean and $\operatorname{Var}\left(Y_{j}\right)=\left(\sigma_{j}^{y}\right)^{2}<\infty$ and $\operatorname{Var}\left(Z_{i}\right)=\left(\sigma_{i}^{z}\right)^{2}<\infty$ for $j=1, \ldots, L$ and $i=1, \ldots, M$. Moreover, we assume that $f_{j} \in L^{2}(D)$ for $j=0,1, \ldots, M$. The two parameters $\varepsilon$ and $\delta$ control the amount of randomness in $a$ and $f$, respectively.

Remark 2.2.1. The case where only the forcing term is affected by uncertainty can be easily deduced from the one considered here by setting $\varepsilon=0$.

Let $\mathbf{Y}=\left(Y_{1}, \ldots, Y_{L}\right), \mathbf{Z}=\left(Z_{1}, \ldots, Z_{M}\right)$ and $\mathbf{R}=(\mathbf{Y}, \mathbf{Z})$. For $j=1, \ldots, L$, let $\Gamma_{j}^{y}$ denote the bounded image in $\mathbb{R}$ of $Y_{j}$ and for $i=1, \ldots, M$ let $\Gamma_{i}^{z}$ be the image in $\mathbb{R}$ of $Z_{i}$. Moreover, we write $\rho_{j}^{y}$ and $\rho_{i}^{z}$ their probability density function. Let $\Gamma=\Gamma^{y} \times \Gamma^{z}=\Gamma_{1}^{y} \times \ldots \Gamma_{L}^{y} \times \Gamma_{1}^{z} \times \ldots \times \Gamma_{M}^{z}$. Thanks to the independence of the random variables, the joint density function $\rho: \Gamma \rightarrow \mathbb{R}^{+}$of the random vector $\mathbf{R}$ is given by $\rho(\mathbf{r})=\rho^{y}(\mathbf{y}) \rho^{z}(\mathbf{z})=\Pi_{j=1}^{L} \rho_{j}^{y}\left(y_{j}\right) \Pi_{i=1}^{M} \rho_{i}^{z}\left(z_{i}\right)$ for all $\mathbf{r}=(\mathbf{y}, \mathbf{z}) \in \Gamma$ with $\mathbf{y}=\left(y_{1}, \ldots, y_{L}\right) \in \Gamma^{y}$ and $\mathbf{z}=\left(z_{1}, \ldots, z_{M}\right) \in \Gamma^{z}$. By definition, for any measurable function $g: \Gamma \rightarrow \mathbb{R}$, the expected value of the random variable $g(\mathbf{R})$ is $\mathbb{E}[g(\mathbf{R})]=\int_{\Gamma} g(\mathbf{r}) \rho(\mathbf{r}) d \mathbf{r}$. The finite dimensional noise assumption implies that the random solution $u$ of problem (2.13) can be described by $L+M$ random variables

$$
u(\mathbf{x}, \omega)=u\left(\mathbf{x}, Y_{1}(\omega), \ldots, Y_{L}(\omega), Z_{1}(\omega), \ldots, Z_{M}(\omega)\right)
$$

Therefore, the solution $u$ can be sought in the probability space $(\Omega, \mathscr{F}, P)$ or equivalently in $(\Gamma, B(\Gamma), \rho(\mathbf{r}) d \mathbf{r})$. The problem (2.13) can indeed be equivalently written in the following deterministic parametric form:
find $u: D \times \Gamma^{y} \times \Gamma^{z} \rightarrow \mathbb{R}$ such that $\rho$-a.e. in $\Gamma$ it holds:

$$
\left\{\begin{align*}
-\operatorname{div}(a(\mathbf{x}, \mathbf{y}) \nabla u(\mathbf{x}, \mathbf{y}, \mathbf{z})) & =f(\mathbf{x}, \mathbf{z}) & & \mathbf{x} \in D  \tag{2.16}\\
u(\mathbf{x}, \mathbf{y}, \mathbf{z}) & =0 & & \mathbf{x} \in \partial D .
\end{align*}\right.
$$

The pointwise weak formulation of (2.16) reads:
find $u \in L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)$ such that

$$
\begin{equation*}
\mathscr{A}(u(\cdot, \mathbf{y}, \mathbf{z}), v ; \mathbf{y})=F(v ; \mathbf{z}) \quad \forall v \in H_{0}^{1}(D), \rho \text {-a.e. in } \Gamma, \tag{2.17}
\end{equation*}
$$

where

$$
\begin{align*}
\mathscr{A}(u(\cdot, \mathbf{y}, \mathbf{z}), v ; \mathbf{y}) & =\int_{D} a(\mathbf{x}, \mathbf{y}) \nabla u(\mathbf{x}, \mathbf{y}, \mathbf{z}) \cdot \nabla v(\mathbf{x}) d \mathbf{x}  \tag{2.18}\\
F(v ; \mathbf{z}) & =\int_{D} f(\mathbf{x}, \mathbf{z}) v(\mathbf{x}) d \mathbf{x} . \tag{2.19}
\end{align*}
$$

The well-posedness of problem (2.17) can be shown using Lax-Milgram's lemma. In particular, the assumptions on $f_{0}, f_{i}$ and $Z_{i}, i=1, \ldots, M$, ensure that $f \in L_{\rho}^{2}\left(\Gamma ; L^{2}(D)\right)$.

We assume small uncertainty and use a perturbation approach expanding $u$ with respect to $\varepsilon$ and $\delta$ as

$$
\begin{align*}
u(\mathbf{x}, \mathbf{Y}(\omega), \mathbf{Z}(\omega))= & u_{0}(\mathbf{x})+\varepsilon u_{1}^{y}(\mathbf{x}, \mathbf{Y}(\omega))+\delta u_{1}^{z}(\mathbf{x}, \mathbf{Z}(\omega)) \\
& +\varepsilon^{2} u_{2}^{y}(\mathbf{x}, \mathbf{Y}(\omega))+\varepsilon \delta u_{2}^{y z}(\mathbf{x}, \mathbf{Y}(\omega), \mathbf{Z}(\omega))+\delta^{2} u_{2}^{z}(\mathbf{x}, \mathbf{Z}(\omega))+\ldots \tag{2.20}
\end{align*}
$$

Notice that similarly to Section 2.1, there will be no term of higher order than 1 in $\delta$, i.e. $u_{2}^{z}$ vanishes, due to the linear dependence of $u$ with respect to $f$.

The problem for $u_{0}$ is given by:
find $u_{0}: D \rightarrow \mathbb{R}$ such that

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{0}(\mathbf{x}) \nabla u_{0}(\mathbf{x})\right) & =f_{0}(\mathbf{x}) & & \mathbf{x} \in D  \tag{2.21}\\
u_{0}(\mathbf{x}) & =0 & & \mathbf{x} \in \partial D .
\end{array}\right.
$$

Writing then $u_{1}^{y}(\mathbf{x}, \mathbf{Y}(\omega))=\sum_{j=1}^{L} U_{j}^{y}(\mathbf{x}) Y_{j}(\omega)$ and $u_{1}^{z}(\mathbf{x}, \mathbf{Z}(\omega))=\sum_{j=1}^{M} U_{j}^{z}(\mathbf{x}) Z_{j}(\omega)$, the first order term in (2.20) is obtained by solving the following $L+M$ deterministic uncoupled problems:
find $U_{j}^{y}: D \rightarrow \mathbb{R}$ such that

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{0}(\mathbf{x}) \nabla U_{j}^{y}(\mathbf{x})+a_{j}(\mathbf{x}) \nabla u_{0}(\mathbf{x})\right) & = & 0 & \mathbf{x} \in D  \tag{2.22}\\
U_{j}^{y}(\mathbf{x}) & = & 0 & \mathbf{x} \in \partial D
\end{array} \quad j=1, \ldots, L\right.
$$

and
find $U_{j}^{z}: D \rightarrow \mathbb{R}$ such that

$$
\left\{\begin{array}{rlll}
-\operatorname{div}\left(a_{0}(\mathbf{x}) \nabla U_{j}^{z}(\mathbf{x})\right) & =f_{j}(\mathbf{x}) & \mathbf{x} \in D  \tag{2.23}\\
U_{j}^{z}(\mathbf{x}) & =0 & & \mathbf{x} \in \partial D \quad j=1, \ldots, M .
\end{array}\right.
$$

Notice that the solution $u_{0}$ of problem (2.21) is required in problem (2.22) but not in (2.23).

## Error $u-u_{0, h}$

Let $u_{0, h}$ be the $\mathbb{P}_{1}$ finite element approximation of $u_{0}$, i.e. the solution of

$$
\begin{equation*}
\text { find } u_{0, h} \in V_{h}: \quad \int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v_{h}=\int_{D} f_{0} v_{h} \quad \forall v_{h} \in V_{h} \tag{2.24}
\end{equation*}
$$

where $V_{h}=\left\{v \in C^{0}(\bar{D}): v_{\mid K} \in \mathbb{P}_{1} \forall K \in \mathscr{T}_{h}\right\} \cap V$ and $\mathscr{T}_{h}$ is a regular triangulation of $D$. The following proposition gives an a posteriori error estimation of the error $u-u_{0, h}$ in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm.

Proposition 2.2.2. Let $u$ be the weak solution of problem (2.13) and let $u_{0, h}$ be the solution of problem (2.24). There exists a constant $C>0$ depending only the mesh aspect ratio such that

$$
\begin{equation*}
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{\sqrt{3}}{a_{\min }}\left[C \eta_{h}^{2}+\eta_{\varepsilon}^{2}+C_{P}^{2} \eta_{\delta}^{2}\right]^{\frac{1}{2}}, \tag{2.25}
\end{equation*}
$$

where $C_{P}$ is the Poincaré constant and

$$
\begin{align*}
& \eta_{h}^{2}:=\sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2} \text { with } \eta_{K}^{2}=h_{K}^{2}\left\|f_{0}+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \subset \partial K} h_{e}\left\|\frac{1}{2}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2} \\
& \eta_{\varepsilon}^{2}:=\varepsilon^{2} \sum_{j=1}^{L}\left(\sigma_{j}^{y}\right)^{2}\left\|a_{j} \nabla u_{0, h}\right\|_{L^{2}(D)}^{2}  \tag{2.26}\\
& \eta_{\delta}^{2}:=\delta^{2} \sum_{j=1}^{M}\left(\sigma_{j}^{z}\right)^{2}\left\|f_{j}\right\|_{L^{2}(D)}^{2} . \tag{2.28}
\end{align*}
$$

Proof. For any $\nu \in H_{0}^{1}(D)$ and a.s. in $\Omega$ we have

$$
\begin{equation*}
\int_{D} a \nabla\left(u-u_{0, h}\right) \cdot \nabla v=\underbrace{\int_{D} f_{0} v-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v}_{=: \mathrm{I}}+\underbrace{\int_{D}\left(f-f_{0}\right) v-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla v}_{=: \mathrm{II}} \tag{2.29}
\end{equation*}
$$

The term I is nothing else but the residual for $u_{0, h}$ and we have
$\mathrm{I} \leq\left(C \sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2}\right)^{\frac{1}{2}}\|\nabla v\|_{L^{2}(D)}, \quad \eta_{K}^{2}=h_{K}^{2}\left\|f+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \subset \partial K} h_{e}\left\|\frac{1}{2}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2}$
with $C$ an interpolation constant which depends only on the interpolation constants in (1.26) and (1.28). For the second term, thanks to Cauchy-Schwarz and Poincaré inequalities we have the bound

$$
\mathrm{II} \leq\left(C_{P}\left\|f-f_{0}\right\|_{L^{2}(D)}+\left\|\left(a-a_{0}\right) \nabla u_{0, h}\right\|_{L^{2}(D)}\right)\|\nabla v\|_{L^{2}(D)}
$$

where $C_{P}$ denotes the constant in Poincaré's inequality. Using the lower bound on $a$, we thus
obtain

$$
\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)} \leq \frac{1}{a_{\min }}\left[\left(C \sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2}\right)^{\frac{1}{2}}+C_{P}\left\|f-f_{0}\right\|_{L^{2}(D)}+\left\|\left(a-a_{0}\right) \nabla u_{0, h}\right\|_{L^{2}(D)}\right]
$$

The result follows from taking the expected value on the square of the last inequality.

As we will see in the numerical results, the loss due to the use of the Poincaré inequality for the source term is dependent on the input data. In other words, the efficiency of the estimator $\eta_{\delta}$ in (2.28), for which the Poincaré inequality has been used, will be different from one case to another. A way to skirt this drawback is to replace $\eta_{\delta}$ by an implicit estimator obtained by computing (approximately) the dual norm of a residual to be defined. The price to pay is that the computation of this estimator, given in the following proposition, requires the resolution of $M$ additional (Poisson) problems.

Proposition 2.2.3. Let $u$ be the weak solution of problem (2.13) and let $u_{0, h}$ be the solution of problem (2.24). There exists a constant $C>0$ depending only on the mesh aspect ratio such that

$$
\begin{equation*}
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{\sqrt{3}}{a_{\min }}\left[C \eta_{h}^{2}+\eta_{\varepsilon}^{2}+\hat{\eta}_{\delta}^{2}\right]^{\frac{1}{2}}+\text { h.o.t. } \tag{2.31}
\end{equation*}
$$

where $\eta_{h}$ and $\eta_{\varepsilon}$ are as in (2.26) and (2.27), respectively, and

$$
\begin{equation*}
\hat{\eta}_{\delta}^{2}=\delta^{2} \sum_{j=1}^{M}\left(\sigma_{j}^{z}\right)^{2}\left\|\nabla W_{j, h}\right\|_{L^{2}(D)}^{2} \tag{2.32}
\end{equation*}
$$

with $W_{j, h} \in V_{h}$ the solution of

$$
\int_{D} \nabla W_{j, h} \cdot \nabla v_{h}=\int_{D} f_{j} v_{h} \quad \forall v_{h} \in V_{h} .
$$

Proof. The only difference with respect to the proof of Proposition 2.2 .2 is how we bound the term II of (2.29) due to the uncertainty in the input data, more precisely the part due to the forcing term. Let us introduce for any $\mathbf{z} \in \Gamma^{z}$ the operator $R(\cdot ; \mathbf{z}): H_{0}^{1}(D) \rightarrow \mathbb{R}$ defined by

$$
R(v ; \mathbf{z}):=\int_{D}\left(f(\cdot ; \mathbf{z})-f_{0}\right) v=\delta \sum_{j=1}^{M} z_{j} \int_{D} f_{j} v .
$$

The dual norm of $R$ is then given by $\|R(\cdot ; \mathbf{z})\|_{H^{-1}(D)}=\|\nabla w(\cdot ; \mathbf{z})\|_{L^{2}(D)}$ with $w$ the Riesz representant of $R$, i.e. $w(\cdot ; \mathbf{z}) \in H_{0}^{1}(D)$ is such that $\int_{D} \nabla w \cdot \nabla v=R(\nu ; \mathbf{z})$ for all $v \in H_{0}^{1}(D)$ and $\rho^{z}$-a.e. in $\Gamma^{z}$. We can write $w=w(\mathbf{x}, \mathbf{Z}(\omega))=\delta \sum_{j=1}^{M} W_{j}(\mathbf{x}) Z_{j}(\omega)$ with $W_{j} \in H_{0}^{1}(D)$ the solution of

$$
\begin{equation*}
\int_{D} \nabla W_{j} \cdot \nabla v=\int_{D} f_{j} v \quad \forall v \in H_{0}^{1}(D) \tag{2.33}
\end{equation*}
$$

from which we deduce

$$
\mathbb{E}\left[\|R\|_{H^{-1}(D)}^{2}\right]=\delta^{2} \sum_{j=1}^{M}\left(\sigma_{j}^{z}\right)^{2}\left\|\nabla W_{j}\right\|_{L^{2}(D)}^{2}
$$

Since the solution of (2.33) can not be computed exactly, we can replace it by its finite element approximation $W_{j, h} \in V_{h}$. Doing so introduce an error of higher order, the proof being similar to that of Proposition 1.B.1.

We mention that the computational cost to get the error estimator $\hat{\eta}_{\delta}$ is the same as that needed to get the finite element approximation $u_{1, h}^{z}$ of the term $u_{1}^{z}$ in the expansion (2.20). Since the solution $u$ depends linearly on the input $f$, there is no term of order $\delta^{2}$ and it would thus be better to simply add the term $\delta u_{1, h}^{z}$ to $u_{0, h}$. The quantification of the error in $\mathscr{O}(\delta h)$ so introduced is made precisely in Proposition 2.2 .5 , see the term $\eta_{\delta h}$. As mentioned in Chapter 1 , the computational cost might be reduced introducing auxiliary local problems defined on an element or a small subdomain.

Remark 2.2.4. Notice that we could use the same procedure as used in Proposition 2.2.3 for the whole term II, and not only the part due to $f$, by considering the residual defined for all $v \in H_{0}^{1}(D)$ and $(\mathbf{y}, \mathbf{z}) \in \Gamma$ by

$$
R(v ; \mathbf{y}, \mathbf{z})=\int_{D}\left(f(\cdot, \mathbf{z})-f_{0}\right) v-\int_{D}\left(a(\cdot, \mathbf{y})-a_{0}\right) \nabla u_{0, h} \cdot \nabla v
$$

The dual norm of $R$ is then given by $\|R(\cdot ; \mathbf{y}, \mathbf{z})\|_{H^{-1}(D)}=\|\nabla w(\cdot ; \mathbf{y}, \mathbf{z})\|_{L^{2}(D)}$ where $w(\cdot ; \mathbf{y}, \mathbf{z}) \in H_{0}^{1}(D)$ $\rho$-a.e. in $\Gamma$ writes

$$
w(\mathbf{x} ; \mathbf{Y}(\omega), \mathbf{Z}(\omega))=\varepsilon \sum_{j=1}^{L} W_{j}^{y}(\mathbf{x}) Y_{j}(\omega)+\delta \sum_{j=1}^{M} W_{j}^{z}(\mathbf{x}) Z_{j}(\omega)
$$

with $W_{j}^{y}$ and $W_{j}^{z}$ the solutions of

$$
\int_{D} \nabla W_{j}^{y} \cdot \nabla v=-\int_{D} a_{j} \nabla u_{0, h} \cdot \nabla v \quad \forall v \in H_{0}^{1}(D)
$$

and

$$
\int_{D} \nabla W_{j}^{z} \cdot \nabla v=\int_{D} f_{j} v \quad \forall v \in H_{0}^{1}(D)
$$

respectively. Writing $W_{j, h}^{y}$ and $W_{j, h}^{z}$ the finite element approximations of $W_{j}^{y}$ and $W_{j}^{z}$, respectively, the error estimate reads then

$$
\begin{equation*}
\mathbb{E}\left[\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{\sqrt{3}}{a_{m i n}}\left[C \eta_{h}^{2}+\hat{\eta}_{\varepsilon}^{2}+\hat{\eta}_{\delta}^{2}\right]^{\frac{1}{2}}+\text { h.o.t., } \tag{2.34}
\end{equation*}
$$

with $\eta_{h}$ defined in (2.26) and

$$
\begin{equation*}
\hat{\eta}_{\varepsilon}^{2}:=\varepsilon^{2} \sum_{j=1}^{L}\left(\sigma_{j}^{y}\right)^{2}\left\|\nabla W_{j, h}^{y}\right\|_{L^{2}(D)}^{2} \quad \text { and } \quad \hat{\eta}_{\delta}^{2}=\delta^{2} \sum_{j=1}^{M}\left(\sigma_{j}^{z}\right)^{2}\left\|\nabla W_{j, h}^{z}\right\|_{L^{2}(D)}^{2} . \tag{2.35}
\end{equation*}
$$

Error $u-\left(u_{0, h}+\varepsilon u_{1, h}^{y}+\delta u_{1, h}^{z}\right)$
Let us write $u_{h}^{1}=u_{0, h}+\varepsilon u_{1, h}^{y}+\delta u_{1, h}^{z}$, where $u_{1, h}^{y}=\sum_{j=1}^{L} U_{j, h}^{y} Y_{j}, u_{1, h}^{z}=\sum_{i=1}^{M} U_{i, h}^{z} Z_{i}$ and, for $j=1, \ldots, L$ and $i=1, \ldots, M, U_{j, h}^{y}$ and $U_{i, h}^{z}$ are the solutions of respectively

$$
\begin{equation*}
\int_{D}\left(a_{0} \nabla U_{j, h}^{y}+a_{j} \nabla u_{0, h}\right) \cdot \nabla v_{h}=0 \quad \forall v_{h} \in V_{h} \tag{2.36}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{D} a_{0} \nabla U_{i, h}^{z} \cdot \nabla v_{h}=\int_{D} f_{i} v_{h} \quad \forall v_{h} \in V_{h} . \tag{2.37}
\end{equation*}
$$

To simplify the notation, we write $w_{j, h}=a_{0} \nabla U_{j, h}^{y}+a_{j} \nabla u_{0, h}$. The following proposition gives an a posteriori error estimation of the error $u-u_{h}^{1}$ in the $L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)$ norm. Notice that in particular, there is no term of order $\delta^{2}$. Indeed, we deduce from Proposition 2.2.5 that $\left\|u-u_{h}^{1}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}=\mathscr{O}\left(h+h(\varepsilon+\delta)+\varepsilon^{2}+\varepsilon \delta\right)$ if $u$ is regular enough in the physical space.

Proposition 2.2.5. Let $u$ be the weak solution of problem (2.13) and let $u_{0, h}, U_{j, h}^{y}, j=1, \ldots, L$ and $U_{i, h}^{z}, i=1, \ldots, M$, be the solutions of problems (2.24), (2.36) and (2.37), respectively. There exist constants $C_{1}, C_{2}, C_{3}>0$ depending only on the mesh aspect ratio such that

$$
\begin{equation*}
\mathbb{E}\left[\left\|\nabla\left(u-u_{h}^{1}\right)\right\|_{L^{2}(D)}^{2}\right]^{\frac{1}{2}} \leq \frac{2}{a_{\min }}\left[C_{1} \eta_{h}^{2}+C_{2} \eta_{\varepsilon h}^{2}+C_{3} \eta_{\delta h}^{2}+2 \eta_{\varepsilon \delta}^{2}\right]^{\frac{1}{2}}, \tag{2.38}
\end{equation*}
$$

where

$$
\begin{aligned}
\eta_{h}^{2}= & \sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2} \text { with } \eta_{K}^{2}=h_{K}^{2}\left\|f_{0}+\nabla \cdot\left(a_{0} \nabla u_{0, h}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \subset \partial K} h_{e}\left\|\frac{1}{2}\left[a_{0} \nabla u_{0, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2} \\
\eta_{\varepsilon h}^{2}= & \varepsilon^{2} \sum_{K \in \mathscr{T}_{h}} \sum_{j=1}^{L}\left(\sigma_{j}^{y}\right)^{2} \theta_{K, j}^{2} \quad \text { with } \quad \theta_{K, j}^{2}=h_{K}^{2}\left\|\nabla \cdot w_{j, h}\right\|_{L^{2}(K)}^{2}+\sum_{e \subset \partial K} h_{e}\left\|\frac{1}{2}\left[w_{j, h} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2} \\
\eta_{\delta h}^{2}= & \delta^{2} \sum_{K \in \mathscr{T}_{h}} \sum_{j=1}^{M}\left(\sigma_{j}^{z}\right)^{2} v_{K, j}^{2} \quad \text { with } \quad \vartheta_{K, j}^{2}=h_{K}^{2}\left\|f_{j}+\nabla \cdot\left(a_{0} \nabla U_{j, h}^{z}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \subset \partial K} h_{e}\left\|\frac{1}{2}\left[a_{0} \nabla U_{j, h}^{z} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2} \\
\eta_{\varepsilon \delta}^{2}= & \varepsilon^{4}\left(\int_{D} \sum_{i=1}^{L} a_{i}^{2}\left|\nabla U_{i, h}^{y}\right|^{2} \mathbb{E}\left[Y_{i}^{4}\right]+\int_{D} \sum_{\substack{i, j=1 \\
i \neq j}}^{L}\left(\sigma_{i}^{y} \sigma_{j}^{y}\right)^{2}\left[a_{i}^{2}\left|\nabla U_{j, h}^{y}\right|^{2}+2 a_{i} a_{j} \nabla U_{i, h}^{y} \cdot \nabla U_{j, h}^{y}\right]\right) \\
& +(\varepsilon \delta)^{2} \sum_{j=1}^{L} \sum_{i=1}^{M}\left(\sigma_{j}^{y} \sigma_{i}^{z}\right)^{2}\left\|a_{j} \nabla U_{i, h}^{z}\right\|_{L^{2}(D)}^{2} .
\end{aligned}
$$

Proof. The proof can be easily obtained from the relation

$$
\int_{D} a \nabla\left(u-u_{h}^{1}\right) \cdot \nabla v=\mathrm{I}+\mathrm{II}+\mathrm{III}+\mathrm{IV} \quad \forall v \in H_{0}^{1}(D), \text { a.s. in } \Omega
$$

with

$$
\begin{aligned}
\mathrm{I} & =\int_{D} f_{0} v-\int_{D} a_{0} \nabla u_{0, h} \cdot \nabla v \\
\text { II } & =-\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla v-\varepsilon \int_{D} a_{0} \nabla u_{1, h}^{y} \cdot \nabla v \\
\text { III } & =\int_{D}\left(f-f_{0}\right) v-\delta \int_{D} a_{0} \nabla u_{1, h}^{z} \cdot \nabla v \\
\mathrm{IV} & =-\varepsilon \int_{D}\left(a-a_{0}\right) \nabla u_{1, h}^{y} \cdot \nabla v-\delta \int_{D}\left(a-a_{0}\right) \nabla u_{1, h}^{z} \cdot \nabla v
\end{aligned}
$$

bounding then each term separately.

### 2.3 Numerical results

We consider one-dimensional examples with $D=(0,1)$. In the results below, the true error is computed with the standard Monte Carlo method with a sample size of $K=10000$ and a reference solution computed on a uniform partition with mesh size $h_{r e f}=2^{-12}$.

## Random forcing term

We consider first the case where only the forcing term is random, that is we set $\varepsilon=0$ in (2.14). As mentioned above, the efficiency of the stochastic estimator $\eta_{\delta}$ in (2.28) depends on the input data, due to the use of the Poincaré inequality for the forcing term. To observe this behaviour, we consider the following two cases

$$
\begin{equation*}
f(x, \omega)=1+\delta \sum_{j=1}^{M} f_{j}(x) Z_{j}(\omega), \quad f_{j}(x)=\frac{\sin (2 \pi j x)}{j} \tag{2.39}
\end{equation*}
$$

and

$$
\begin{equation*}
f(x, \omega)=1+\delta \sum_{j=1}^{M} f_{j}(x) Z_{j}(\omega), \quad f_{j}(x)=0.5 j^{-\frac{1}{2}} e^{-50 j(x-0.5)^{2}} \tag{2.40}
\end{equation*}
$$

where $Z_{j}, j=1, \ldots, M$, are uniform random variables in $[-\sqrt{3}, \sqrt{3}]$.
The plot of several realizations of the forcing term for the case (2.39) with $M=6$ and $M=50$ is given in Figures 2.1 and 2.2, respectively, where the corresponding solution is also depicted. The forcing term contains much more high oscillating features with $M=50$ than in the case $M=6$. The difference between the two cases for the corresponding solutions is not noticeable, but is indeed present.

Chapter 2. Elliptic model problems with other sources of uncertainty


Figure 2.1: Six realizations of the random forcing term $f$ given in (2.39) with $\delta=0.5$ and $M=6$ (left) and the corresponding solution (right).


Figure 2.2: Six realizations of the random forcing term $f$ given in (2.39) with $\delta=0.5$ and $M=50$ (left) and the corresponding solution (right).

Recall that we have set $\varepsilon=0$ here, namely only the forcing term is affected by uncertainty, and thus $\eta_{\varepsilon}=\hat{\eta}_{\varepsilon}=0$. We give in Table 2.1 the error $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$ and the estimators $\eta=\left(\eta_{h}^{2}+\eta_{\delta}^{2}\right)^{\frac{1}{2}}$ and $\hat{\eta}=\left(C_{H_{0}^{1}}^{2} \eta_{h}^{2}+\hat{\eta}_{\delta}^{2}\right)^{\frac{1}{2}}$ with $C_{H_{0}^{1}}=1 / 3.46$ for the first case (2.39), where $\eta_{h}, \eta_{\delta}$ and $\hat{\eta}_{\delta}$ are given in (2.26), (2.28) and (2.32), respectively.

|  | $\delta$ | error | $\eta_{\delta}$ | $\eta$ | $\eta /$ error | $\hat{\eta}_{\delta}$ | $\hat{\eta}$ | $\hat{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\aleph$ | $2^{0}$ | $1.1692 \mathrm{e}-1$ | $8.6354 \mathrm{e}-1$ | $8.6357 \mathrm{e}-1$ | 7.3859 | $1.1700 \mathrm{e}-1$ | $1.1702 \mathrm{e}-1$ | 1.0008 |
|  | $2^{-2}$ | $2.9361 \mathrm{e}-2$ | $2.1588 \mathrm{e}-1$ | $2.1603 \mathrm{e}-1$ | 7.3575 | $2.9250 \mathrm{e}-2$ | $2.9337 \mathrm{e}-2$ | 0.9993 |
|  | $2^{-4}$ | $7.6029 \mathrm{e}-3$ | $5.3971 \mathrm{e}-2$ | $5.4534 \mathrm{e}-2$ | 7.1727 | $7.3124 \mathrm{e}-3$ | $7.6531 \mathrm{e}-3$ | 1.0066 |
|  | $2^{-6}$ | $2.9040 \mathrm{e}-3$ | $1.3493 \mathrm{e}-2$ | $1.5591 \mathrm{e}-2$ | 5.3689 | $1.8281 \mathrm{e}-3$ | $2.9052 \mathrm{e}-3$ | 1.0004 |
|  | $2^{-8}$ | $2.3004 \mathrm{e}-3$ | $3.3732 \mathrm{e}-3$ | $8.5096 \mathrm{e}-3$ | 3.6992 | $4.5703 \mathrm{e}-4$ | $2.3037 \mathrm{e}-3$ | 1.0015 |


|  | $\delta$ | error | $\eta_{\delta}$ | $\eta$ | $\eta / \mathrm{error}$ | $\hat{\eta}_{\delta}$ | $\hat{\eta}$ | $\hat{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\curvearrowright$ | $2^{0}$ | $1.1745 \mathrm{e}-1$ | $9.0142 \mathrm{e}-1$ | $9.0146 \mathrm{e}-1$ | 7.6750 | $1.1706 \mathrm{e}-1$ | $1.1708 \mathrm{e}-1$ | 0.9969 |
|  | $2^{-2}$ | $2.9515 \mathrm{e}-2$ | $2.2536 \mathrm{e}-1$ | $2.2549 \mathrm{e}-1$ | 7.6400 | $2.9266 \mathrm{e}-2$ | $2.9353 \mathrm{e}-2$ | 0.9945 |
|  | $2^{-4}$ | $7.6573 \mathrm{e}-3$ | $5.6339 \mathrm{e}-2$ | $5.6878 \mathrm{e}-2$ | 7.4280 | $7.3164 \mathrm{e}-3$ | $7.6569 \mathrm{e}-3$ | 0.9999 |
|  | $2^{-6}$ | $2.8939 \mathrm{e}-3$ | $1.4085 \mathrm{e}-2$ | $1.6106 \mathrm{e}-2$ | 5.5656 | $1.8291 \mathrm{e}-3$ | $2.9058 \mathrm{e}-3$ | 1.0041 |
|  | $2^{-8}$ | $2.2996 \mathrm{e}-3$ | $3.5212 \mathrm{e}-3$ | $8.5694 \mathrm{e}-3$ | 3.7264 | $4.5728 \mathrm{e}-4$ | $2.3038 \mathrm{e}-3$ | 1.0018 |

Table 2.1: Efficiency of the two error estimator $\eta$ and $\hat{\eta}$ for the case (2.39) with $h=2^{-7}$ ( $\eta_{h}=7.8125 \mathrm{e}-3$ ).

We see that similar results are obtained for the two cases $M=6$ and $M=50$. Moreover, the efficiency of the error estimator $\eta$ varies between 3.7 and 7.7. More precisely, we recover the value of $C_{H_{0}^{1}}$ in a physical space error dominant regime while it is about 7.7 when the stochastic error is dominant. The second error estimator $\hat{\eta}$, obtained by taking into account the constant $C_{H_{0}^{1}}$ for $\eta_{h}$ and by computing $M$ additional Poisson problems (see Proposition 2.2.3), yields an effectivity index close to 1 . The results for the second case (2.40), see Figure 2.3 for a plot of some realizations for $f$ and the corresponding solutions, are given in Table 2.2.


Figure 2.3: Seven realizations of the random forcing term $f$ given in (2.40) with $\delta=0.5$ and $M=50$ (left) and the corresponding solution (right).

In this case, the effectivity index of the error estimator $\eta$ is about 4.5 when the stochastic error

Chapter 2. Elliptic model problems with other sources of uncertainty

|  | $\delta$ | error | $\eta_{\delta}$ | $\eta$ | $\eta$ /error | $\hat{\eta}_{\delta}$ | $\hat{\eta}$ | $\hat{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\begin{aligned} & \text { in } \\ & 11 \\ & \sum \end{aligned}$ | $2^{0}$ | $7.2668 \mathrm{e}-2$ | $3.2138 \mathrm{e}-1$ | $3.2148 \mathrm{e}-1$ | 4.4239 | $7.2070 \mathrm{e}-2$ | $7.2106 \mathrm{e}-2$ | 0.9923 |
|  | $2^{-2}$ | $1.8098 \mathrm{e}-2$ | 8.0346e-2 | $8.0725 \mathrm{e}-2$ | 4.4605 | $1.8018 \mathrm{e}-2$ | $1.8159 \mathrm{e}-2$ | 1.0034 |
|  | $2^{-4}$ | $5.0669 \mathrm{e}-3$ | $2.0086 \mathrm{e}-2$ | $2.1552 \mathrm{e}-2$ | 4.2536 | $4.5044 \mathrm{e}-3$ | $5.0386 \mathrm{e}-3$ | 0.9944 |
|  | $2^{-6}$ | $2.5199 \mathrm{e}-3$ | $5.0216 \mathrm{e}-3$ | $9.2872 \mathrm{e}-3$ | 3.6856 | 1.1261e-3 | $2.5232 \mathrm{e}-3$ | 1.0013 |
|  | $2^{-8}$ | $2.2718 \mathrm{e}-3$ | $1.2554 \mathrm{e}-3$ | $7.9127 \mathrm{e}-3$ | 3.4830 | $2.8152 \mathrm{e}-4$ | $2.2754 \mathrm{e}-3$ | 1.0016 |

Table 2.2: Efficiency of the two error estimator $\eta$ and $\hat{\eta}$ for the case (2.40) with $h=2^{-7}$ ( $\eta_{h}=7.8125 \mathrm{e}-3$ ).
is dominant, to be compared to about 7.7 for the first example. This highlight the dependence of the efficiency of $\eta$ with respect to the input data, due to the different loss when using the Poincaré inequality. On the contrary, the second error estimator $\hat{\eta}$ is also very close to 1 for this second example.

## Random forcing term and diffusion coefficient

Let us now consider the case of two random inputs with

$$
\begin{equation*}
a(x, \omega)=1+\varepsilon \sum_{j=1}^{50} a_{j}(x) Y_{j}(\omega), \quad a_{j}(x)=\frac{\sin (2 \pi j x)}{(\pi j)^{2}}, \quad Y_{j} \sim \mathscr{U}[-\sqrt{3}, \sqrt{3}] \tag{2.41}
\end{equation*}
$$

and

$$
f(x, \omega)=1+\delta \sum_{j=1}^{50} f_{j}(x) Z_{j}(\omega), \quad f_{j}(x)=0.5 j^{-\frac{1}{2}} e^{-50 j(x-0.5)^{2}}, \quad Z_{j} \sim \mathscr{N}(0,1) .
$$

Remark 2.3.1. We mention that the choice of the $a_{j}$ in (2.41) is the one for which we obtained the largest effectivity index for the stochastic error estimator $\eta_{\varepsilon}$, namely the ratio of $\eta_{\varepsilon}$ over the error is about 1.8 in the pure stochastic error case (with $\delta=0$ ). It is still an open question, at least to us, to show if there are cases for which we get a larger constant, i.e. for which the loss due to the use of Cauchy-Schwarz inequality in

$$
\int_{D}\left(a-a_{0}\right) \nabla u_{0, h} \cdot \nabla\left(u-u_{0, h}\right) \leq\left\|\left(a-a_{0}\right) \nabla u_{0, h}\right\|_{L^{2}(D)}\left\|\nabla\left(u-u_{0, h}\right)\right\|_{L^{2}(D)}
$$

is bigger.

We give in Tables 2.3 and 2.4 the results obtained for the cases $h=2^{-5}$ and $h=2^{-7}$, respectively. We report the error $\left\|u-u_{0, h}\right\|_{L_{P}^{2}\left(\Omega ; H_{0}^{1}(D)\right)}$, the estimators $\eta_{h}, \eta_{\varepsilon}$ and $\eta_{\delta}$ defined in (2.26), (2.27) and (2.28), respectively, and the effectivity index of the full estimator $\eta=\left(\eta_{h}^{2}+\eta_{\varepsilon}^{2}+\eta_{\delta}^{2}\right)^{\frac{1}{2}}$. We also give the efficiency of the implicit estimator $\hat{\eta}=\left(C_{H_{0}^{1}}^{2} \eta_{h}^{2}+\hat{\eta}_{\varepsilon}^{2}+\hat{\eta}_{\delta}^{2}\right)^{\frac{1}{2}}$ with $\hat{\eta}_{\varepsilon}$ and $\hat{\eta}_{\delta}$ defined in (2.35) and $C_{H_{0}^{1}}=1 / 3.46$.

From the results of Tables 2.3 and 2.4, we see that the efficiency of the full error estimator $\eta$ is

| $\varepsilon$ | $\delta$ | error | $\eta_{\varepsilon}$ | $\eta_{\delta}$ | $\eta / \mathrm{error}$ | $\hat{\eta}_{\varepsilon}$ | $\hat{\eta}_{\delta}$ | $\hat{\eta} / \mathrm{error}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $2^{0}$ | $2^{0}$ | $7.3783 \mathrm{e}-2$ | $1.9954 \mathrm{e}-2$ | $3.2249 \mathrm{e}-1$ | 4.3995 | $1.1331 \mathrm{e}-2$ | $7.2013 \mathrm{e}-2$ | 0.9956 |
| $2^{-2}$ | $2^{0}$ | $7.3191 \mathrm{e}-2$ | $4.9885 \mathrm{e}-3$ | $3.2249 \mathrm{e}-1$ | 4.4272 | $2.8328 \mathrm{e}-3$ | $7.2013 \mathrm{e}-2$ | 0.9924 |
| $2^{-4}$ | $2^{0}$ | $7.2246 \mathrm{e}-2$ | $1.2471 \mathrm{e}-3$ | $3.2249 \mathrm{e}-1$ | 4.4847 | $7.0821 \mathrm{e}-4$ | $7.2013 \mathrm{e}-2$ | 1.0046 |
| $2^{-6}$ | $2^{0}$ | $7.2718 \mathrm{e}-2$ | $3.1178 \mathrm{e}-4$ | $3.2249 \mathrm{e}-1$ | 4.4555 | $1.7705 \mathrm{e}-4$ | $7.2013 \mathrm{e}-2$ | 0.9981 |
|  |  |  |  |  |  |  |  |  |
| $2^{0}$ | $2^{-2}$ | $2.3233 \mathrm{e}-2$ | $1.9954 \mathrm{e}-2$ | $8.0622 \mathrm{e}-2$ | 3.8195 | $1.1331 \mathrm{e}-2$ | $1.8003 \mathrm{e}-2$ | 0.9947 |
| $2^{-2}$ | $2^{-2}$ | $2.0159 \mathrm{e}-2$ | $4.9885 \mathrm{e}-3$ | $8.0622 \mathrm{e}-2$ | 4.2964 | $2.8328 \mathrm{e}-3$ | $1.8003 \mathrm{e}-2$ | 1.0090 |
| $2^{-4}$ | $2^{-2}$ | $2.0186 \mathrm{e}-2$ | $1.2471 \mathrm{e}-3$ | $8.0622 \mathrm{e}-2$ | 4.2840 | $7.0821 \mathrm{e}-4$ | $1.8003 \mathrm{e}-2$ | 0.9984 |
| $2^{-6}$ | $2^{-2}$ | $2.0131 \mathrm{e}-2$ | $3.1178 \mathrm{e}-4$ | $8.0622 \mathrm{e}-2$ | 4.2952 | $1.7705 \mathrm{e}-4$ | $1.8003 \mathrm{e}-2$ | 1.0006 |
|  |  |  |  |  |  |  |  |  |
| $2^{0}$ | $2^{-4}$ | $1.5453 \mathrm{e}-2$ | $1.9954 \mathrm{e}-2$ | $2.0155 \mathrm{e}-2$ | 2.7309 | $1.1331 \mathrm{e}-2$ | $4.5008 \mathrm{e}-3$ | 0.9819 |
| $2^{-2}$ | $2^{-4}$ | $1.0487 \mathrm{e}-2$ | $4.9885 \mathrm{e}-3$ | $2.0155 \mathrm{e}-2$ | 3.5776 | $2.8328 \mathrm{e}-3$ | $4.5008 \mathrm{e}-3$ | 0.9994 |
| $2^{-4}$ | $2^{-4}$ | $1.0114 \mathrm{e}-2$ | $1.2471 \mathrm{e}-3$ | $2.0155 \mathrm{e}-2$ | 3.6789 | $7.0821 \mathrm{e}-4$ | $4.5008 \mathrm{e}-3$ | 1.0002 |
| $2^{-6}$ | $2^{-4}$ | $1.0068 \mathrm{e}-2$ | $3.1178 \mathrm{e}-4$ | $2.0155 \mathrm{e}-2$ | 3.6937 | $1.7705 \mathrm{e}-4$ | $4.5008 \mathrm{e}-3$ | 1.0025 |
|  |  |  |  |  |  |  |  |  |
| $2^{0}$ | $2^{-6}$ | $1.4804 \mathrm{e}-2$ | $1.9954 \mathrm{e}-2$ | $5.0388 \mathrm{e}-3$ | 2.5276 | $1.1331 \mathrm{e}-2$ | $1.1252 \mathrm{e}-3$ | 0.9818 |
| $2^{-2}$ | $2^{-6}$ | $9.5369 \mathrm{e}-3$ | $4.9885 \mathrm{e}-3$ | $5.0388 \mathrm{e}-3$ | 3.3600 | $2.8328 \mathrm{e}-3$ | $1.1252 \mathrm{e}-3$ | 0.9995 |
| $2^{-4}$ | $2^{-6}$ | $9.1184 \mathrm{e}-3$ | $1.2471 \mathrm{e}-3$ | $5.0388 \mathrm{e}-3$ | 3.4741 | $7.0821 \mathrm{e}-4$ | $1.1252 \mathrm{e}-3$ | 1.0012 |
| $2^{-6}$ | $2^{-6}$ | $9.0934 \mathrm{e}-3$ | $3.1178 \mathrm{e}-4$ | $5.0388 \mathrm{e}-3$ | 3.4811 | $1.7705 \mathrm{e}-4$ | $1.1252 \mathrm{e}-3$ | 1.0011 |

Table 2.3: for $h=2^{-5}\left(\eta_{h}=3.125 \mathrm{e}-2\right)$

| $\varepsilon$ | $\delta$ | error | $\eta_{\varepsilon}$ | $\eta_{\delta}$ | $\eta /$ error | $\hat{\eta}_{\varepsilon}$ | $\hat{\eta}_{\delta}$ | $\hat{\eta} /$ error |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $2^{0}$ | $2^{0}$ | $7.3022 \mathrm{e}-2$ | $1.9923 \mathrm{e}-2$ | $3.2138 \mathrm{e}-1$ | 4.4109 | $1.1490 \mathrm{e}-2$ | $7.2070 \mathrm{e}-2$ | 0.9999 |
| $2^{-2}$ | $2^{0}$ | $7.2376 \mathrm{e}-2$ | $4.9806 \mathrm{e}-3$ | $3.2138 \mathrm{e}-1$ | 4.4423 | $2.8724 \mathrm{e}-3$ | $7.2070 \mathrm{e}-2$ | 0.9971 |
| $2^{-4}$ | $2^{0}$ | $7.2361 \mathrm{e}-2$ | $1.2452 \mathrm{e}-3$ | $3.2138 \mathrm{e}-1$ | 4.4428 | $7.1811 \mathrm{e}-4$ | $7.2070 \mathrm{e}-2$ | 0.9965 |
| $2^{-6}$ | $2^{0}$ | $7.1792 \mathrm{e}-2$ | $3.1129 \mathrm{e}-4$ | $3.2138 \mathrm{e}-1$ | 4.4779 | $1.7953 \mathrm{e}-4$ | $7.2070 \mathrm{e}-2$ | 1.0044 |
|  |  |  |  |  |  |  |  |  |
| $2^{0}$ | $2^{-2}$ | $2.1710 \mathrm{e}-2$ | $1.9923 \mathrm{e}-2$ | $8.0346 \mathrm{e}-2$ | 3.8299 | $1.1490 \mathrm{e}-2$ | $1.8018 \mathrm{e}-2$ | 0.9898 |
| $2^{-2}$ | $2^{-2}$ | $1.8452 \mathrm{e}-2$ | $4.9806 \mathrm{e}-3$ | $8.0346 \mathrm{e}-2$ | 4.3832 | $2.8724 \mathrm{e}-3$ | $1.8018 \mathrm{e}-2$ | 0.9963 |
| $2^{-4}$ | $2^{-2}$ | $1.8183 \mathrm{e}-2$ | $1.2452 \mathrm{e}-3$ | $8.0346 \mathrm{e}-2$ | 4.4401 | $7.1811 \mathrm{e}-4$ | $1.8018 \mathrm{e}-2$ | 0.9994 |
| $2^{-6}$ | $2^{-2}$ | $1.7873 \mathrm{e}-2$ | $3.1129 \mathrm{e}-4$ | $8.0346 \mathrm{e}-2$ | 4.5165 | $1.7953 \mathrm{e}-4$ | $1.8018 \mathrm{e}-2$ | 1.0160 |
|  |  |  |  |  |  |  |  |  |
| $2^{0}$ | $2^{-4}$ | $1.2685 \mathrm{e}-2$ | $1.9923 \mathrm{e}-2$ | $2.0086 \mathrm{e}-2$ | 2.3138 | $1.1490 \mathrm{e}-2$ | $4.5044 \mathrm{e}-3$ | 0.9890 |
| $2^{-2}$ | $2^{-4}$ | $5.7768 \mathrm{e}-3$ | $4.9806 \mathrm{e}-3$ | $2.0086 \mathrm{e}-2$ | 3.8291 | $2.8724 \mathrm{e}-3$ | $4.5044 \mathrm{e}-3$ | 1.0040 |
| $2^{-4}$ | $2^{-4}$ | $5.0541 \mathrm{e}-3$ | $1.2452 \mathrm{e}-3$ | $2.0086 \mathrm{e}-2$ | 4.2715 | $7.1811 \mathrm{e}-4$ | $4.5044 \mathrm{e}-3$ | 1.0070 |
| $2^{-6}$ | $2^{-4}$ | $5.0988 \mathrm{e}-3$ | $3.1129 \mathrm{e}-4$ | $2.0086 \mathrm{e}-2$ | 4.2274 | $1.7953 \mathrm{e}-4$ | $4.5044 \mathrm{e}-3$ | 0.9888 |
|  |  |  |  |  |  |  |  |  |
| $2^{0}$ | $2^{-6}$ | $1.1897 \mathrm{e}-2$ | $1.9923 \mathrm{e}-2$ | $5.0216 \mathrm{e}-3$ | 1.8476 | $1.1490 \mathrm{e}-2$ | $1.1261 \mathrm{e}-3$ | 0.9888 |
| $2^{-2}$ | $2^{-6}$ | $3.8361 \mathrm{e}-3$ | $4.9806 \mathrm{e}-3$ | $5.0216 \mathrm{e}-3$ | 2.7471 | $2.8724 \mathrm{e}-3$ | $1.1261 \mathrm{e}-3$ | 0.9966 |
| $2^{-4}$ | $2^{-6}$ | $2.6217 \mathrm{e}-3$ | $1.2452 \mathrm{e}-3$ | $5.0216 \mathrm{e}-3$ | 3.5742 | $7.1811 \mathrm{e}-4$ | $1.1261 \mathrm{e}-3$ | 1.0007 |
| $2^{-6}$ | $2^{-6}$ | $2.5186 \mathrm{e}-3$ | $3.1129 \mathrm{e}-4$ | $5.0216 \mathrm{e}-3$ | 3.6895 | $1.7953 \mathrm{e}-4$ | $1.1261 \mathrm{e}-3$ | 1.0043 |

Table 2.4: for $h=2^{-7}\left(\eta_{h}=7.8125 \mathrm{e}-3\right)$
comprised between the efficiency of each of its parts $\eta_{h}, \eta_{\varepsilon}$ and $\eta_{\delta}$, depending on which is the predominant source of error. For instance, the effectivity index tends to the value 3.46 when the FE error is dominant, see e.g. Table 2.3 with $\delta=2^{-6}$ and $\varepsilon=2^{-4}$ or $2^{-6}$, while it is about 4.5 in a $\delta$-error dominant regime as in the similar case (2.40) considered above. Finally, we see that if the error due to the uncertainty in the diffusion coefficient $a$ is largest, the effectivity index tends to the value 1.8 indicated in Remark 2.3.1. In all cases, the implicit error estimator $\hat{\eta}$ has an effectivity index close to 1 , but more work is required to compute it.

## Conclusions

In this chapter, we have extended the results we obtained in Chapter 1 for the linear model problem to include other sources of uncertainty. More precisely, we have considered first the case of Neumann random boundary conditions and then the combination of two random input data, namely the diffusion coefficient and the forcing term. For the latter case, two different sets of random variables have been used to characterize the data affected by uncertainty. We have shown that when the random solution depends linearly on the random input, as it is the case for Neumann boundary conditions or the source term, then the solution is fully described by the first two terms in the expansion, the remaining terms being zero. Moreover, we have seen that when the Poincaré inequality is required in the estimation, the efficiency of the error estimator might change when modifying the input data, even though it has the optimal convergence rate. The same behaviour was observed when considering the error $u-u_{0, h}$ in the $L^{2}(D)$ norm in Chapter 1 . As a remedy to the sensitivity of the error estimator to the input data, we have proposed a second error estimator, see Proposition 2.2.3. It is obtained by solving additional (Poisson) problems, as many as the number of random variables used to characterized the uncertainty in the data. However, we can use the same spatial mesh than the one for $u_{0, h}$ to solve these problems approximately.

## 3 PDEs in random domains

In this chapter, we consider nonlinear PDEs defined on random domains. The first part consists of the analysis of a 1D problem, namely the viscous Burgers' equation to be solved on an interval of random length. This equation, first introduced by Bateman in [19] and then used by Burgers in [34] for modelling turbulence, can be seen as a simplification of the Navier-Stokes equations to the one-dimensional case. In the second part, whose material in mainly taken from the submitted paper [75], we consider the more involved incompressible Navier-Stokes equations in random domains. We restrict ourselves here to the stationary formulation of these equations.

For both problems, we use the so-called domain mapping method [125]: we introduce a random mapping that transforms the deterministic PDEs defined on a random domain into PDEs on a fixed reference domain with random coefficients. For simplicity, we assume that the uncertainty in the system is only due to the random domain, but the analysis can be straightforwardly extended to include other sources of randomness.

## Introduction

Several approaches have been developed to perform analysis and numerical approximation of PDEs in random domains, such as the fictitious domain method [40], the perturbation method based on shape calculus [77] and the domain mapping method initially proposed by [125] and also used for instance in [39, 43, 76]. In the first approach, the PDEs are extended to a fixed reference domain, the so-called fictitious domain, which contains all the random domains. The original boundary condition is then imposed through a Lagrange multiplier yielding a saddle-point problem to be solved in the fictitious domain. In the perturbation method, which is suitable for small perturbations only, the solution is represented using a shape Taylor expansion with respect to the (random) perturbation field of the boundary of the domain. Finally, the domain mapping approach, which is the one considered in this work, transforms the deterministic PDEs defined on a random domain into PDEs on a fixed reference domain with random coefficients via a random mapping. We give in Figure 3.1 an
illustration of the mapping for a given $\omega$ between the physical domain and the reference one, supplemented with some notation.


Figure 3.1: Illustration and notation for the domain mapping approach.

Contrary to the method based on shape derivatives, our approach requires the construction of a random mapping defined in the whole domain consistent with the random perturbation of the boundary. If the random mapping is not given analytically, it can be obtained by solving appropriate equations, e.g. Laplace equation as it is done in [125]. The domain mapping method prevents the need of remeshing and can make use of the well-developed theory for PDEs on deterministic domains with random coefficients. Numerical approximation of the solution on the fixed reference domain can indeed be obtained through any of the well-known techniques, such as Monte-Carlo methods [63] and their generalizations as quasi-Monte Carlo [38,54,70] and multi-level Monte-Carlo [17,52,68,79], or the stochastic spectral methods comprising the stochastic Galerkin [10, 11, 21, 64, 67] and the stochastic collocation [7,97, 124] methods.

The (weak) formulation on the reference domain can be obtained using two strategies, as illustrated in Figure 3.2. In general, the two strategies are not equivalent. They yield the same result only in particular cases, for instance if the Jacobian of the mapping does not depend on the physical space variable. In this work, we will use the first strategy $s_{1}$, that is the formulation on the reference domain is obtained performing the change of variables on the weak formulation of the problem on the random domain. We refer for instance to [36] for a version where the second strategy $s_{2}$ is used.


Figure 3.2: Two strategies $s_{1}$ and $s_{2}$ for the (strong) formulation on the reference domain.

For the stochastic space approximation, we proceed as in the previous chapters and use a perturbation approach [82] to expand the exact random solution with respect to a parameter $\varepsilon$ that controls the level of uncertainty in the problem. This approach yields uncoupled deterministic problems for each term in the expansion, which can be solved using for instance the finite element (FE) method. The main goal here is to perform an a posteriori error analysis for the error between the exact random solution and the finite element approximation of the first term in the expansion, that is the solution corresponding to the case $\varepsilon=0$. The error estimators we obtain are made of two parts, namely one part due to the physical space discretization and another one due to the uncertainty. Their computation requires only the FE approximation of the solution of the problem for $\varepsilon=0$ and the Jacobian matrix of the mapping between the reference domain and the physical random domain. These estimators can be used for instance to adaptively determine a mesh that yields a numerical accuracy comparable with the model uncertainty. Notice that the error estimates we get here using the domain mapping method combined with a perturbation technique are defined for any fixed $\varepsilon$. The only restriction is that $\varepsilon$ is sufficiently small for the problem to be well-posed. The more common perturbation method is to use shape calculus [77], thus avoiding to recast the equations in a reference domain. However, the derivation of a posteriori error estimates for a fixed value of $\varepsilon$ is, in our opinion, not obvious in this context and, to the best of our knowledge, it is still an open question.

We mention that the formulation we obtained in Section 3.2.2 for the Navier-Stokes equations on the reference domain is similar to the one obtained for instance in [71] where a fluidstructure interaction problem is considered or in $[91,108]$ where the Navier-Stokes equations in parametrized domains are solved approximately using the Reduced Basis Method.

### 3.1 Steady-state viscous Burgers' equation in random intervals

To start with, we consider a 1D problem on a random domain, namely the (nonlinear) steadystate viscous Burgers' equation. This equation can be viewed as a simplification of the NavierStokes equations in the one-dimensional case. We consider a physical domain with uncertain geometry, which reduces here to an interval of random length. We study first the deterministic case, considering the Burgers' equation on a fixed domain, say $[0,1]$.

### 3.1.1 Deterministic case

We consider the following nonlinear deterministic problem with mixed Neumann-Dirichlet homogeneous boundary conditions:
find $u:(0,1) \rightarrow \mathbb{R}$ such that

$$
\left\{\begin{align*}
-a u^{\prime \prime}+b u u^{\prime} & =f \text { in }(0,1)  \tag{3.1}\\
u(0) & =0 \\
u^{\prime}(1) & =0
\end{align*}\right.
$$

where $a$ and $b$ are positive constants and $f \in L^{2}(0,1)$. It can be written in conservation form as

$$
-a u^{\prime \prime}+\frac{b}{2}\left(u^{2}\right)^{\prime}=f \quad \text { in }(0,1)
$$

Let $V=\left\{v \in H^{1}(0,1): v(0)=0\right\}$ that we endow with the norm $\|\cdot\|_{V}:=|\cdot|_{H^{1}(0,1)}$. This is possible thanks to the Friedrich-Poincaré inequality, see (2.5), which reads here

$$
\begin{equation*}
\|u\|_{L^{2}(0,1)} \leq C_{F}\left\|u^{\prime}\right\|_{L^{2}(0,1)} \tag{3.2}
\end{equation*}
$$

and holds for instance for $C_{F}=\frac{1}{\sqrt{2}} \leq 1$. The weak form of problem (3.1) is given by

$$
\begin{equation*}
\text { find } u \in V: \int_{0}^{1} a u^{\prime} v^{\prime} d x+\int_{0}^{1} b u u^{\prime} v d x=\int_{0}^{1} f v d x \quad \forall v \in V \text {. } \tag{3.3}
\end{equation*}
$$

We first show, under suitable conditions on the data, that the problem (3.1) is well-posed. Since we do not have an a priori estimate, due to the mixed Neumann-Dirichlet boundary conditions ${ }^{1}$, we restrict ourselves to the set of functions whose norm is bounded by a certain constant. More precisely, we consider

$$
\mathscr{M}:=\left\{v \in V:\left\|v^{\prime}\right\|_{L^{2}(0,1)} \leq r\right\} \quad \text { with } \quad r=\sqrt{\frac{C_{F}}{b}\|f\|_{L^{2}(0,1)}}
$$

Since $\mathscr{M}$ is a closed ball in $V$, it is bounded, convex and closed in $V$. The well-posedness of the problem under certain assumptions on the data is proved in the following proposition.

Proposition 3.1.1. If $\frac{a}{b} \geq 2 r$, then there exists a solution $u \in \mathscr{M}$ to problem (3.3). Moreover, if $\frac{a}{b}>2 r$, then such solution is unique.

Strictly speaking, it is enough to assume $\frac{a}{b}>r$ to prove the existence of a solution in $\mathscr{M}$. Using the definition of $r$, the condition $\frac{a}{b}>2 r$ can be expressed more explicitly in terms of the given

[^3]data by $C_{F}\|f\|_{L^{2}(0,1)}<\frac{a^{2}}{4 b}$, which coincides with the one given in [28] replacing $C_{F}\|f\|_{L^{2}(0,1)}$ by the dual norm $\|f\|_{V^{\prime}}$. The proof of Proposition 3.1.1, given below for completeness, uses the Schauder's fixed point theorem for the existence and is inspired by the one given in [120]. The uniqueness is proved using a variational argument.

Proof. Existence: we define the mapping $T: \mathscr{M} \rightarrow V, u \mapsto T u=: w$, where $w \in V$ is the unique solution of

$$
\begin{equation*}
\text { find } w \in V: \mathscr{A}_{u}(w, v)=F(v) \quad \forall v \in V \tag{3.4}
\end{equation*}
$$

with

$$
\mathscr{A}_{u}(w, v):=\int_{0}^{1} a w^{\prime} v^{\prime} d x+\int_{0}^{1} b u w^{\prime} v d x \text { and } F(v):=\int_{0}^{1} f v d x
$$

We show that $T$ is well-defined, maps $\mathscr{M}$ to $\mathscr{M}$ and is compact. Let $u \in \mathscr{M}$, i.e. $\left\|u^{\prime}\right\|_{L^{2}(0,1)} \leq r$. The fact that $T: \mathscr{M} \rightarrow V$ is well-defined follows directly from Lax-Milgram's lemma. Indeed, for any $v, w \in V$ we have
$\mathscr{A}_{u}(w, v) \leq a\left\|w^{\prime}\right\|_{L^{2}(0,1)}\left\|v^{\prime}\right\|_{L^{2}(0,1)}+b\|u\|_{L^{4}(0,1)}\left\|w^{\prime}\right\|_{L^{2}(0,1)}\|v\|_{L^{4}(0,1)} \leq(a+b r)\left\|w^{\prime}\right\|_{L^{2}(0,1)}\left\|v^{\prime}\right\|_{L^{2}(0,1)}$
using successively Cauchy-Schwarz and Hölder's inequalities and the fact that

$$
\begin{equation*}
\|v\|_{L^{4}(0,1)} \leq C\left\|\nu^{\prime}\right\|_{L^{2}(0,1)} \quad \text { holds with } C=1 \tag{3.5}
\end{equation*}
$$

Moreover, since $u \in \mathscr{M}$ and $\frac{a}{b} \geq 2 r$ by assumption, we have

$$
-\int_{0}^{1} b u w^{\prime} w d x \leq b\left\|u^{\prime}\right\|_{L^{2}(0,1)}\left\|w^{\prime}\right\|_{L^{2}(0,1)}^{2} \leq b r\left\|w^{\prime}\right\|_{L^{2}(0,1)}^{2} \leq \frac{a}{2}\left\|w^{\prime}\right\|_{L^{2}(0,1)}^{2}
$$

and thus

$$
\mathscr{A}_{u}(w, w)=a\left\|w^{\prime}\right\|_{L^{2}(0,1)}^{2}+\int_{0}^{1} b u w^{\prime} w d x \geq \frac{a}{2}\left\|w^{\prime}\right\|_{L^{2}(0,1)}^{2}
$$

Finally, thanks to (3.2) we get

$$
F(v) \leq C_{F}\|f\|_{L^{2}(0,1)}\left\|v^{\prime}\right\|_{L^{2}(0,1)}
$$

and the assumptions of Lax-Milgram's lemma are satisfied. We now show that $T$ maps $\mathscr{M}$ to itself, i.e. $T u=w \in \mathscr{M}$. Thanks to the coercivity of $\mathscr{A}_{u}$ and the continuity of $F$, taking $v=w$ in (3.4) yields

$$
\frac{a}{2}\left\|w^{\prime}\right\|_{L^{2}(0,1)}^{2} \leq \int_{0}^{1} f w d x \leq C_{F}\|f\|_{L^{2}(0,1)}\left\|w^{\prime}\right\|_{L^{2}(0,1)}
$$

and thus

$$
\left\|w^{\prime}\right\|_{L^{2}(0,1)} \leq \frac{2}{a} C_{F}\|f\|_{L^{2}(0,1)}=\frac{2}{a} b r^{2} \leq r
$$

We finally show that $T$ is compact. Let $\left(u_{n}\right)_{n \in \mathbb{N}}$ be a bounded sequence in $\mathscr{M}$. Since $H^{1}(0,1)$ is compactly embedded in $L^{4}(0,1)$, there exists a subsequence $\left(u_{n_{j}}\right)_{j \in \mathbb{N}}$ which converges in $L^{4}(0,1)$. Let $u_{n}$ and $u_{m}$ be two elements of this subsequence and write $w_{n}$ and $w_{m}$ the
corresponding images under $T$. We have

$$
\int_{0}^{1} a\left(w_{n}^{\prime}-w_{m}^{\prime}\right) v^{\prime} d x+\int_{0}^{1} b\left[u_{n}\left(w_{n}^{\prime}-w_{m}^{\prime}\right)+w_{m}^{\prime}\left(u_{n}-u_{m}\right)\right] v d x=0 \quad \forall v \in V .
$$

If we take $v=w_{n}-w_{m}$, using that $u_{n} \in \mathscr{M}$ and $b r \leq \frac{a}{2}$ we can easily show that

$$
\frac{a}{2}\left\|w_{n}^{\prime}-w_{m}^{\prime}\right\|_{L^{2}(0,1)} \leq b\left\|w_{m}^{\prime}\right\|_{L^{2}(0,1)}\left\|u_{n}-u_{m}\right\|_{L^{4}(0,1)}
$$

and thus

$$
\left\|w_{n}^{\prime}-w_{m}^{\prime}\right\|_{L^{2}(0,1)} \leq\left\|u_{n}-u_{m}\right\|_{L^{4}(0,1)}
$$

since $w_{m} \in \mathscr{M}$. Therefore, $\left(w_{n_{j}}\right)_{j \in \mathbb{N}}$ is a Cauchy sequence in $V$ and thus converges.
Uniqueness: we use a variational argument. Let $u_{1}, u_{2} \in \mathscr{M}$ be two solutions of problem (3.3). We have

$$
\int_{0}^{1} a\left(u_{1}^{\prime}-u_{2}^{\prime}\right) v d x+\int_{0}^{1} b\left(u_{1} u_{1}^{\prime}-u_{2} u_{2}^{\prime}\right) v d x=0 \quad \forall v \in V .
$$

If we take $v=u_{1}-u_{2}$, we obtain

$$
\begin{aligned}
a\left\|u_{1}^{\prime}-u_{2}^{\prime}\right\|_{L^{2}(0,1)}^{2}= & -\int_{0}^{1} b\left(u_{1}\left(u_{1}^{\prime}-u_{2}^{\prime}\right)+u_{2}^{\prime}\left(u_{1}-u_{2}\right)\right)\left(u_{1}-u_{2}\right) d x \\
\leq & b\left\|u_{1}\right\|_{L^{4}(0,1)}\left\|u_{1}-u_{2}\right\|_{L^{2}(0,1)}\left\|u_{1}-u_{2}\right\|_{L^{4}(0,1)} \\
& +b\left\|u_{2}^{\prime}\right\|_{L^{2}(0,1)}\left\|u_{1}-u_{2}\right\|_{L^{4}(0,1)}^{2} \\
\leq & b\left(\left\|u_{1}^{\prime}\right\|_{L^{2}(0,1)}+\left\|u_{2}^{\prime}\right\|_{L^{2}(0,1)}\right)\left\|u_{1}^{\prime}-u_{2}^{\prime}\right\|_{L^{2}(0,1)}^{2} \\
\leq & 2 b r\left\|u_{1}^{\prime}-u_{2}^{\prime}\right\|_{L^{2}(0,1)}^{2}
\end{aligned}
$$

and thus

$$
(a-2 b r)\left\|u_{1}^{\prime}-u_{2}^{\prime}\right\|_{L^{2}(0,1)}^{2} \leq 0
$$

Since $\frac{a}{b}>2 r$ by assumption, the last inequality implies $u_{1}^{\prime}=u_{2}^{\prime}$. The fact that $u_{1}(0)=u_{2}(0)$ allows us to conclude that $u_{1}=u_{2}$.

Remark 3.1.2. If the solution is assumed to be in $H^{2}(0,1)$, we can alternatively use Schaefer's fixed point theorem [62] to prove the existence of a solution to problem (3.1).

We now give an a posteriori estimate of the error in the $V$ norm between the exact solution $u$ and its finite element approximation. We thus consider

$$
0=x_{0}<x_{1}<\ldots<x_{N}<x_{N+1}=1
$$

a partition of $[0,1]$ and let $h_{i}=x_{i+1}-x_{i}$ for $i=0, \ldots, N$. Let $V_{h} \subset V$ be the finite dimensional space of continuous piecewise polynomials of degree less or equal to one associated to this
partition (the usual hat functions). The finite element approximation of problem (3.3) reads

$$
\begin{equation*}
\text { find } u_{h} \in V_{h}: \int_{0}^{1} a u_{h}^{\prime} v_{h}^{\prime} d x+\int_{0}^{1} b u_{h} u_{h}^{\prime} v_{h} d x=\int_{0}^{1} f v_{h} d x \quad \forall v_{h} \in V_{h} \text {. } \tag{3.6}
\end{equation*}
$$

Similarly to the continuous case, we can show that there exists a unique solution $u_{h} \in \mathscr{M}_{h}$ to problem (3.6) if $\frac{a}{b}>2 r$, with $\mathscr{M}_{h}=\left\{v_{h} \in V_{h}:\left\|v_{h}^{\prime}\right\|_{L^{2}(0,1)} \leq r\right\} \subset \mathscr{M}$. Moreover, if we take $v=v_{h}$ in (3.3) and subtract (3.6), we get the following so-called Galerkin orthogonality property

$$
\begin{equation*}
\int_{0}^{1} a\left(u^{\prime}-u_{h}^{\prime}\right) v_{h}^{\prime} d x+\int_{0}^{1} b\left(u u^{\prime}-u_{h} u_{h}^{\prime}\right) v_{h} d x=0 \quad \forall v_{h} \in V_{h} . \tag{3.7}
\end{equation*}
$$

Proposition 3.1.3. If $a, b$ and $f$ are such that $\frac{a}{b}>2 r$, i.e. $\frac{4 b}{a^{2}} C_{F}\|f\|_{L^{2}(0,1)}<1$, then there exists a constant $C>0$ independent of $h$ and $u$ such that

$$
\begin{equation*}
\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L^{2}(0,1)} \leq \frac{C}{a}\left(\sum_{i=0}^{N} \eta_{i}^{2}\right)^{\frac{1}{2}} \tag{3.8}
\end{equation*}
$$

with

$$
\begin{equation*}
\eta_{i}^{2}=h_{i}^{2} \int_{x_{i}}^{x_{i+1}}\left(f-b u_{h} u_{h}^{\prime}+a u_{h}^{\prime \prime}\right)^{2} d x, \quad i=0, \ldots, N \tag{3.9}
\end{equation*}
$$

Proof. For any $v \in V$, let $\left\langle\mathscr{R}\left(u_{h}\right), v\right\rangle=\int_{0}^{1}\left(f v-b u_{h} u_{h}^{\prime} v-a u_{h}^{\prime} \nu^{\prime}\right) d x$ denote the residual for $u_{h}$. We have

$$
\begin{aligned}
\int_{0}^{1} a\left(u^{\prime}-u_{h}^{\prime}\right) v d x & =\int_{0}^{1} f v d x-\int_{0}^{1} b u u^{\prime} v d x-\int_{0}^{1} a u_{h}^{\prime} \nu^{\prime} d x \\
& =\left\langle\mathscr{R}\left(u_{h}\right), v\right\rangle-\int_{0}^{1} b\left(u u^{\prime}-u_{h} u_{h}^{\prime}\right) v d x .
\end{aligned}
$$

If we take $v=u-u_{h}$, the second term can be bounded by

$$
-\int_{0}^{1} b\left(u u^{\prime}-u_{h} u_{h}^{\prime}\right) v d x \leq 2 b r\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L^{2}(0,1)}^{2}
$$

Therefore

$$
\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L^{2}(0,1)}^{2} \leq \frac{1}{a}\left\langle\mathscr{R}\left(u_{h}\right), u-u_{h}\right\rangle+\frac{2 b r}{a}\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L^{2}(0,1)}^{2}
$$

Since $\frac{a}{b}>2 r$ by assumption, there exists $\gamma>0$ such that $\frac{2 b r}{a} \leq 1-\gamma$. Therefore, we have

$$
\begin{equation*}
\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L^{2}(0,1)}^{2} \leq \frac{1}{a \gamma} \mathscr{R}\left(u-u_{h}\right) . \tag{3.10}
\end{equation*}
$$

It only remains to give an estimation of the residual. First note that

$$
\left\langle\mathscr{R}\left(u_{h}\right), v_{h}\right\rangle=0 \quad \forall v_{h} \in V_{h} .
$$

Taking $v_{h}=r_{h} v$ the Lagrange interpolant of $v$ and using standard techniques, we get

$$
\begin{equation*}
\left\langle\mathscr{R}\left(u_{h}\right), v\right\rangle \leq C_{I}\left(\sum_{i=0}^{N} h_{i}^{2} \int_{x_{i}}^{x_{i+1}}\left(f-b u_{h} u_{h}^{\prime}+a u_{h}^{\prime \prime}\right)^{2} d x\right)^{\frac{1}{2}}\left\|v^{\prime}\right\|_{L^{2}(0,1)} \tag{3.11}
\end{equation*}
$$

where $C_{I}$ is the constant (independent of $h$ and $v$ ) in the interpolation error bound

$$
\begin{equation*}
\left\|v-r_{h} v\right\|_{L^{2}\left(x_{i}, x_{i+1}\right)} \leq C_{I} h_{i}\left\|\nu^{\prime}\right\|_{L^{2}\left(x_{i}, x_{i+1}\right)} . \tag{3.12}
\end{equation*}
$$

For instance, we can take $C_{I}=\sqrt{\frac{49}{30}}$. Inserting (3.11) in (3.10) yields (3.8) with $C=\frac{C_{I}}{\gamma}$.
Remark 3.1.4. The a posteriori error estimate (3.8) holds under the constraint $\frac{2 b r}{a}<1$, i.e. $\frac{2 b r}{a} \leq 1-\gamma$ for a certain $\gamma>0$. However, if $\gamma$ is chosen too small, then the constant $C$ explodes. In practice, it is common to assume that the input data are such that $\frac{2 b r}{a} \leq \frac{1}{2}$ holds.

### 3.1.2 Random case

Let $(\Omega, \mathscr{F}, P)$ be a complete probability space and for any $\omega \in \Omega$ let $D_{\omega}:=(0, s(w)) \subseteq \hat{D}$ be an interval of random length $s(w)$. To simplify the notation, the set

$$
\left\{(x, \omega): x \in D_{\omega}, \omega \in \Omega\right\}
$$

will be denoted by $D_{\omega} \times \Omega$ in the sequel. The goal is to solve the problem:
find $\tilde{u}: D_{w} \times \Omega \rightarrow \mathbb{R}$ such that a.s. in $\Omega$

$$
\left\{\begin{align*}
-a \frac{\partial^{2}}{\partial x^{2}} \tilde{u}(x, \omega)+b \tilde{u}(x, \omega) \frac{\partial}{\partial x} \tilde{u}(x, \omega) & =\tilde{f}(x) \quad x \in D_{\omega}  \tag{3.13}\\
\tilde{u}(0, \omega) & =0 \\
\frac{\partial}{\partial x} \tilde{u}(s(\omega), \omega) & =0
\end{align*}\right.
$$

where $a$ and $b$ are positive constants and $\tilde{f} \in L^{2}(\hat{D})$ is a deterministic forcing term. Let $\tilde{V}_{\omega}=\left\{\tilde{v} \in H^{1}\left(D_{w}\right): \tilde{v}(0, \omega)=0\right.$ a.s. in $\left.\Omega\right\}$. The pointwise weak form of problem (3.13) reads:
find $\tilde{u}(\cdot, \omega) \in \tilde{V}_{\omega}$ such that

$$
\begin{equation*}
\int_{0}^{s(\omega)} a \frac{\partial \tilde{u}(\cdot, \omega)}{\partial x} \frac{\partial \tilde{v}}{\partial x} d x+\int_{0}^{s(\omega)} b \tilde{u}(\cdot, \omega) \frac{\partial \tilde{u}(\cdot, \omega)}{\partial x} \tilde{v} d x=\int_{0}^{s(\omega)} \tilde{f} \tilde{v} d x \quad \forall \tilde{v} \in \tilde{V}_{\omega} \tag{3.14}
\end{equation*}
$$

For ease of presentation, we will use the short hand notation $\tilde{u}(\omega)=\tilde{u}(\cdot, \omega)$ when no confusion arises. Instead of solving this problem on the stochastic domain $D_{\omega}$, we will solve it on a fixed reference domain, namely $D=(0,1)$, by considering the change of variable $x=s(\omega) \xi$. Therefore, assuming $s(\omega)>0$ a.s. in $\Omega$ we define the (random) mapping

$$
\begin{align*}
g_{\omega}: \quad D_{\omega} & \rightarrow D \\
x & \mapsto \xi=g_{\omega}(x)=\frac{x}{s(\omega)} \tag{3.15}
\end{align*}
$$

whose inverse is given by

$$
\begin{aligned}
g_{\omega}^{-1}: \quad D & \rightarrow D_{\omega} \\
\xi & \mapsto x=g_{\omega}^{-1}(\xi)=s(\omega) \xi .
\end{aligned}
$$

Let $u(\xi, \omega)=\tilde{u}(x, \omega)$ and $f(\xi, \omega)=\tilde{f}(x, \omega)$ denote respectively the velocity and the forcing term on the fixed domain $D$, i.e. $u(\xi, \omega)=\tilde{u}\left(g_{\omega}^{-1}(\xi), \omega\right)$ and $f(\xi, \omega)=\tilde{f}\left(g_{\omega}^{-1}(\xi)\right)$. Finally, let $V=\left\{\nu \in H^{1}(D): \nu(0)=0\right\}$. Applying the standard chain rule and the change of variable formula, the pointwise weak problem (3.14) can then be rewritten:
find $u(\omega) \in V$ such that

$$
\begin{equation*}
\int_{0}^{1} \frac{a}{s(\omega)} \frac{\partial u(\omega)}{\partial \xi} \frac{\partial v}{\partial \xi} d \xi+\int_{0}^{1} b u(\omega) \frac{\partial u(\omega)}{\partial \xi} v d \xi=\int_{0}^{1} s(\omega) f(\omega) v d \xi \quad \forall v \in V . \tag{3.16}
\end{equation*}
$$

The strong form of the problem on the reference domain can be stated as:
find $u: D \times \Omega \rightarrow \mathbb{R}$ such that a.s. in $\Omega$

$$
\left\{\begin{align*}
-\frac{a}{s(\omega)^{2}} \frac{\partial^{2}}{\partial \xi^{2}} u(\xi, \omega)+\frac{b}{s(\omega)} u(\xi, \omega) \frac{\partial}{\partial \xi} u(\xi, \omega) & =f(\xi, \omega) \quad \xi \in D  \tag{3.17}\\
u(0, \omega) & =0 \\
\frac{\partial}{\partial \xi} u(1, \omega) & =0 .
\end{align*}\right.
$$

Notice that here, performing the change of variable on the variational formulation (3.14) of the problem or directly on the strong formulation (3.13) yields the same result, which is not the case in general. This is due to the fact that $s$ does not depend on the physical variable plus the fact that we are considering the pointwise (in $\omega$ ) weak formulation.

From now on, we assume that the random length of interval $s(\omega)$ has the form

$$
s(\omega)=s_{0}+\varepsilon Y(\omega),
$$

where $Y$ is a random variable with zero mean, unit variance and bounded image $\Gamma$. Moreover, we assume that $Y$ is such that $s(\omega)$ is bounded almost surely from below and above by respectively $s_{\min }$ and $s_{\max }$. More precisely, we assume that

$$
\begin{equation*}
\exists 0<s_{\min } \leq s_{\max }<\infty: \quad P\left(\omega \in \Omega: s_{\min } \leq s(\omega) \leq s_{\max }\right)=1 . \tag{3.18}
\end{equation*}
$$

Due to the Doob-Dynkin lemma, the solution $u$ of (3.17) depends on the same random variable as $s$, i.e. $u(\xi, \omega)=u(\xi, Y(\omega))$. Let $\rho: \Gamma \rightarrow \mathbb{R}^{+}$denotes the density function of $Y$. The solution of problem (3.17) can then be sought either in the probability space $(\Omega, \mathscr{F}, P)$ or in its image space $(\Gamma, B(\Gamma), \rho(y) d y)$. The stochastic problem (3.17) can indeed be written in the following
deterministic parametric form:
find $u: D \times \Gamma \rightarrow \mathbb{R}$ such that $\rho$-a.e. in $\Gamma$ we have

$$
\left\{\begin{align*}
-\frac{a}{s(y)^{2}} \frac{\partial^{2}}{\partial \xi^{2}} u(\xi, y)+\frac{b}{s(y)} u(\xi, y) \frac{\partial}{\partial \xi} u(\xi, y) & =f(\xi, y) \quad \xi \in D  \tag{3.19}\\
u(0, y) & =0 \\
\frac{\partial}{\partial \xi} u(1, y) & =0 .
\end{align*}\right.
$$

From now on, we will drop the dependence of the functions on either $\xi, \omega$ or $y$ when no confusion is possible. Furthermore, we will write $u^{\prime}$ for $\frac{\partial}{\partial \xi}$. Since $s$ is expanded as sum of coefficients, it is more convenient to have all its occurrences in the numerator rather than having division by $s$. Therefore, we will consider the following weak form of problem (3.19):
find $u(y) \in V$ such that

$$
\begin{equation*}
\int_{0}^{1} a u^{\prime} v^{\prime} d \xi+\int_{0}^{1} b s(y) u u^{\prime} v d \xi=\int_{0}^{1} s^{2}(y) f v d \xi \quad \forall v \in V, \rho \text {-a.e. in } \Gamma \text {. } \tag{3.20}
\end{equation*}
$$

Before giving an a posteriori error estimation for the problem (3.17), and thus for the problem (3.13), we briefly give a condition on the given data that ensures the well-posedness of the problem. Recall that $f=\tilde{f} \circ g_{\omega}^{-1}$, i.e. $f(\xi, \omega)=\tilde{f}(s(\omega) \xi)$. Thanks to the uniform bounds on $s$, we have in particular $s^{k} f \in L_{P}^{2}\left(\Omega ; L^{2}(D)\right)$ for any $k$. Notice that it can be shown using only the lower bound $s_{\min }$ or the upper bound $s_{\max }$ depending on the sign of $k$. For instance, we have for the right-hand side of (3.20)

$$
\left\|s^{2}(\omega) f(\omega)\right\|_{L^{2}(D)}=s^{\frac{3}{2}}(\omega)\|\tilde{f}\|_{L^{2}\left(D_{\omega}\right)} \leq s_{\max }^{\frac{3}{2}}\|\tilde{f}\|_{L^{2}(\hat{D})}<\infty \quad \text { a.s. in } \Omega
$$

More generally, we can easily show that the assumption (3.18) ensures that the spaces $L^{2}\left(D_{\omega}\right)$ and $L^{2}(D)$, respectively $\tilde{V}_{\omega}$ and $V$, are isomorphic. This is precisely stated in the following proposition.

Proposition 3.1.5. Under assumption (3.18), for any $\tilde{f} \in L^{2}\left(D_{\omega}\right)$ and any $\tilde{v} \in \tilde{V}_{\omega}$ we have a.s. in $\Omega$

$$
\sqrt{s_{\min }}\|f\|_{L^{2}(D)} \leq\|\tilde{f}\|_{L^{2}\left(D_{\omega}\right)} \leq \sqrt{s_{\max }}\|f\|_{L^{2}(D)}
$$

and

$$
\frac{1}{\sqrt{s_{\max }}}\left\|\frac{\partial v}{\partial \xi}\right\|_{L^{2}(D)} \leq\left\|\frac{\partial \tilde{v}}{\partial x}\right\|_{L^{2}\left(D_{\omega}\right)} \leq \frac{1}{\sqrt{s_{\min }}}\left\|\frac{\partial v}{\partial \xi}\right\|_{L^{2}(D)}
$$

with $f=\tilde{f} \circ g_{\omega}^{-1}$ and $v=\tilde{v} \circ g_{\omega}^{-1}$. The same relations hold for any $f \in L^{2}(D)$ and any $v \in V$ with $\tilde{f}=f \circ g_{\omega}$ and $\tilde{v}=\nu \circ g_{\omega}$.

Similarly to the deterministic problem (3.1), we restrict ourselves to the solutions which lie in $\mathscr{M}$ defined by

$$
\begin{equation*}
\mathscr{M}:=\left\{v \in L_{P}^{2}(\Omega ; V):\left\|v(\omega)^{\prime}\right\|_{L^{2}(0,1)} \leq r_{\omega} \text { a.s. in } \Omega\right\} \tag{3.21}
\end{equation*}
$$

with $r_{\omega}=\sqrt{\frac{s(\omega)}{b} C_{F}\|f(\omega)\|_{L^{2}(0,1)}}$, where $C_{F}$ is the Friedrich-Poincaré constant on the reference interval $D$ given in (3.2). Since

$$
s(\omega)\|f(\omega)\|_{L^{2}(0,1)}=\frac{s(\omega)}{\sqrt{s(\omega)}}\|\tilde{f}\|_{L^{2}\left(D_{\omega}\right)} \leq \sqrt{s_{\max }}\|\tilde{f}\|_{L^{2}(\hat{D})}<\infty \quad \text { a.s. in } \Omega
$$

we have $r_{\omega} \in L_{P}^{\infty}(\Omega)$. Therefore, since $L_{P}^{\infty}(\Omega) \subset L_{P}^{2}(\Omega), \mathscr{M}$ is a closed ball in $L_{P}^{2}(\Omega ; V)$ and thus $\mathscr{M}$ is bounded, convex an closed in $L_{P}^{2}(\Omega ; V)$.

The well-posedness of the stochastic problem can thus be proved following a reasoning similar to the one used in the deterministic case.

Proposition 3.1.6. If bs $(\omega) r_{\omega} \leq \frac{a}{2}$ a.s. in $\Omega$, or in other words if $\frac{4 b s_{\max }^{3}}{a^{2}} C_{F}\|f(\omega)\|_{L^{2}(D)} \leq 1$ a.s. in $\Omega$, then there exists a solution $u \in \mathscr{M}$ to problem (3.20). Furthermore, if the inequality is strict, then the solution is unique.

Remark 3.1.7. We can show the well-posedness of the problem under the slightly less restrictive assumption

$$
\begin{equation*}
\frac{4 C_{F} b s_{\max }^{5 / 2}}{a^{2}}\|\tilde{f}\|_{L^{2}\left(D_{\omega}\right)}<1 \quad \text { a.s. } \operatorname{in} \Omega \tag{3.22}
\end{equation*}
$$

setting then $r_{\omega}=\sqrt{\frac{\sqrt{s(\omega)}}{b}} C_{F}\|f(\omega)\|_{L^{2}(0,1)}$ in (3.21). The inequality (3.22) holds true if the input data satisfy the assumption of Proposition 3.1.6 since $\|f(\omega)\|_{L^{2}(D)} \geq s_{\max }^{-\frac{1}{2}}\|\tilde{f}\|_{L^{2}\left(D_{\omega}\right)}$ by Proposition 3.1.5. We refer to Remark 3.2.9 for the same discussion about the small data assumption for the well-posedness of the Navier-Stokes problem and we mention that the assumption of Proposition 3.1.6 and (3.22) are consistent with (3.43) and (3.41), respectively.

We use a perturbation approach and write

$$
u(\xi, Y(\omega))=u_{0}(\xi)+\varepsilon u_{1}(\xi, Y(\omega))+\mathscr{O}\left(\varepsilon^{2}\right)
$$

with $\varepsilon$ a small parameter that controls the amplitude of the variation of $s$. The goal is now to derive an a posteriori error estimate for the approximation $u \approx u_{0, h}$ with $u_{0, h}$ the finite element approximation of $u_{0}$. We assume that $\tilde{f} \in H^{2}(\hat{D})$ which allows us to write $f=f(\xi, Y(\omega))$ as $f=f_{0}+\varepsilon f_{1} Y+\varepsilon^{2} f_{2} Y^{2}$ with

$$
f_{0}(\xi)=\tilde{f}\left(s_{0} \xi\right), \quad f_{1}(\xi)=\frac{\partial \tilde{f}}{\partial x}\left(s_{0} \xi\right) \xi \quad \text { and } \quad f_{2}(\xi, Y(\omega))=\xi^{2} \int_{0}^{1}(1-t) \frac{\partial^{2} \tilde{f}}{\partial x^{2}}\left(s_{0} \xi+\varepsilon Y(\omega) \xi t\right) d t
$$

using a Taylor expansion with integral remainder of $\tilde{f}(s \xi), s=s_{0}+\varepsilon Y$. The deterministic part $u_{0}$ of the solution can be found by solving

$$
\left\{\begin{align*}
-a u_{0}^{\prime \prime}+b s_{0} u_{0} u_{0}^{\prime} & =s_{0}^{2} f_{0} \quad \text { in } D  \tag{3.23}\\
u_{0}(0) & =0 \\
u_{0}^{\prime}(1) & =0
\end{align*}\right.
$$

Remark 3.1.8. Notice that we could also choose to take $\left(0, s_{0}\right)$ as the reference domain, $i . e$. the interval corresponding to the case $\varepsilon=0$, using then the mapping $g_{\omega}(x)=\frac{s_{0} x}{s(\omega)}$ instead of (3.15). In this case, the problem for $u_{0}$ would not contain the coefficient $s_{0}$, contrary to (3.23). We should then be careful when using for instance the Friedrich-Poincaré inequality (3.2) which holds on $\left(0, s_{0}\right)$ up to a factor $s_{0}$.

We use the finite element method to approximate numerically the solution $u_{0}$ of problem (3.23). To this aim, we consider $0=\xi_{0}<\xi_{1}<\ldots<\xi_{N}<\xi_{N+1}=1$ a partition of $D$ and let $h_{i}=\xi_{i+1}-\xi_{i}$ for $i=0, \ldots, N$. Then, we consider $V_{h}$ the finite dimensional space of $V$ constituted of the corresponding continuous, piecewise linear finite element functions that vanish in 0 . We now give an a posteriori estimate of the error between the exact solution $u$ and the finite element approximation $u_{0, h}$ of $u_{0}$ in the $L_{P}^{2}(\Omega ; V)$ norm.

Proposition 3.1.9. If $\frac{2 b s(\omega) r_{\omega}}{a} \leq \frac{1}{2}$ a.s. in $\Omega$, then there exists a constant $C>0$ depending only on $s_{0}, f_{0}, f_{1}$ and $\mathbb{E}\left[Y^{k} f_{2}^{p}\right]$ for $p=0,1,2$ and some $3 \leq k \leq 8$ such that

$$
\begin{equation*}
\mathbb{E}\left[\left\|u^{\prime}-u_{0, h}^{\prime}\right\|_{L^{2}(0,1)}^{2}\right]^{\frac{1}{2}} \leq \frac{2 \sqrt{2}}{a}\left[\eta_{h}^{2}+\eta_{\varepsilon}^{2}\right]^{\frac{1}{2}}+C \varepsilon^{2} \tag{3.24}
\end{equation*}
$$

with

$$
\begin{align*}
\eta_{h}^{2} & :=C_{I}^{2} \sum_{i=0}^{N} h_{i}^{2} \int_{\xi_{i}}^{\xi_{i+1}}\left(s_{0}^{2} f_{0}-b s_{0} u_{0, h} u_{0, h}^{\prime}+a u_{0, h}^{\prime \prime}\right)^{2} d \xi  \tag{3.25}\\
\eta_{\varepsilon}^{2} & :=\varepsilon^{2} C_{F}^{2}\left\|2 s_{0} f_{0}+s_{0}^{2} f_{1}-b u_{0, h} u_{0, h}^{\prime}\right\|_{L^{2}(D)}^{2} \tag{3.26}
\end{align*}
$$

where $C_{I}$ and $C_{F}$ are the constants in (3.12) and (3.2), respectively.
Remark 3.1.10. The factor 2 in (3.24) comes from the assumption $\frac{2 b s(\omega) r_{\omega}}{a} \leq \frac{1}{2}$ on the input data, which is imposed so that the constant does not explode, see also Remark 3.1.4.

Proof. For any $v \in V$ and a.s. in $\Omega$ we can decompose

$$
\begin{aligned}
\int_{0}^{1} a\left(u^{\prime}-u_{0, h}^{\prime}\right) v^{\prime} d \xi= & \underbrace{\int_{0}^{1}\left(s_{0}^{2} f_{0} v-b s_{0} u_{0, h} u_{0, h}^{\prime} v-a u_{0, h}^{\prime} v^{\prime}\right) d \xi}_{A_{1}(v)}+\underbrace{\int_{0}^{1}\left(s^{2} f-s_{0}^{2} f_{0}\right) v d \xi}_{A_{3}(\nu)} \\
& \underbrace{-\int_{0}^{1} b s\left(u u^{\prime}-u_{0, h} u_{0, h}^{\prime}\right) v d}_{A_{2}(\nu)} \underbrace{-\int_{0}^{1} b\left(s-s_{0}\right) u_{0, h} u_{0, h}^{\prime} v d \xi}_{A_{4}(\nu)}
\end{aligned}
$$

and thus

$$
\left\|u^{\prime}-u_{0, h}^{\prime}\right\|_{L^{2}(D)}^{2}=\frac{1}{a}\left[A_{1}\left(u-u_{0, h}\right)+A_{2}\left(u-u_{0, h}\right)+A_{3}\left(u-u_{0, h}\right)+A_{4}\left(u-u_{0, h}\right)\right] .
$$

Let us consider each term separately. First of all, note that the first term $A_{1}$ corresponds to the residual for $u_{0, h}$, the finite element approximation of problem (3.23). Using a standard
procedure, it can be bounded by

$$
A_{1}(\nu) \leq\left(C_{I}^{2} \sum_{i=0}^{N} h_{i}^{2} \int_{\xi_{i}}^{\xi_{i+1}}\left(s_{0}^{2} f_{0}-b s_{0} u_{0, h} u_{0, h}^{\prime}+a u_{0, h}^{\prime \prime}\right)^{2} d \xi\right)^{\frac{1}{2}}\left\|v^{\prime}\right\|_{L^{2}(D)}
$$

with $C_{I}$ the constant in (3.12). Thanks to the Cauchy-Schwarz and Friedrich-Poincaré inequalities, the second and fourth terms, that we keep together for sharpness ${ }^{2}$, can be bounded by

$$
A_{2}(v)+A_{4}(v) \leq C_{F}\left\|s^{2} f-s_{0}^{2} f_{0}-b\left(s-s_{0}\right) u_{0, h} u_{0, h}\right\|_{L^{2}(D)}\left\|v^{\prime}\right\|_{L^{2}(D)}
$$

Finally, we consider the term $A_{3}$ which is due to the nonlinear part of the problem. If we take $v=u-u_{0, h} \in V$ a.s. in $\Omega$, it can be bounded by

$$
A_{3}\left(u-u_{0, h}\right) \leq 2 b s(\omega) r_{\omega}\left\|u^{\prime}-u_{0, h}^{\prime}\right\|_{L^{2}(D)}^{2}
$$

using Hölder's inequality, Sobolev embedded theorem and the fact that $\left\|u^{\prime}\right\|_{L^{2}(0,1)}$ and $\left\|u_{0, h}^{\prime}\right\|_{L^{2}(0,1)}$ are bounded by $r_{\omega}$ a.s. in $\Omega$. Thanks to the assumption that $\frac{2 b s(\omega) r_{\omega}}{a} \leq \frac{1}{2}$ a.s. in $\Omega$, we have

$$
\frac{1}{a} A_{3}\left(u-u_{0, h}\right) \leq \frac{1}{2}\left\|u^{\prime}-u_{0, h}^{\prime}\right\|_{L^{2}(D)}^{2}
$$

Altogether, we obtain

$$
\begin{aligned}
\left\|u^{\prime}-u_{0, h}^{\prime}\right\|_{L^{2}(D)} \leq & \frac{2}{a}\left[\left(C_{I}^{2} \sum_{i=0}^{N} h_{i}^{2} \int_{\xi_{i}}^{\xi_{i+1}}\left(s_{0}^{2} f_{0}-b s_{0} u_{0, h} u_{0, h}^{\prime}+a u_{0, h}^{\prime \prime}\right)^{2} d \xi\right)^{\frac{1}{2}}\right. \\
& \left.+C_{F}\left\|s^{2} f-s_{0}^{2} f_{0}-b\left(s-s_{0}\right) u_{0, h} u_{0, h}^{\prime}\right\|_{L^{2}(D)}\right]
\end{aligned}
$$

which yields

$$
\begin{aligned}
\left\|u^{\prime}-u_{0, h}^{\prime}\right\|_{L^{2}(D)}^{2} \leq & \frac{8}{a^{2}}\left[C_{I}^{2} \sum_{i=0}^{N} h_{i}^{2} \int_{\xi_{i}}^{\xi_{i+1}}\left(s_{0}^{2} f-b s_{0} u_{0, h} u_{0, h}^{\prime}+a u_{0, h}^{\prime \prime}\right)^{2} d \xi\right. \\
& \left.+C_{F}^{2}\left\|s^{2} f-s_{0}^{2} f_{0}-b\left(s-s_{0}\right) u_{0, h} u_{0, h}^{\prime}\right\|_{L^{2}(D)}^{2}\right] .
\end{aligned}
$$

Since $Y$ has zero mean and unit variance, the result follows taking first the expected value and then the square root on both sides of last inequality. Indeed, we have

$$
s^{2} f-s_{0}^{2} f_{0}=\varepsilon\left(2 s_{0} f_{0}+s_{0}^{2} f_{1}\right) Y+\varepsilon^{2}\left(f_{0}+2 s_{0} f_{1}+s_{0}^{2} f_{2}\right) Y^{2}+\varepsilon^{3}\left(f_{1}+2 s_{0} f_{2}\right) Y^{3}+\varepsilon^{4} f_{2} Y^{4}
$$

from which we deduce, recalling that $s-s_{0}=\varepsilon Y$,

$$
\mathbb{E}\left[\left\|s^{2} f-s_{0}^{2} f_{0}-b\left(s-s_{0}\right) u_{0, h} u_{0, h}^{\prime}\right\|_{L^{2}(D)}^{2}\right]=\varepsilon^{2}\left\|2 s_{0} f_{0}+s_{0}^{2} f_{1}-b u_{0, h} u_{0, h}^{\prime}\right\|_{L^{2}(D)}^{2}+C_{2} \varepsilon^{3}
$$

[^4]where $C_{2}$ depends only on $s_{0}, f_{0}, f_{1}, \mathbb{E}\left[Y^{k}\right]$ for $k=3, \ldots, 6, \mathbb{E}\left[Y^{k} f_{2}\right]$ for $k=3, \ldots, 7$ and $\mathbb{E}\left[Y^{k} f_{2}^{2}\right]$ for $k=4, \ldots, 8$.

Notice that we have used the Friedrich-Poincaré inequality to bound the terms $A_{2}$ and $A_{4}$ due to the forcing and nonlinear terms, for which 1 is a uniform bound for $D=(0,1)$. The loss due to the use of this inequality is different from case to case, therefore affecting the efficiency of the estimator $\eta_{\varepsilon}$ when changing the input data.

### 3.1.3 Numerical results

We consider here two numerical examples for the Burgers' equation. We choose $s_{0}=1$ for simplicity. We start with the results for the deterministic case presented in Section 3.1.1.

## Deterministic case

Let $a=b=1$. For the first example, we consider

$$
\begin{equation*}
\tilde{u}(x)=-0.3 \tanh (x)+0.3 \operatorname{sech}(1)^{2} x, \quad x \in(0,1) \tag{3.27}
\end{equation*}
$$

and compute the corresponding right-hand side $\tilde{f}=-a \tilde{u}^{\prime \prime}+b \tilde{u} u$ and for the second example we set the source term to

$$
\begin{equation*}
\tilde{g}(x)=\sin (\pi x) . \tag{3.28}
\end{equation*}
$$

Notice that $\tilde{g}$ does not satisfy the bound $C_{F}\|\tilde{g}\|_{L^{2}(0,1)}<\frac{a^{2}}{4 b}=0.25$ with $C_{F}=1 / \sqrt{2}$ since $C_{F}\|\tilde{g}\|_{L^{2}(0,1)}=0.5$. We give in Table 3.1 the results for these two cases considering various (uniform) partitions of $[0,1]$. Here, error stands for the error $\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L^{2}(0,1)}$, while

$$
\eta=\frac{1}{a}\left(\sum_{i=0}^{N} \eta_{i}^{2}\right)^{\frac{1}{2}} \quad \text { with } \eta_{i} \text { in (3.9) }
$$

and e.i. denotes the ratio between the estimator $\eta$ and the error. The error is computed with the exact solution for the first case (3.27) and with respect to the reference solution obtained with $h_{r e f}=2^{-12}$ for the second case (3.28).

By looking at the effectivity index for both cases, we see that for $h$ small enough, we recover the value $3.46 \approx 2 \sqrt{3}$ obtained in the one-dimensional numerical examples of the previous chapters, see also Appendix 1.C. The slight increase of e.i. for small value of $h$ in the second case (3.28) is due to the fact that the error is computed with respect to a reference solution.

## Random case

We consider now the case of random interval $D_{\omega}=(0, s(\omega))$ with $s(\omega)=s_{0}+\varepsilon Y(\omega)=1+\varepsilon Y(\omega)$, where $Y$ is a uniform random variable in $[-\sqrt{3}, \sqrt{3}]$. Considering $\tilde{f}$ and $\tilde{g}$ defined above

|  | $f$ |  |  | $\tilde{g}$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $h$ | error | $\eta$ | e.i. | error | $\eta$ | e.i. |
| $1 / 4$ | $1.3440 \mathrm{e}-2$ | $4.6302 \mathrm{e}-2$ | 3.4452 | $5.0731 \mathrm{e}-2$ | $1.6499 \mathrm{e}-1$ | 3.2523 |
| $1 / 8$ | $6.7111 \mathrm{e}-3$ | $2.3217 \mathrm{e}-2$ | 3.4594 | $2.4118 \mathrm{e}-2$ | $8.2193 \mathrm{e}-2$ | 3.4080 |
| $1 / 16$ | $3.3545 \mathrm{e}-3$ | $1.1616 \mathrm{e}-2$ | 3.4629 | $1.1901 \mathrm{e}-2$ | $4.1057 \mathrm{e}-2$ | 3.4499 |
| $1 / 32$ | $1.6771 \mathrm{e}-3$ | $5.8092 \mathrm{e}-3$ | 3.4638 | $5.9306 \mathrm{e}-3$ | $2.0524 \mathrm{e}-2$ | 3.4606 |
| $1 / 64$ | $8.3855 \mathrm{e}-4$ | $2.9047 \mathrm{e}-3$ | 3.4640 | $2.9626 \mathrm{e}-3$ | $1.0261 \mathrm{e}-2$ | 3.4636 |
| $1 / 128$ | $4.1927 \mathrm{e}-4$ | $1.4524 \mathrm{e}-3$ | 3.4641 | $1.4804 \mathrm{e}-3$ | $5.1306 \mathrm{e}-3$ | 3.4656 |
| $1 / 256$ | $2.0964 \mathrm{e}-4$ | $7.2620 \mathrm{e}-4$ | 3.4641 | $7.3909 \mathrm{e}-4$ | $2.5653 \mathrm{e}-3$ | 3.4708 |
| $1 / 512$ | $1.0482 \mathrm{e}-4$ | $3.6340 \mathrm{e}-4$ | 3.4641 | $3.6736 \mathrm{e}-4$ | $1.2826 \mathrm{e}-3$ | 3.4915 |
| $1 / 1024$ | $5.2409 \mathrm{e}-5$ | $1.8155 \mathrm{e}-4$ | 3.4641 | $1.7925 \mathrm{e}-4$ | $6.4132 \mathrm{e}-4$ | 3.5777 |

Table 3.1: Error, estimator and effectivity index for the deterministic Burgers' equation with mesh size $2^{-2} \leq h \leq 2^{-10}$.
as (deterministic) forcing terms for the problems on the physical random domain $D_{\omega}$, the corresponding right-hand sides for the problems on the reference interval $(0,1)$ are then given by $f(\xi, \omega)=\tilde{f}(s(\omega) \xi)$ and $g(\xi, \omega)=\tilde{g}(s(\omega) \xi)$, respectively. We give in Figure 3.3 the graph of the function $f$ and the corresponding solution $u$ of problem (3.20) for different values of $s$ and the results for the second case $g$ can be found in Figure 3.4.


Figure 3.3: Function $f$ and corresponding solution $u$ for various values of $s$.

We give then in Table 3.2 the error $\left\|u-u_{0, h}\right\|_{L_{P}^{2}(\Omega ; V)}$, the estimators $\eta_{h}$ and $\eta_{\varepsilon}$ defined in (3.25) and (3.26), respectively, and the effectivity index for the first case $f$. Notice that the error has been computed with the Monte-Carlo method with a sample size $K=1000$ using a reference solution obtained with $h_{\text {ref }}=2^{-12}$. The results for the second case $g$ are provided in Table 3.3.

As anticipated in the theoretical results, the efficiency of the error estimator $\eta_{\varepsilon}$ is sensitive to the input data. Indeed, it is about 1.6 and 4.7 for the cases $f$ and $g$, respectively. One remedy would be to consider an implicit error estimator for $\eta_{\varepsilon}$, proceeding similarly to what is done in Proposition 3.2.16 for the Navier-Stokes equations or in Proposition 2.2.3 for the model problem with random forcing term.


Figure 3.4: Function $g$ and corresponding solution $u$ for various values of $s$.

|  |  | $\varepsilon=0.005$ |  |  | $\varepsilon=0.00125$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $h$ | $\eta_{h}$ | $\eta_{\varepsilon}$ | error | e.i. | $\eta_{\varepsilon}$ | error | e.i. |
| $1 / 4$ | $4.6302 \mathrm{e}-2$ | $1.9365 \mathrm{e}-3$ | $1.3522 \mathrm{e}-2$ | 3.4273 | $4.8413 \mathrm{e}-4$ | $1.3469 \mathrm{e}-2$ | 3.4379 |
| $1 / 8$ | $2.3217 \mathrm{e}-2$ | $1.9439 \mathrm{e}-3$ | $6.8240 \mathrm{e}-3$ | 3.4141 | $4.8598 \mathrm{e}-4$ | $6.7214 \mathrm{e}-3$ | 3.4549 |
| $1 / 16$ | $1.1616 \mathrm{e}-2$ | $1.9458 \mathrm{e}-3$ | $3.5705 \mathrm{e}-3$ | 3.2988 | $4.8644 \mathrm{e}-4$ | $3.3693 \mathrm{e}-3$ | 3.4508 |
| $1 / 32$ | $5.8092 \mathrm{e}-3$ | $1.9462 \mathrm{e}-3$ | $2.0823 \mathrm{e}-3$ | 2.9422 | $4.8656 \mathrm{e}-4$ | $1.7037 \mathrm{e}-3$ | 3.4218 |
| $1 / 64$ | $2.9047 \mathrm{e}-3$ | $1.9464 \mathrm{e}-3$ | $1.4531 \mathrm{e}-3$ | 2.4063 | $4.8659 \mathrm{e}-4$ | $8.9101 \mathrm{e}-4$ | 3.3055 |
| $1 / 128$ | $1.4524 \mathrm{e}-3$ | $1.9464 \mathrm{e}-3$ | $1.2835 \mathrm{e}-3$ | 1.8922 | $4.8660 \mathrm{e}-4$ | $5.1987 \mathrm{e}-4$ | 2.9464 |
| $1 / 256$ | $7.2620 \mathrm{e}-4$ | $1.9464 \mathrm{e}-3$ | $1.2266 \mathrm{e}-3$ | 1.6936 | $4.8660 \mathrm{e}-4$ | $3.6041 \mathrm{e}-4$ | 2.4254 |
| $1 / 512$ | $3.6310 \mathrm{e}-4$ | $1.9464 \mathrm{e}-3$ | $1.1935 \mathrm{e}-3$ | 1.6590 | $4.8660 \mathrm{e}-4$ | $3.1774 \mathrm{e}-4$ | 1.9108 |
| $1 / 1024$ | $1.8155 \mathrm{e}-4$ | $1.9464 \mathrm{e}-3$ | $1.2322 \mathrm{e}-3$ | 1.5865 | $4.8660 \mathrm{e}-4$ | $3.0688 \mathrm{e}-4$ | 1.6924 |

Table 3.2: Error, estimators and effectivity index for the Burgers' equation in random intervals for the first case $f$ with $\varepsilon=0.005$ and 0.00125 .

|  |  | $\varepsilon=0.01$ |  |  | $\varepsilon=0.0025$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $h$ | $\eta_{h}$ | $\eta_{\varepsilon}$ | error | e.i. | $\eta_{\varepsilon}$ | error | e.i. |
| $1 / 4$ | $1.6499 \mathrm{e}-1$ | $1.7575 \mathrm{e}-2$ | $5.0817 \mathrm{e}-2$ | 3.2652 | $4.3939 \mathrm{e}-3$ | $5.0736 \mathrm{e}-2$ | 3.2532 |
| $1 / 8$ | $8.2193 \mathrm{e}-2$ | $1.6843 \mathrm{e}-2$ | $2.4358 \mathrm{e}-2$ | 3.4444 | $4.2108 \mathrm{e}-3$ | $2.4133 \mathrm{e}-2$ | 3.4103 |
| $1 / 16$ | $4.1057 \mathrm{e}-2$ | $1.6664 \mathrm{e}-2$ | $1.2396 \mathrm{e}-2$ | 3.5745 | $4.1659 \mathrm{e}-3$ | $1.1932 \mathrm{e}-2$ | 3.4587 |
| $1 / 32$ | $2.0524 \mathrm{e}-2$ | $1.6619 \mathrm{e}-2$ | $6.9139 \mathrm{e}-3$ | 3.8196 | $4.1548 \mathrm{e}-3$ | $5.9977 \mathrm{e}-3$ | 3.4914 |
| $1 / 64$ | $1.0261 \mathrm{e}-2$ | $1.6608 \mathrm{e}-2$ | $4.6153 \mathrm{e}-3$ | 4.2299 | $4.1520 \mathrm{e}-3$ | $3.0952 \mathrm{e}-3$ | 3.5763 |
| $1 / 128$ | $5.1306 \mathrm{e}-3$ | $1.6605 \mathrm{e}-2$ | $3.8080 \mathrm{e}-3$ | 4.5640 | $4.1513 \mathrm{e}-3$ | $1.7249 \mathrm{e}-3$ | 3.8261 |
| $1 / 256$ | $2.5653 \mathrm{e}-3$ | $1.6604 \mathrm{e}-2$ | $3.6616 \mathrm{e}-3$ | 4.5885 | $4.1511 \mathrm{e}-3$ | $1.1517 \mathrm{e}-3$ | 4.2369 |
| $1 / 512$ | $1.2826 \mathrm{e}-3$ | $1.6604 \mathrm{e}-2$ | $3.5817 \mathrm{e}-3$ | 4.6496 | $4.1510 \mathrm{e}-3$ | $9.6613 \mathrm{e}-4$ | 4.4970 |
| $1 / 1024$ | $6.4132 \mathrm{e}-4$ | $1.6604 \mathrm{e}-2$ | $3.5264 \mathrm{e}-3$ | 4.7121 | $4.1510 \mathrm{e}-3$ | $8.9897 \mathrm{e}-4$ | 4.6724 |

Table 3.3: Error, estimators and effectivity index for the Burgers' equation in random intervals for the second case $g$ with $\varepsilon=0.01$ and 0.0025 .

### 3.2 Steady-state incompressible Navier-Stokes equations in random domains

We consider now the steady-state incompressible Navier-Stokes equations in random domains. We start with the statement of the problem in Section 3.2.1. We introduce in Section 3.2.2 the corresponding problem on a fixed reference domain using a random mapping and show its well-posedness in Section 3.2.3 under the small data assumption and suitable assumptions on the mapping. A specific but rather general form of the random mapping is introduced in Section 3.2.4, namely that it depends linearly on finite (but arbitrary large) number of independent random variables. In Section 3.2.5, which is the core part, an a posteriori error analysis is performed with the derivation of two a posteriori error estimates for the first order approximation. Finally, numerical experiments are presented in Section 3.2.6 and agree with the theoretical results.

### 3.2.1 Problem statement

Let $D_{\omega} \subseteq \hat{D} \subset \mathbb{R}^{d}, d=2,3$, be an open bounded domain with Lipschitz continuous boundary that depends on a random parameter $\omega \in \Omega$, where $\hat{D}$ is a fixed bounded domain that contains $D_{\omega}$ for all $\omega \in \Omega$. Here $(\Omega, \mathscr{F}, P)$ denotes a complete probability space, where $\Omega$ is the set of outcomes, $\mathscr{F} \subset 2^{\Omega}$ is the $\sigma$-algebra of events and $P: \mathscr{F} \rightarrow[0,1]$ is a probability measure. By a slight abuse of notations, we will denote

$$
D_{\omega} \times \Omega:=\left\{(\mathbf{x}, \omega): \mathbf{x} \in D_{\omega}, \omega \in \Omega\right\}
$$

We consider the steady incompressible Navier-Stokes equations in $D_{\omega}$ :
find a velocity $\tilde{\mathbf{u}}: D_{\omega} \times \Omega \rightarrow \mathbb{R}^{d}$ and a pressure $\tilde{p}: D_{\omega} \times \Omega \rightarrow \mathbb{R}$ such that $P$-almost everywhere (a.e.) in $\Omega$, or in other words almost surely (a.s.), the following equations hold

$$
\left\{\begin{array}{rlll}
-v \Delta_{\mathbf{x}} \tilde{\mathbf{u}}+\left(\tilde{\mathbf{u}} \cdot \nabla_{\mathbf{x}}\right) \tilde{\mathbf{u}}+\nabla_{\mathbf{x}} \tilde{p} & = & \tilde{\mathbf{f}} & \mathbf{x} \in D_{\omega}  \tag{3.29}\\
\nabla_{\mathbf{x}} \cdot \tilde{\mathbf{u}} & = & 0 & \mathbf{x} \in D_{\omega} \\
\tilde{\mathbf{u}} & =\mathbf{0} & \mathbf{x} \in \partial D_{\omega}
\end{array}\right.
$$

where $v$ is the kinematic viscosity, $\tilde{\mathbf{f}} \in\left[L^{2}(\hat{D})\right]^{d}$ is the external force field per unit mass that we assume to be deterministic and well-defined for all $\mathbf{x} \in \hat{D}$. Note that $\tilde{p}$ is the pressure divided by the density of the fluid. We consider homogeneous Dirichlet boundary conditions for the sake of simplicity. Should we consider non-homogeneous conditions, a lifting of the boundary conditions could be used which only modifies the right-hand side of the equations. However, the lifting has to satisfy some assumptions for the problem to be well-posed (see [116] for a complete discussion in the deterministic case). In particular, the forcing term would no longer be deterministic. In (3.29), we have used the following notation: if we write $\mathbf{x}=\left(x_{1}, \ldots, x_{d}\right)$ and
$\tilde{\mathbf{u}}=\left(\tilde{u}_{1}, \ldots, \tilde{u}_{d}\right)^{T}$ then for $i, j=1, \ldots, d$

$$
\nabla_{\mathbf{x}} \tilde{p}=\left(\frac{\partial \tilde{p}}{\partial x_{1}}, \ldots, \frac{\partial \tilde{p}}{\partial x_{d}}\right)^{T}, \quad\left(\nabla_{\mathbf{x}} \tilde{\mathbf{u}}\right)_{i j}=\frac{\partial \tilde{u}_{i}}{\partial x_{j}}, \quad \nabla_{\mathbf{x}} \cdot \tilde{\mathbf{u}}=\sum_{i=1}^{d} \frac{\partial \tilde{u}_{i}}{\partial x_{i}}
$$

and

$$
\left(\Delta_{\mathbf{x}} \tilde{\mathbf{u}}\right)_{i}=\left(\nabla_{\mathbf{x}} \cdot \nabla_{\mathbf{x}} \tilde{\mathbf{u}}\right)_{i}=\sum_{j=1}^{d} \frac{\partial}{\partial x_{j}} \frac{\partial \tilde{u}_{i}}{\partial x_{j}}=\Delta_{\mathbf{x}} \tilde{u}_{i}, \quad\left[\left(\tilde{\mathbf{u}} \cdot \nabla_{\mathbf{x}}\right) \tilde{\mathbf{u}}\right]_{i}=\sum_{j=1}^{d} \tilde{u}_{j} \frac{\partial \tilde{u}_{i}}{\partial x_{j}} .
$$

Note that we will use the same notation to denote the norm of a scalar, vector or matrixvalued function, with the natural extension $\|\mathbf{v}\|^{2}=\sum_{i=1}^{d}\left\|v_{i}\right\|^{2}$ (Euclidean norm) and $\|B\|^{2}=$ $\sum_{i, j=1}^{d}\left\|B_{i j}\right\|^{2}$ (Frobenius norm) for any vector $\mathbf{v}=\left(\nu_{1}, \ldots, v_{d}\right) \in \mathbb{R}^{d}$ and any matrix $B=\left(B_{i j}\right)_{i, j=1}^{d} \in$ $\mathbb{R}^{d \times d}$. In order to write the weak formulation of the problem, we need to introduce some functional spaces. For a given Banach space $W$ with norm $\|\cdot\|_{W}$, we define the Bochner space

$$
L_{P}^{2}(\Omega ; W):=\left\{v: \Omega \rightarrow W, v \text { is strongly measurable and }\|v\|_{L_{P}^{2}(\Omega ; W)}<+\infty\right\}
$$

where $\|v\|_{L_{P}^{2}(\Omega ; W)}^{2}:=\int_{\Omega}\|\nu(\omega)\|_{W}^{2} d P(\omega)=\mathbb{E}\left[\|v\|_{W}^{2}\right]$ using the shorthand notation $v(\omega)=v(\cdot, \omega)$ for ease of presentation. Notice that if $W$ is a separable Hilbert space, then $L_{P}^{2}(\Omega ; W)$ is isomorphic [10] to the tensor product space $L_{P}^{2}(\Omega) \otimes W$. Finally, we define $\tilde{V}_{\omega}=\left[H_{0}^{1}\left(D_{\omega}\right)\right]^{d}$ equipped with the gradient norm $\|\cdot\|_{\tilde{V}_{\omega}}:=\left\|\nabla_{\mathbf{x}} \cdot\right\|_{L^{2}\left(D_{\omega}\right)}$ and $\tilde{Q}_{\omega}=L^{2}\left(D_{\omega}\right)$. Note that unless otherwise clearly stated, the Lebesgue measure is used in $D_{\omega}$. The (pointwise in $\omega$ ) weak formulation of problem (3.29) reads:
find $(\tilde{\mathbf{u}}(\omega), \tilde{p}(\omega)) \in \tilde{V}_{\omega} \times \tilde{Q}_{\omega}$ such that

$$
\left\{\begin{align*}
v \int_{D_{\omega}} \nabla_{\mathbf{x}} \tilde{\mathbf{u}}: \nabla_{\mathbf{x}} \tilde{\mathbf{v}} d \mathbf{x}+\int_{D_{\omega}}\left[\left(\tilde{\mathbf{u}} \cdot \nabla_{\mathbf{x}}\right) \tilde{\mathbf{u}}\right] \cdot \tilde{\mathbf{v}} d \mathbf{x} & -\int_{D_{\omega}} \tilde{p} \nabla_{\mathbf{x}} \cdot \tilde{\mathbf{v}} d \mathbf{x}=\int_{D_{\omega}} \tilde{\mathbf{f}} \cdot \tilde{\mathbf{v}} d \mathbf{x}  \tag{3.30}\\
& -\int_{D_{\omega}} \tilde{q} \nabla_{\mathbf{x}} \cdot \tilde{\mathbf{u}} d \mathbf{x}=0
\end{align*}\right.
$$

for all $(\tilde{\mathbf{v}}, \tilde{q}) \in \tilde{V}_{\omega} \times \tilde{Q}_{\omega}$ and a.s. in $\Omega$. Since we impose Dirichlet conditions on the whole boundary, the pressure is only defined up to an additive constant. We come back to this point in the next section (see Remark 3.2.1). Under the assumption of small data, the well-posedness of the problem on the family of random domains $\left(D_{\omega}\right)_{\omega \in \Omega}$ can be proved using two different approaches. The first one would be to consider the Navier-Stokes equations directly on $D_{\omega} \times \Omega$. Another approach, adopted here, consists in mapping the random domain to a reference one, yielding PDEs on a (fixed, deterministic) reference domain with random coefficients.

### 3.2.2 Formulation on a reference domain

Let $D \subset \mathbb{R}^{d}$ be an open bounded reference domain with Lipschitz continuous boundary $\partial D$. We assume that there exists a mapping $\mathbf{x}: D \times \Omega \rightarrow \mathbb{R}^{d}$ that transforms $D$ into $D_{\omega}$ : for each
$\omega \in \Omega$

$$
\begin{aligned}
& \mathbf{x}_{\omega}: \quad D \quad \rightarrow \quad D_{\omega} \\
& \boldsymbol{\xi} \mapsto \mathbf{x}=\mathbf{x}_{\omega}(\boldsymbol{\xi})
\end{aligned}
$$

where the notation $\mathbf{x}_{\omega}(\boldsymbol{\xi})$ stands for $\mathbf{x}(\boldsymbol{\xi}, \omega)$. We assume that for any $\omega \in \Omega, \mathbf{x}_{\omega}$ is invertible and sufficiently regular so that everything that follows makes sense, the precise regularity assumptions on the random mapping $\mathbf{x}$ being given in Section 3.2.3. Let $\boldsymbol{\xi}_{\omega}$ be the inverse of $\mathbf{x}_{\omega}$ defined by

$$
\begin{aligned}
\boldsymbol{\xi}_{\omega}: D_{\omega} & \rightarrow D \\
\mathbf{x} & \mapsto \boldsymbol{\xi}=\boldsymbol{\xi}_{\omega}(\mathbf{x}) .
\end{aligned}
$$

We also introduce the $d \times d$ Jacobian matrices $A^{-1}=A^{-1}(\boldsymbol{\xi}, \omega)$ and $\tilde{A}=\tilde{A}(\mathbf{x}, \omega)$ corresponding respectively to the random transformations $\mathbf{x}_{\omega}$ and $\boldsymbol{\xi}_{\omega}$ and defined by

$$
A^{-1}=\left(A_{i j}^{-1}\right)_{1 \leq i, j \leq d} \quad \text { with } \quad A_{i j}^{-1}:=\frac{\partial\left(\mathbf{x}_{\omega}\right)_{i}}{\partial \xi_{j}}
$$

and

$$
\tilde{A}=\left(\tilde{A}_{i j}\right)_{1 \leq i, j \leq d} \quad \text { with } \quad \tilde{A}_{i j}:=\frac{\partial\left(\boldsymbol{\xi}_{\omega}\right)_{i}}{\partial x_{j}}
$$

We mention that the matrix $A^{-1}$ is often denoted $F$ in the continuum mechanics literature. For any function $\tilde{g}$ defined on $D_{\omega} \times \Omega$, we denote by $g=\tilde{g} \circ \mathbf{x}_{\omega}$ its corresponding function on $D \times \Omega$, i.e. $g(\boldsymbol{\xi}, \omega)=\tilde{g}(\mathbf{x}, \omega)$ with $\mathbf{x}=\mathbf{x}_{\omega}(\boldsymbol{\xi})$. Notice that the matrix $A=\tilde{A} \circ \mathbf{x}_{\omega}$ is the inverse (in the matrix sense) of $A^{-1}$. From the chain rule, the following relations hold true

$$
\nabla_{\mathbf{x}}=\tilde{A}^{T} \nabla_{\boldsymbol{\xi}} \quad \text { and } \quad \nabla_{\mathbf{x}} \tilde{\mathbf{u}}=\left(\nabla_{\xi} \mathbf{u} \circ \boldsymbol{\xi}_{\omega}\right) \tilde{A},
$$

where $\tilde{A}^{T} \nabla_{\xi}$ is a matrix-vector product. For the sake of notation, we will write $\nabla$ instead of $\nabla_{\xi}$ from now on and use the notation

$$
[(B \nabla) p]_{i}=\sum_{j=1}^{d} B_{i j} \frac{\partial p}{\partial \xi_{j}}, \quad(B \nabla) \cdot \mathbf{u}=\sum_{i, j=1}^{d} B_{i j} \frac{\partial u_{i}}{\partial \xi_{j}}=B: \nabla \mathbf{u}
$$

and

$$
[(B \nabla) \mathbf{u}]_{i j}=\sum_{k=1}^{d} B_{j k} \frac{\partial u_{i}}{\partial \xi_{k}}, \quad[(\mathbf{u} \cdot B \nabla) \mathbf{v}]_{i}=\sum_{j, k=1}^{d} u_{j} B_{j k} \frac{\partial \nu_{i}}{\partial \xi_{k}}
$$

for a $d \times d$ matrix $B=\left(B_{i j}\right)_{1 \leq i, j \leq d}$. Note that $(A \nabla) p=A(\nabla p)$. Moreover, let $J_{\mathbf{x}}=\operatorname{det}\left(A^{-1}\right)$ denotes the determinant of the Jacobian matrix $A^{-1}$ associated to $\mathbf{x}_{\omega}$. Finally, we introduce the spaces $V=\left[H_{0}^{1}(D)\right]^{d}$ and $Q=L_{0}^{2}(D)=\left\{q \in L^{2}(D): \int_{D} q d \xi=0\right\}$.

Remark 3.2.1. We choose to fix the constant part of the pressure by imposing zero average on $D$ and not on $D_{\omega}$, the goal being not to estimate this constant when performing the error analysis. Notice that if we fix $\tilde{p}$ with zero average on $D_{\omega}$, then the average of the corresponding pressure $p=\tilde{p} \circ \mathbf{x}_{\omega}$ on $D$ would be small when $\mathbf{x}_{\omega}$ is a small perturbation of the identity mapping. Indeed, we have $\int_{D} p d \boldsymbol{\xi}=\int_{D} p d \boldsymbol{\xi}-\int_{D_{\omega}} \tilde{p} d \mathbf{x}=\int_{D} p\left(1-\left|J_{\mathbf{x}}\right|\right) d \boldsymbol{\xi}$.

We are now able to write the weak formulation of problem (3.29) on the reference domain, using the change of variable $\mathbf{x}=\mathbf{x}_{\omega}(\boldsymbol{\xi})$ :
find $(\mathbf{u}(\omega), p(\omega)) \in V \times Q$ such that

$$
\left\{\begin{align*}
a(\mathbf{u}, \mathbf{v} ; \omega)+c(\mathbf{u}, \mathbf{u}, \mathbf{v} ; \omega)+b(\mathbf{v}, p ; \omega) & =F(\mathbf{v} ; \omega)  \tag{3.31}\\
b(\mathbf{u}, q ; \omega) & =0
\end{align*}\right.
$$

for all $(\mathbf{v}, q) \in V \times Q$ and a.s. in $\Omega$, where

$$
\begin{align*}
a(\mathbf{u}, \mathbf{v} ; \omega) & :=v \int_{D}(\nabla \mathbf{u} A(\omega)):(\nabla \mathbf{v} A(\omega)) J_{\mathbf{x}}(\omega) d \boldsymbol{\xi} \\
b(\mathbf{v}, q ; \omega) & :=-\int_{D} q J_{\mathbf{x}}(\omega)\left(A(\omega)^{T} \nabla\right) \cdot \mathbf{v} d \boldsymbol{\xi}  \tag{3.32}\\
c(\mathbf{u}, \mathbf{v}, \mathbf{w} ; \omega) & :=\int_{D}\left[\left(\mathbf{u} \cdot A(\omega)^{T} \nabla\right) \mathbf{v}\right] \cdot \mathbf{w} J_{\mathbf{x}}(\omega) d \boldsymbol{\xi} \\
F(\mathbf{v} ; \omega) & :=\int_{D} \mathbf{f}(\omega) \cdot \mathbf{v} J_{\mathbf{x}}(\omega) d \boldsymbol{\xi} .
\end{align*}
$$

Using the relations (see Appendix 3.C for proofs)

$$
\begin{equation*}
(\nabla \mathbf{u} A):(\nabla \mathbf{u} A)=\left(\nabla \mathbf{u} A A^{T}\right):(\nabla \mathbf{u}), \quad \nabla \mathbf{u} A=\left(A^{T} \nabla\right) \mathbf{u} \tag{3.33}
\end{equation*}
$$

and

$$
\begin{equation*}
-\int_{D} q J_{\mathbf{x}}\left(A^{T} \nabla\right) \cdot \mathbf{v} d \boldsymbol{\xi}=\int_{D} J_{\mathbf{x}}\left(A^{T} \nabla q\right) \cdot \mathbf{v} d \boldsymbol{\xi} \tag{3.34}
\end{equation*}
$$

the strong form of (3.31) can be written:
find $\mathbf{u}: D \times \Omega \rightarrow \mathbb{R}^{d}$ and $p: D \times \Omega \rightarrow \mathbb{R}$ such that $P$-almost everywhere in $\Omega$ there holds:

$$
\left\{\begin{align*}
-v \nabla \cdot\left[\left(J_{\mathbf{x}} A A^{T} \nabla\right) \mathbf{u}\right]+\left(\mathbf{u} \cdot J_{\mathbf{x}} A^{T} \nabla\right) \mathbf{u}+\left(J_{\mathbf{x}} A^{T} \nabla\right) p & =\mathbf{f} J_{\mathbf{x}} & \boldsymbol{\xi} \in D  \tag{3.35}\\
\left(J_{\mathbf{x}} A^{T} \nabla\right) \cdot \mathbf{u} & =0 & \boldsymbol{\xi} \in D \\
\mathbf{u} & =\mathbf{0} & \boldsymbol{\xi} \in \partial D .
\end{align*}\right.
$$

Notice that similarly to the formulation in [71], the continuity equation can be equivalently written $\nabla \cdot\left(J_{\mathbf{x}} A \mathbf{u}\right)$ thanks to Piola's identity (see Appendix 3.C).

Remark 3.2.2. If homogeneous Neumann boundary conditions $v \frac{\partial \tilde{\mathbf{u}}}{\partial \tilde{\mathbf{n}}}-\tilde{p} \tilde{\mathbf{n}}=\mathbf{0}$ are prescribed for problem (3.29) on a part of the boundary $\partial D_{\omega}$, typically at the outflow part of the boundary, the corresponding boundary conditions for the problem on the reference domain $D$ read $v J_{\mathbf{x}} \nabla \mathbf{u} A A^{T} \mathbf{n}-p J_{\mathbf{x}} A^{T} \mathbf{n}=\mathbf{0}$. However, the problem might no longer be well-posed due to the loss of (uniform) coercivity of $a(\cdot, \cdot ; \omega)+c(\cdot, \cdot, \cdot ; \omega)$ or its counter part on $D_{\omega}$. Indeed, we are not able to control the negative part of the boundary integral. Braack and al. proved in [29] the existence and uniqueness of a solution to the Navier-Stokes equations with small data and homogeneous Neumann conditions on a part of the boundary after introducing what they called a directed-do-nothing condition, adding a (boundary integral) term in the weak
formulation of the problem. From a physical point of view, a force per unit area is prescribed by imposing $v\left(\nabla_{\mathbf{x}} \tilde{\mathbf{u}}+\left(\nabla_{\mathbf{x}} \tilde{\mathbf{u}}\right)^{T}\right) \tilde{\mathbf{n}}-\tilde{p} \tilde{\mathbf{n}}=\tilde{\mathbf{g}}$, corresponding to $v J_{\mathbf{x}}\left(\nabla \mathbf{u} A+(\nabla \mathbf{u} A)^{T}\right) A^{T} \mathbf{n}-p J_{\mathbf{x}} A^{T} \mathbf{n}=\mathbf{g}$ on the reference domain. In such a case, $\Delta_{\mathbf{x}} \tilde{\mathbf{u}}$ in (3.29) should be replaced by $\nabla_{\mathbf{x}} \cdot\left(\nabla_{\mathbf{x}} \tilde{\mathbf{u}}+\left(\nabla_{\mathbf{x}} \tilde{\mathbf{u}}\right)^{T}\right)$.

### 3.2.3 Well-posedness of the problem

The goal is now to show the well-posedness of problem (3.29), under suitable conditions on the family of random mapping $\left(\mathbf{x}_{\omega}\right)_{\omega \in \Omega}$ and restriction on the input data. We will show that there exists a unique solution $(\mathbf{u}, p)$ to problem (3.31), the weak solution of problem (3.29) being then given by $(\tilde{\mathbf{u}}, \tilde{p})=\left(\mathbf{u} \circ \boldsymbol{\xi}_{\omega}, p \circ \boldsymbol{\xi}_{\omega}\right)$.

For any $\omega \in \Omega$, we assume that $\mathbf{x}_{\omega}: D \rightarrow D_{\omega}$, with $D_{\omega}=\mathbf{x}_{\omega}(D)$, is a one-to-one mapping such that $\mathbf{x}_{\omega} \in\left[W^{1, \infty}(D)\right]^{d}, \boldsymbol{\xi}_{\omega} \in\left[W^{1, \infty}\left(D_{\omega}\right)\right]^{d}$ and $D_{\omega}$ is bounded with Lipschitz continuous boundary $\partial D_{\omega}$. Since $\mathbf{x}_{\omega}$ is invertible, the determinant $J_{\mathbf{x}}$ of its Jacobian matrix $A^{-1}$ does not vanish. Without loss of generality, we can assume that $J_{\mathbf{x}}>0$, namely that the mapping is orientation-preserving. Moreover, we make the following assumption [43, 76] on the singular values $\sigma_{i}$ of $A^{-1}$ : there exist two constants $\sigma_{\min }, \sigma_{\max }$ such that for $i=1, \ldots, d$

$$
\begin{equation*}
0<\sigma_{\min } \leq \sigma_{i}\left(A^{-1}(\boldsymbol{\xi}, \omega)\right) \leq \sigma_{\max }<\infty \quad \text { a.e. in } D \text { and a.s. in } \Omega . \tag{3.36}
\end{equation*}
$$

Notice that the singular values of $A$ are then bounded uniformly from below and above by $\sigma_{\text {max }}^{-1}$ and $\sigma_{\text {min }}^{-1}$, respectively. Therefore, the random mapping $\mathbf{x}$ have finite moment of any order and with the above regularity assumption we have $\mathbf{x} \in L_{P}^{\infty}\left(\Omega ;\left[W^{1, \infty}(D)\right]^{d}\right)$. Moreover, the following properties are immediate consequences of assumption (3.36).

Proposition 3.2.3. Under assumption (3.36), we have a.e. in $D$ and a.s. in $\Omega$

- $\sigma_{\text {min }}^{d} \leq \operatorname{det}\left(A^{-1}\right) \leq \sigma_{\text {max }}^{d}$,
- $\sigma_{\text {max }}^{-2} \leq \lambda_{i}\left(A A^{T}\right) \leq \sigma_{\text {min }}^{-2}$ for $i=1, \ldots, d$,
where $\lambda_{i}\left(A A^{T}\right), i=1, \ldots, d$, denote the eigenvalues of $A A^{T}$.

Proof. Since the eigenvalues of $A^{-1} A^{-T}$ (and thus of the so-called (right) Cauchy-Green strain tensor $A^{-T} A^{-1}$ ) are the square of the singular values of $A^{-1}$, the first relation follows directly from (3.36) and the fact that

$$
\operatorname{det}\left(A^{-1}\right)=\sqrt{\operatorname{det}\left(A^{-1} A^{-T}\right)}=\sqrt{\Pi_{i=1}^{d} \lambda_{i}\left(A^{-1} A^{-T}\right)}=\Pi_{i=1}^{d} \sigma_{i}\left(A^{-1}\right)
$$

The second relation is just a consequence of $\lambda_{i}\left(A A^{T}\right)=\sigma_{i}(A)^{2}$.

The following proposition ensures that the spaces $L^{2}\left(D_{w}\right)$ and $L^{2}(D)$, respectively $\left[H_{0}^{1}\left(D_{w}\right)\right]^{d}$ and $\left[H_{0}^{1}(D)\right]^{d}$, are isomorphic.

Proposition 3.2.4. Under assumption (3.36), for any $\tilde{g} \in L^{2}\left(D_{\omega}\right)$ and $\tilde{\mathbf{v}} \in\left[H^{1}\left(D_{\omega}\right)\right]^{d}$ we have a.s. in $\Omega$

$$
\begin{equation*}
\sigma_{\text {min }}^{\frac{d}{2}}\|g\|_{L^{2}(D)} \leq\|\tilde{g}\|_{L^{2}\left(D_{\omega}\right)} \leq \sigma_{\text {max }}^{\frac{d}{2}}\|g\|_{L^{2}(D)} \tag{3.37}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\sigma_{\text {min }}^{\frac{d}{2}}}{\sigma_{\text {max }}}\|\nabla \mathbf{v}\|_{L^{2}(D)} \leq\left\|\nabla_{\mathbf{x}} \tilde{\mathbf{v}}\right\|_{L^{2}\left(D_{\omega}\right)} \leq \frac{\sigma_{\text {max }}^{\frac{d}{2}}}{\sigma_{\text {min }}}\|\nabla \mathbf{v}\|_{L^{2}(D)} \tag{3.38}
\end{equation*}
$$

with $g=\tilde{g} \circ \mathbf{x}_{\omega}$ and $\mathbf{v}=\tilde{\mathbf{v}} \circ \mathbf{x}_{\omega}$. The same relations hold true for any $g \in L^{2}(D)$ and $\mathbf{v} \in\left[H^{1}(D)\right]^{d}$ with $\tilde{g}=g \circ \boldsymbol{\xi}_{\omega}$ and $\tilde{\mathbf{v}}=\mathbf{v} \circ \boldsymbol{\xi}_{\omega}$.

Proof. Let $\tilde{g} \in L^{2}\left(D_{\omega}\right)$ and $\tilde{\mathbf{v}} \in\left[H^{1}\left(D_{\omega}\right)\right]^{d}$. The proof of (3.37) is immediate using the uniform bounds on $\operatorname{det}\left(A^{-1}\right)$ given by Proposition 3.2.3. For (3.38), we use the fact that $\sigma_{\min }^{d} \sigma_{\max }^{-2}$ and $\sigma_{\text {max }}^{d} \sigma_{\text {min }}^{-2}$ are uniform bounds for the eigenvalues (or equivalently singular values) of the symmetric positive definite matrix $\operatorname{det}\left(A^{-1}\right) A A^{T}$ and the relation

$$
\left\|\nabla_{\mathbf{x}} \tilde{\mathbf{u}}\right\|_{L^{2}\left(D_{\omega}\right)}^{2}=\int_{D}(\nabla \mathbf{u} A):(\nabla \mathbf{u} A) \operatorname{det}\left(A^{-1}\right) d \boldsymbol{\xi}=\int_{D} \sum_{i=1}^{d}\left(\operatorname{det}\left(A^{-1}\right) A A^{T} \nabla u_{i}\right) \cdot \nabla u_{i} d \boldsymbol{\xi}
$$

The proof of (3.37) and (3.38) for the case $g \in L^{2}(D)$ and $\mathbf{v} \in\left[H^{1}(D)\right]^{d}$ is similar using the relations $\sigma_{\text {max }}^{-d} \leq \operatorname{det}(A) \leq \sigma_{\text {min }}^{-d}$ and $\sigma_{\text {max }}^{-2} \sigma_{\text {min }}^{2} \leq \lambda_{i}\left(\operatorname{det}(A) A^{-1} A^{-T}\right) \leq \sigma_{\text {min }}^{-d} \sigma_{\text {max }}^{2}$ a.e. in $D$ and a.s. in $\Omega$.

To show the well-posedness of problem (3.31), the forms $a, b$ and $c$ defined in (3.32) have to satisfy (uniformly) some properties, which we verify in the following proposition.

Proposition 3.2.5. For any $\mathbf{u}, \mathbf{v}, \mathbf{w} \in V$ and any $q \in L^{2}(D)$ we have a.s. in $\Omega$

- $a$ is continuous: $|a(\mathbf{u}, \mathbf{v} ; \omega)| \leq v M\|\nabla \mathbf{u}\|_{L^{2}(D)}\|\nabla \mathbf{v}\|_{L^{2}(D)}$ with $M=\sigma_{\min }^{-2} \sigma_{\max }^{d}$,
- $a$ is coercive: $a(\mathbf{v}, \mathbf{v} ; \omega) \geq v \alpha\|\nabla \mathbf{v}\|_{L^{2}(D)}^{2}$ with $\alpha=\sigma_{\text {max }}^{-2} \sigma_{\text {min }}^{d}$,
- $b$ is continuous: $|b(\mathbf{v}, q ; \omega)| \leq \sigma_{\text {max }}^{d} \sigma_{\text {min }}^{-1}\|q\|_{L^{2}(D)}\|\nabla \mathbf{v}\|_{L^{2}(D)}$,
- cis continuous: $|c(\mathbf{u}, \mathbf{v}, \mathbf{w} ; \omega)| \leq \hat{C}\|\nabla \mathbf{u}\|_{L^{2}(D)}\|\nabla \mathbf{v}\|_{L^{2}(D)}\|\nabla \mathbf{w}\|_{L^{2}(D)}$ with $\hat{C}=C_{I}^{2} \sigma_{\text {max }}^{d} \sigma_{\text {min }}^{-1}$,
where $C_{I}=C_{I}(D)$ is the constant in $\|\mathbf{v}\|_{L^{4}(D)} \leq C_{I}\|\nabla \mathbf{v}\|_{L^{2}(D)}$, resulting from Sobolev embedding's theorem and Poincarés inequality on $D$.

Proof. The proof is immediate from Proposition 3.2.3, Hölder's inequality and the Sobolev embedding theorem. The relation (see e.g. [106])

$$
\int_{D}(\nabla \cdot \mathbf{v})(\nabla \cdot \mathbf{v}) d \boldsymbol{\xi}+\int_{D}(\nabla \times \mathbf{v}) \cdot(\nabla \times \mathbf{v}) d \boldsymbol{\xi}=\int_{D} \nabla \mathbf{v}: \nabla \mathbf{v} d \boldsymbol{\xi} \quad \forall \mathbf{v} \in V
$$

where $\nabla \times \mathbf{v}$ denotes the curl of $\mathbf{v}$, is used to prove the continuity of $b$.

Notice that we do not include the parameter $v$ in the constants $\alpha$ and $M$ linked to the coercivity and continuity of $a$, respectively, because we will track its occurrence in the derivation of our $a$ posteriori error estimates, the goal being to minimize the sensitivity of the effectivity index with respect to $v$. We mention that $b$ is also continuous on $\left[H^{1}(D)\right]^{d}$ with the same constant as in Proposition 3.2.5 up to a multiplication by a factor $\sqrt{d}$ and satisfies the so-called
 $D_{\omega}$ is a Lipschitz domain. Moreover, we assume that there exists a constant $\beta>0$ such that the inf-sup condition holds uniformly with respect to $\omega$, i.e.

$$
\begin{equation*}
\inf _{q \in Q} \sup _{\mathbf{v} \in V} \frac{b(\mathbf{v}, q ;(\omega)}{\|q\|_{L^{2}(D)}\|\nabla \mathbf{v}\|_{L^{2}(D)}} \geq \beta \text { a.s. in } \Omega . \tag{3.39}
\end{equation*}
$$

Remark 3.2.6. The inf-sup condition (3.39) can be easily shown under the assumption that the mapping satisfies $\mathbf{x} \in L_{P}^{\infty}\left(\Omega ;\left[W^{2, \infty}(D)\right]^{d}\right)$, proceeding similarly to [71]. Indeed, for any $q \in Q$ there exists $\mathbf{z} \in V$ such that $\nabla \cdot \mathbf{z}=q$ and $\|\nabla \mathbf{z}\|_{L^{2}(D)} \leq C_{1}\|q\|_{L^{2}(D)}$ with a constant $C_{1}$ depending only on the reference domain $D$, see for instance [69]. Setting $\mathbf{v}=-\left(J_{\mathbf{x}} A\right)^{-1} \mathbf{z}$ we get a.s. in $\Omega$

$$
b(\mathbf{v}, q ; \omega)=\|q\|_{L^{2}(D)}^{2} \geq \frac{1}{C_{1}}\|q\|_{L^{2}(D)}\|\nabla \mathbf{z}\|_{L^{2}(D)} \text { and }\|\nabla \mathbf{v}\|_{L^{2}(D)} \leq C_{2}\left\|\left(J_{\mathbf{x}} A\right)^{-1}\right\|_{W^{1, \infty}(D)}\|\nabla \mathbf{z}\|_{L^{2}(D)}
$$

where $C_{2}$ depends only on the Poincaré constant on $D$. From these two inequalities, we deduce that $\frac{b(\mathbf{V}, q ;(\omega)}{\|\mathbf{V}\|_{L^{2}(D)}} \geq \beta\|q\|_{L^{2}(D)}$.s. in $\Omega$ with $\beta^{-1}=C_{1} C_{2}\left\|\left(J_{\mathbf{x}} A\right)^{-1}\right\|_{L_{P}^{\infty}\left(\Omega ;\left[W^{1, \infty}(D)\right]^{d \times d)}\right.}$.

Let us introduce the subspace $\tilde{V}_{\text {div, } \omega} \subset \tilde{V}_{\omega}$ constituted of all (weakly) divergence-free functions of $\tilde{V}_{\omega}$, and its counterpart on $D$ given by

$$
V_{\mathrm{div}, \omega}:=\{\mathbf{v} \in V: b(\mathbf{v}, q ; \omega)=0 \quad \forall q \in Q \text {, a.s. in } \Omega\} .
$$

We can then formulate the (reduced, pointwise in $\omega$ ) weak formulation of problem (3.31):
find $\mathbf{u}(\omega) \in V_{\text {div, } \omega}$ such that

$$
\begin{equation*}
a(\mathbf{u}, \mathbf{v} ; \omega)+c(\mathbf{u}, \mathbf{u}, \mathbf{v} ; \omega)=F(\mathbf{v} ; \omega) \quad \forall \mathbf{v} \in V_{\text {div, }, \omega}, \text { a.s. in } \Omega . \tag{3.40}
\end{equation*}
$$

Proposition 3.2.7. For $\mathbf{u}(\omega) \in V_{\text {div, } \omega}$ solution of (3.40), there exists a unique pressure $p(\omega) \in Q$ so that $(\mathbf{u}, p)$ is a solution of (3.31), a.s in $\Omega$.

Proof. Follows from the inf-sup condition (see [69, p.283]).

Therefore, to show the well-posedness of problem (3.31), and thus of the original problem
(3.30), it only remains to prove that the nonlinear problem (3.40) admits a unique solution. Recalling that $F$ is defined in (3.32) with $\mathbf{f}=\tilde{\mathbf{f}}_{\circ} \mathbf{x}_{\omega}$, the following proposition gives a sufficient condition on the input data so that problem (3.40) is well-posed.

Proposition 3.2.8. If there exists $\theta \in[0,1)$ such that

$$
\begin{equation*}
\frac{C_{P} C_{I}^{2} \sigma_{\max }^{\frac{3 d}{2}+4}}{v^{2} \sigma_{\min }^{2 d+1}}\|\tilde{\mathbf{f}}\|_{L^{2}\left(D_{\omega}\right)} \leq \theta<1 \quad \text { a.s. } \operatorname{in} \Omega, \tag{3.41}
\end{equation*}
$$

where $C_{P}=C_{P}(D)$ denotes the Poincaré constant on $D$, then problem (3.40) has a unique solution. Moreover, its solution satisfies

$$
\begin{equation*}
\|\nabla \mathbf{u}(\omega)\|_{L^{2}(D)} \leq \theta \frac{v \sigma_{\min }^{d+1}}{C_{I}^{2} \sigma_{\max }^{d+2}}=\theta \frac{v \alpha}{\hat{C}} \quad \text { a.s. } i n \Omega \tag{3.42}
\end{equation*}
$$

with $\alpha$ and $\hat{C}$ defined in Proposition 3.2.5.
Remark 3.2.9. Notice that if condition (3.41) holds, then $\frac{\hat{C}}{(v \alpha)^{2}}\|F(\cdot ; \omega)\|_{V_{\text {div, } \omega}^{\prime}}<1$ a.s. in $\Omega$, where the norm on the dual space is defined in the usual way, which is nothing else but the standard small data assumption for uniqueness (see e.g. [60, 69, 116]). Indeed, we have

$$
\frac{\hat{C}}{(v \alpha)^{2}}\|F(\cdot ; \omega)\|_{V_{\mathrm{div}, \omega}^{\prime}}=\frac{\hat{C}}{(v \alpha)^{2}} \sup _{\mathbf{v} \in V_{\mathrm{div}, \omega}} \frac{|F(\mathbf{v} ; \omega)|}{\|\nabla \mathbf{v}\|_{L^{2}(D)}} \leq \frac{C_{P} C_{I}^{2} \sigma_{\max }^{\frac{3 d}{2}+4}}{v^{2} \sigma_{\min }^{2 d+1}}\|\tilde{\mathbf{f}}\|_{L^{2}\left(D_{\omega}\right)} \quad \text { a.s. in } \Omega,
$$

where for the last inequality we used the relation

$$
|F(\mathbf{v} ; \omega)| \leq \sigma_{\text {max }}^{\frac{d}{2}}\left\|\mathbf{f} J_{\mathbf{x}}^{\frac{1}{2}}\right\|_{L^{2}(D)}\|\mathbf{v}\|_{L^{2}(D)} \leq C_{P} \sigma_{\text {max }}^{\frac{d}{2}}\|\tilde{\mathbf{f}}\|_{L^{2}\left(D_{\omega)}\right)}\|\nabla \mathbf{v}\|_{L^{2}(D)} \quad \text { a.s. in } \Omega .
$$

Moreover, instead of (3.41), we could impose that

$$
\begin{equation*}
\frac{C_{P} C_{I}^{2} \sigma_{\max }^{2(d+2)}}{v^{2} \sigma_{\min }^{2 d+1}}\|\mathbf{f}(\omega)\|_{L^{2}(D)} \leq \theta<1 \quad \text { a.s. } \operatorname{in} \Omega \tag{3.43}
\end{equation*}
$$

since $\|\mathbf{f}(\omega)\|_{L^{2}(D)} \geq \sigma_{\max }^{-\frac{d}{2}}\|\tilde{\mathbf{f}}\|_{L^{2}\left(D_{\omega}\right)}$ by Proposition 3.2.4, and thus (3.43) implies (3.41).

The proof of Proposition 3.2.8 follows the same procedure as the one proposed in [109] for deterministic steady Navier-Stokes equations in a given domain and is based on a fixed point argument.

Proof. In this proof, the explicit dependence of the functions with respect to $\omega \in \Omega$ will not necessarily be indicated, unless ambiguity holds. Moreover, with little abuse of notation we define the space

$$
L_{P}^{2}\left(\Omega ; V_{\operatorname{div}, \omega}\right):=\left\{\mathbf{v} \in L_{P}^{2}(\Omega ; V): \mathbf{v}(\omega) \in V_{\operatorname{div}, \omega} \text { a.s. in } \Omega\right\} .
$$

First of all, we can show that

$$
\begin{equation*}
c(\mathbf{u}, \mathbf{v}, \mathbf{v} ; \omega)=0 \quad \forall \mathbf{u} \in V_{\operatorname{div}, \omega}, \forall \mathbf{v} \in V, \text { a.s. in } \Omega \tag{3.44}
\end{equation*}
$$

Indeed, if we write $\tilde{\mathbf{u}}=\mathbf{u} \circ \boldsymbol{\xi}_{\omega}$ and $\tilde{\mathbf{v}}=\mathbf{v} \circ \boldsymbol{\xi}_{\omega}$ then $\tilde{\mathbf{u}} \in \tilde{V}_{\mathrm{div}, \omega}, \tilde{\mathbf{v}} \in \tilde{V}_{\omega}$ and

$$
\begin{aligned}
c(\mathbf{u}, \mathbf{v}, \mathbf{v} ; \omega) & =\int_{D}\left[\left(\mathbf{u} \cdot A^{T} \nabla\right) \mathbf{v}\right] \cdot \mathbf{v} J_{\mathbf{x}} d \boldsymbol{\xi}=\int_{D_{\omega}}\left[\left(\tilde{\mathbf{u}} \cdot \nabla_{\mathbf{x}}\right) \tilde{\mathbf{v}}\right] \cdot \tilde{\mathbf{v}} d \mathbf{x} \\
& =-\frac{1}{2} \int_{D_{\omega}}\left(\nabla_{\mathbf{x}} \cdot \tilde{\mathbf{u}}\right)|\tilde{\mathbf{v}}|^{2} d \mathbf{x}+\frac{1}{2} \int_{\partial D_{\omega}}(\tilde{\mathbf{u}} \cdot \mathbf{n})|\tilde{\mathbf{v}}|^{2} d s=0
\end{aligned}
$$

using the fact that we have imposed homogeneous Dirichlet boundary conditions. Now, for any $\mathbf{u} \in L_{P}^{2}\left(\Omega ; V_{\text {div }}\right)$ we define the (pointwise in $\omega$ ) bilinear form $\mathscr{A}_{\mathbf{u}(\omega)}(\cdot, \cdot ; \omega): V_{\mathrm{div}, \omega} \times V_{\operatorname{div}, \omega} \rightarrow \mathbb{R}$ by

$$
\mathscr{A}_{\mathbf{u}(\omega)}(\mathbf{w}, \mathbf{v} ; \omega):=a(\mathbf{w}, \mathbf{v} ; \omega)+c(\mathbf{u}(\omega), \mathbf{w}, \mathbf{v} ; \omega)
$$

which is uniformly continuous and coercive (on $V$ and thus on $V_{\text {div, } \omega}$ ) thanks to Proposition 3.2.5 and relation (3.44). Since $\left\|\mathbf{f} J_{\mathbf{x}}\right\|_{L^{2}(D)} \leq \sigma_{\max }^{d / 2}\|\tilde{\mathbf{f}}\|_{L^{2}(\hat{D})}<+\infty$ a.s. in $\Omega$, in particular $\mathbf{f} J_{\mathbf{x}} \in$ $L_{P}^{2}\left(\Omega ; L^{2}(D)\right)$ and Lax-Milgram's lemma ensures the existence of a unique solution to the problem:
for every $\omega \in \Omega$, find $\mathbf{w}(\omega) \in V_{\operatorname{div}, \omega}$ such that

$$
\begin{equation*}
\mathscr{A}_{\mathbf{u}(\omega)}(\mathbf{w}, \mathbf{v} ; \omega)=F(\mathbf{v} ; \omega) \quad \forall \mathbf{v} \in V_{\operatorname{div}, \omega}, \text { a.s. in } \Omega . \tag{3.45}
\end{equation*}
$$

Moreover, taking $\mathbf{v}=\mathbf{w}(\omega)$ in (3.45) and using the coercivity of $\mathscr{A}_{\mathbf{u}}(\cdot, \cdot ; \omega)$ we have a.s. in $\Omega$

$$
v \sigma_{\min }^{d} \sigma_{\max }^{-2}\|\nabla \mathbf{w}\|_{L^{2}(D)}^{2} \leq \mathscr{A}_{\mathbf{u}}(\mathbf{w}, \mathbf{w} ; \omega)=F(\mathbf{w} ; \omega) \leq C_{P} \sigma_{\max }^{\frac{d}{2}}\|\tilde{\mathbf{f}}\|_{L^{2}\left(D_{\omega}\right)}\|\nabla \mathbf{w}\|_{L^{2}(D)}
$$

and thus

$$
\begin{equation*}
\|\nabla \mathbf{w}\|_{L^{2}(D)} \leq \frac{C_{P} \sigma_{\max }^{\frac{d}{2}+2}}{v \sigma_{\min }^{d}}\|\tilde{\mathbf{f}}\|_{L^{2}\left(D_{\omega}\right)} \leq \frac{C_{P} \sigma_{\max }^{\frac{d}{2}+2}}{v \sigma_{\min }^{d}}\|\tilde{\mathbf{f}}\|_{L^{2}(\hat{D})}<\infty \tag{3.46}
\end{equation*}
$$

from which we deduce that $\mathbf{w} \in L_{P}^{2}\left(\Omega ; V_{\operatorname{div}, \omega}\right)$. Notice that a fixed point of the application $\Phi: L_{P}^{2}\left(\Omega ; V_{\operatorname{div}, \omega}\right) \rightarrow L_{P}^{2}\left(\Omega ; V_{\operatorname{div}, \omega}\right)$, which maps $\mathbf{u}$ to the unique solution $\mathbf{w}$ of (3.45), is a solution of problem (3.40). Therefore, it only remains to prove that $\Phi$ is a strict contraction. Let $\mathbf{w}=\Phi(\mathbf{u})$ with $\mathbf{u} \in L_{P}^{2}\left(\Omega ; V_{\mathrm{div}, \omega}\right)$. First, using relation (3.46) we directly get that $\Phi\left(L_{P}^{2}\left(\Omega ; V_{\mathrm{div}}\right)\right) \subset \mathscr{M}$, where the ball $\mathscr{M} \subset L_{P}^{2}\left(\Omega ; V_{\mathrm{div}, \omega}\right)$ is defined by

$$
\mathscr{M}:=\left\{\mathbf{v} \in L_{P}^{2}\left(\Omega ; V_{\mathrm{div}, \omega}\right):\|\nabla \mathbf{v}\|_{L^{2}(D)} \leq \frac{C_{P} \sigma_{\max }^{\frac{d}{2}+2}}{v \sigma_{\min }^{d}}\|\tilde{\mathbf{f}}\|_{L^{2}\left(D_{\omega}\right)} \text { a.s. in } \Omega\right\} .
$$

Finally, we show that $\Phi$ is a contraction, i.e. that there exists a constant $0<k<1$ such that

$$
\|\Phi(\mathbf{u})-\Phi(\overline{\mathbf{u}})\|_{L_{P}^{2}(\Omega ; V)} \leq k\|\mathbf{u}-\overline{\mathbf{u}}\|_{L_{P}^{2}(\Omega ; V)} \quad \forall \mathbf{u}, \overline{\mathbf{u}} \in L_{P}^{2}\left(\Omega ; V_{\operatorname{div}, \omega}\right)
$$

Let $\mathbf{w}=\Phi(\mathbf{u})$ and $\overline{\mathbf{w}}=\Phi(\overline{\mathbf{u}})$. Since $\mathbf{w}$ and $\overline{\mathbf{w}}$ satisfy problem (3.45) with $\mathscr{A}_{\mathbf{u}}(\cdot, \cdot ; \omega)$ and $\mathscr{A}_{\overline{\mathbf{u}}}(\cdot, \cdot ; \omega)$, respectively, we have

$$
a(\mathbf{w}-\overline{\mathbf{w}}, \mathbf{v} ; \omega)+c(\mathbf{u}, \mathbf{w}, \mathbf{v} ; \omega)-c(\overline{\mathbf{u}}, \overline{\mathbf{w}}, \mathbf{v} ; \omega)=0 \quad \forall \mathbf{v} \in V_{\operatorname{div}, \omega}, \text { a.s. in } \Omega
$$

from which we deduce

$$
a(\mathbf{w}-\overline{\mathbf{w}}, \mathbf{v} ; \omega)+c(\mathbf{u}-\overline{\mathbf{u}}, \overline{\mathbf{w}}, \mathbf{v} ; \omega)+c(\mathbf{u}, \mathbf{w}-\overline{\mathbf{w}}, \mathbf{v} ; \omega)=0
$$

or in other words

$$
\mathscr{A}_{\mathbf{u}}(\mathbf{w}-\overline{\mathbf{w}}, \mathbf{v} ; \omega)=-c(\mathbf{u}-\overline{\mathbf{u}}, \overline{\mathbf{w}}, \mathbf{v} ; \omega) .
$$

Since $\overline{\mathbf{w}} \in \mathscr{M}$, taking $\mathbf{v}=\mathbf{w}-\overline{\mathbf{w}}$ in the last equation yields a.s. in $\Omega$

$$
\begin{aligned}
v \sigma_{\min }^{d} \sigma_{\max }^{-2}\|\nabla(\mathbf{w}-\overline{\mathbf{w}})\|_{L^{2}(D)}^{2} & \leq \mathscr{A}_{\mathbf{u}}(\mathbf{w}-\overline{\mathbf{w}}, \mathbf{w}-\overline{\mathbf{w}} ; \omega)=-c(\mathbf{u}-\overline{\mathbf{u}}, \overline{\mathbf{w}}, \mathbf{w}-\overline{\mathbf{w}} ; \omega) \\
& \leq C_{I}^{2} \sigma_{\max }^{d} \sigma_{\min }^{-1}\|\nabla(\mathbf{u}-\overline{\mathbf{u}})\|_{L^{2}(D)}\|\nabla \overline{\mathbf{w}}\|_{L^{2}(D)}\|\nabla(\mathbf{w}-\overline{\mathbf{w}})\|_{L^{2}(D)} \\
& \leq \frac{C_{P} C_{I}^{2} \sigma_{\max }^{\frac{33}{2}+2}}{v \sigma_{\min }^{d+1}}\|\tilde{\mathbf{f}}\|_{L^{2}\left(D_{\omega)}\right)}\|\nabla(\mathbf{u}-\overline{\mathbf{u}})\|_{L^{2}(D)}\|\nabla(\mathbf{w}-\overline{\mathbf{w}})\|_{L^{2}(D)} .
\end{aligned}
$$

Therefore

$$
\|\nabla(\mathbf{w}-\overline{\mathbf{w}})\|_{L^{2}(D)} \leq \frac{C_{P} C_{I}^{2} \sigma_{\max }^{\frac{3 d}{2}+4}}{v^{2} \sigma_{\min }^{2 d+1}}\|\tilde{\mathbf{f}}\|_{L^{2}\left(D_{\omega)}\right)}\|\nabla(\mathbf{u}-\overline{\mathbf{u}})\|_{L^{2}(D)} \quad \text { a.s. in } \Omega
$$

which proves that $\Phi$ is a contraction under the assumption that (3.41) holds. By the Banach contraction theorem, we know that there exists a unique fixed point $\mathbf{u}=\Phi(\mathbf{u})$, which is solution of problem (3.40). The fact that any solution of (3.40) is in $\mathscr{M}$ and is a fixed point of $\Phi$ achieves the proof of well-posedness of the problem. Finally, recalling that $\alpha$ and $\hat{C}$ are defined in Proposition 3.2.5, the bound (3.42) is immediate since

$$
\|\nabla \mathbf{u}\|_{L^{2}(D)} \leq \frac{C_{P} \sigma_{\max }^{\frac{d}{2}+2}}{v \sigma_{\min }^{d}}\|\tilde{\mathbf{f}}\|_{L^{2}\left(D_{\omega)}\right.} \leq \theta \frac{v \sigma_{\min }^{d+1}}{C_{I}^{2} \sigma_{\max }^{d+2}}=\theta \frac{v \sigma_{\max }^{-2} \sigma_{\min }^{d}}{C_{I}^{2} \sigma_{\max }^{d} \sigma_{\min }^{-1}}=\theta \frac{v \alpha}{\hat{C}}
$$

where we have used that $\mathbf{u} \in \mathscr{M}$ for the first inequality and relation (3.41) for the second one.

### 3.2.4 Specific form of the random mapping

We assume from now on that the random mapping $\mathbf{x}(\boldsymbol{\xi}, \omega)$ is parametrized by $L$ mutually independent random variables and write $\mathbf{x}(\boldsymbol{\xi}, \omega)=\mathbf{x}\left(\boldsymbol{\xi}, Y_{1}(\omega), \ldots, Y_{L}(\omega)\right)$ with a slight abuse of notation. This assumption with $L$ finite, usually referred to as finite dimensional noise assumption, is necessary to make the problem feasible for numerical simulation. Such approximation of a random field can be achieved by several techniques, for instance using truncated Karhunen-Loève or Fourier expansions. More precisely, we assume that the mapping $\mathbf{x}_{\omega}$ from
$D$ to $D_{\omega}$ writes

$$
\begin{equation*}
\mathbf{x}_{\omega}(\boldsymbol{\xi})=\boldsymbol{\varphi}_{0}(\boldsymbol{\xi})+\varepsilon \sum_{j=1}^{L} \boldsymbol{\varphi}_{j}(\boldsymbol{\xi}) Y_{j}(\omega) \tag{3.47}
\end{equation*}
$$

where the $Y_{j}, j=1, \ldots, L$, are independent random variables with zero mean and unit variance, the deterministic functions $\boldsymbol{\varphi}_{j}: D \rightarrow \mathbb{R}^{d}$ are assumed to be smooth so that $\nabla \boldsymbol{\varphi}_{0} \in$ $\left[W^{1, \infty}(D)\right]^{d \times d}$ and $\nabla \boldsymbol{\varphi}_{j} \in\left[L^{\infty}(D)\right]^{d \times d}$ for $j=1, \ldots, L$, and $\varepsilon \in\left[0, \varepsilon_{\max }\right]$ is a parameter that controls the amount of randomness. We assume that the random variables $Y_{j}, j=1, \ldots, L$, and the functions $\boldsymbol{\varphi}_{j}, j=0,1, \ldots, L$, are independent of $\varepsilon$. Without loss of generality, we can assume that $\boldsymbol{\varphi}_{0}$ is the identity mapping (see [76]), i.e.

$$
\begin{equation*}
\mathbf{x}_{\omega}(\boldsymbol{\xi})=\boldsymbol{\xi}+\varepsilon \sum_{j=1}^{L} \boldsymbol{\varphi}_{j}(\boldsymbol{\xi}) Y_{j}(\omega) . \tag{3.48}
\end{equation*}
$$

The Jacobian matrix $A^{-1}$ associated to $\mathbf{x}_{\omega}$ therefore reads

$$
A^{-1}(\boldsymbol{\xi}, \omega)=I+\varepsilon A_{1}(\boldsymbol{\xi}, \omega) \quad \text { with } \quad A_{1}(\boldsymbol{\xi}, \omega)=\sum_{j=1}^{L} \nabla \boldsymbol{\varphi}_{j}(\boldsymbol{\xi}) Y_{j}(\omega)
$$

where $I$ denotes the $d \times d$ identity matrix and $\nabla \boldsymbol{\varphi}_{j}(\boldsymbol{\xi})$ is the Jacobian matrix of $\boldsymbol{\varphi}_{j}$ for $j=1, \ldots, L$. Finally, we make the following additional assumptions to ensure that (3.36) is satisfied:

$$
\begin{equation*}
Y_{j}(\Omega)=\left[-\gamma_{j}, \gamma_{j}\right]=: \Gamma_{j} \text { with } \gamma_{j}>0, j=1, \ldots, L \tag{3.49}
\end{equation*}
$$

and

$$
\begin{equation*}
\varepsilon_{\max }<\frac{1}{\delta} \text { with } \delta \text { such that } \sum_{j=1}^{L} \gamma_{j}\left\|\nabla \boldsymbol{\varphi}_{j}(\boldsymbol{\xi})\right\|_{2} \leq \delta \quad \text { a.e. in } D, \tag{3.50}
\end{equation*}
$$

where $\|\cdot\|_{2}$ is the spectral norm. It is straightforward to show that under assumptions (3.49) and (3.50), then (3.36) is fullfield for any $\varepsilon \in\left[0, \varepsilon_{\max }\right]$ with $\sigma_{\min }=1-\varepsilon_{\max } \delta$ and $\sigma_{\max }=1+\varepsilon_{\max } \delta$.

Remark 3.2.10. A (truncated) Karhunen-Loève expansion of the random vector field $\mathbf{x}_{\omega}$ (see [76, 87, 88]) yields a characterization of $\mathbf{x}_{\omega}$ that can be recast into the form (3.47). In this case, the functions $\boldsymbol{\varphi}_{j}, j=1, \ldots, L$, write $\boldsymbol{\varphi}_{j}=\sqrt{\lambda_{j}} \boldsymbol{\psi}_{j}$ with $\left\{\lambda_{j}, \boldsymbol{\psi}_{j}\right\}$ the eigenpairs of the (compact, self-adjoint) integral operator associated with the covariance kernel $V: D \times D \rightarrow \mathbb{R}^{d \times d}$ given by

$$
V\left(\boldsymbol{\xi}, \boldsymbol{\xi}^{\prime}\right):=\frac{1}{\varepsilon^{2}} \mathbb{E}\left[\left(\mathbf{x}_{\omega}(\boldsymbol{\xi})-\boldsymbol{\varphi}_{0}(\boldsymbol{\xi})\right)\left(\mathbf{x}_{\omega}\left(\boldsymbol{\xi}^{\prime}\right)-\boldsymbol{\varphi}_{0}\left(\boldsymbol{\xi}^{\prime}\right)\right)^{T}\right] .
$$

We underline that in this work, we do not take into account the error made when the random mapping is approximated via a finite number of random variables. Therefore, we assume here that (3.47) is an exact representation of the random mapping introduced in Section 3.2.2.

Due to the Doob-Dynkin Lemma, the solutions $\mathbf{u}$ and $p$ of (3.35) depend on the same random variables as $\mathbf{x}_{\omega}$. Defining the random vector $\mathbf{Y}=\left(Y_{1}, \ldots, Y_{L}\right)$, we can thus write $\mathbf{u}(\boldsymbol{\xi}, \omega)=$ $\mathbf{u}(\boldsymbol{\xi}, \mathbf{Y}(\omega))$ and $p(\boldsymbol{\xi}, \omega)=p(\boldsymbol{\xi}, \mathbf{Y}(\omega))$. The complete probability space $(\Omega, \mathscr{F}, P)$ can thus be
replaced by $(\Gamma, B(\Gamma), \rho(\mathbf{y}) d \mathbf{y})$, where $\Gamma=\Gamma_{1} \times \ldots \times \Gamma_{L}, B(\Gamma)$ is the Borel $\sigma$-algebra on $\Gamma$ and $\rho(\mathbf{y}) d \mathbf{y}$ is the probability measure of the random vector $\mathbf{Y}$. Notice that since the random variables $Y_{j}, j=1, \ldots, L$, are assumed independent, the joint density function $\rho$ factorizes as $\rho(\mathbf{y})=\Pi_{j=1}^{L} \rho_{j}\left(y_{j}\right)$ for all $\mathbf{y}=\left(y_{1}, \ldots, y_{L}\right) \in \Gamma$. Therefore, for any integrable function $\hat{g}: \Gamma \rightarrow \mathbb{R}$ on $(\Gamma, B(\Gamma), \rho(\mathbf{y}) d \mathbf{y})$, the expectation of the random variable $g=g(\omega)=\hat{g}(\mathbf{Y}(\omega))$ is by definition given by

$$
\mathbb{E}[g]=\int_{\Omega} g(\omega) d P(\omega)=\int_{\Omega} \hat{g}(\mathbf{Y}(\omega)) d P(\omega)=\int_{\Gamma} \hat{g}(\mathbf{y}) \rho(\mathbf{y}) d \mathbf{y}
$$

With a little abuse of notation, we will not distinguish $g$ and $\hat{g}$ in what follows. The problem (3.31) can then be rewritten into the following parametric form:
find $(\mathbf{u}, p) \in L_{\rho}^{2}(\Gamma ; V) \times L_{\rho}^{2}(\Gamma ; Q)$ such that

$$
\left\{\begin{align*}
a(\mathbf{u}(\mathbf{y}), \mathbf{v} ; \mathbf{y})+c(\mathbf{u}(\mathbf{y}), \mathbf{u}(\mathbf{y}), \mathbf{v} ; \mathbf{y})+b(\mathbf{v}, p(\mathbf{y}) ; \mathbf{y}) & =F(\mathbf{v} ; \mathbf{y})  \tag{3.51}\\
b(\mathbf{u}(\mathbf{y}), q ; \mathbf{y}) & =0
\end{align*}\right.
$$

for all $(\mathbf{v}, q) \in V \times Q$ and $\rho$-a.e. in $\Gamma$, where the various forms are defined as in (3.32) with $A(\boldsymbol{\xi}, \omega)$, $A^{-1}(\boldsymbol{\xi}, \omega), J_{\mathbf{x}}(\boldsymbol{\xi}, \omega)$ and $\mathbf{f}(\boldsymbol{\xi}, \omega)$ replaced by $A(\boldsymbol{\xi}, \mathbf{y}), A^{-1}(\boldsymbol{\xi}, \mathbf{y}), J_{\mathbf{x}}(\boldsymbol{\xi}, \mathbf{y})$ and $\mathbf{f}(\boldsymbol{\xi}, \mathbf{y})$, respectively. This problem is well-posed under the so-called small data assumption (3.41) with $\mathbf{f}(\omega)$ replaced by $\mathbf{f}(\mathbf{y})$ and a.s. in $\Omega$ replaced by $\rho$-a.e. in $\Gamma$, the proof being essentially the same as the proof of Proposition 3.2.8. The random weak solution of problem (3.35), i.e. the solution of (3.31), is then given by $(u(\mathbf{Y}(\omega)), p(\mathbf{Y}(\omega)))$ with $(\mathbf{u}, p)$ the parametric solution of (3.51).

Remark 3.2.11. Notice that for any $\mathbf{y} \in \Gamma$, the partial derivative with respect to $y_{j}$ of the solutions $\tilde{\mathbf{u}}=\tilde{\mathbf{u}}(\mathbf{x}, \mathbf{y})$ and $\tilde{p}=\tilde{p}(\mathbf{x}, \mathbf{y})$ of the problem defined on $D_{\mathbf{y}}$ is given for $j=1, \ldots, L$ by

$$
\begin{equation*}
\frac{\partial \tilde{\mathbf{u}}}{\partial y_{j}}=\frac{\partial \mathbf{u}}{\partial y_{j}} \circ \boldsymbol{\xi}_{\mathbf{y}}+\left(\frac{\partial \boldsymbol{\xi}_{\mathbf{y}}}{\partial y_{j}} \cdot \nabla_{\xi}\right) \mathbf{u} \circ \boldsymbol{\xi}_{\mathbf{y}} \quad \text { and } \quad \frac{\partial \tilde{p}}{\partial y_{j}}=\frac{\partial p}{\partial y_{j}} \circ \boldsymbol{\xi}_{\mathbf{y}}+\frac{\partial \boldsymbol{\xi}_{\mathbf{y}}}{\partial y_{j}} \cdot\left(\nabla_{\boldsymbol{\xi}} p \circ \boldsymbol{\xi}_{\mathbf{y}}\right) \tag{3.52}
\end{equation*}
$$

In other words, the (Eulerian) partial derivative with respect to $y_{j}$ of $\tilde{\mathbf{u}}$ (resp. $\tilde{p}$ ) is equal to the material derivative with respect to $y_{j}$ of $\mathbf{u}=\tilde{\mathbf{u}} \circ \mathbf{x}_{\mathbf{y}}\left(\right.$ resp. $\left.p=\tilde{p} \circ \mathbf{x}_{\mathbf{y}}\right)$, transported back to $D_{\mathbf{y}}$. Moreover, we have the relation

$$
\begin{equation*}
\left(\frac{\partial \boldsymbol{\xi}_{\mathbf{y}}}{\partial y_{j}} \cdot \nabla_{\xi}\right) \mathbf{u} \circ \boldsymbol{\xi}_{\mathbf{y}}=-\left(\frac{\partial \mathbf{x}_{\mathbf{y}}}{\partial y_{j}} \circ \boldsymbol{\xi}_{\mathbf{y}} \cdot \nabla_{\mathbf{x}}\right) \tilde{\mathbf{u}} \tag{3.53}
\end{equation*}
$$

and using it in (3.52) we recognize an analogy with the Arbitrary Lagrangian Eulerian (ALE) formulation of PDEs on moving domains [27,56], where the (Eulerian) partial time-derivative is replaced by the partial time-derivative on the ALE frame written in the Eulerian coordinate plus the convective-type term of the right-hand side of (3.53) in which the so-called domain velocity is involved.

### 3.2.5 A posteriori error analysis

To simplify the presentation, we assume from now on that $d=2$ and that $\tilde{\mathbf{f}} \in\left[H^{2}(\hat{D})\right]^{2}$. Since the forcing term on $D$ is given by $\mathbf{f}=\tilde{\mathbf{f}} \circ \mathbf{x}_{\mathbf{Y}}$ and we assumed $\boldsymbol{\varphi}_{0}$ to be the identity mapping, the regularity assumption on $\tilde{\mathbf{f}}$ allows us to write $\mathbf{f}=\mathbf{f}(\boldsymbol{\xi}, \omega)=\mathbf{f}(\boldsymbol{\xi}, \mathbf{Y}(\omega))$ as

$$
\begin{equation*}
\mathbf{f}(\mathbf{Y})=\mathbf{f}_{0}+\varepsilon \mathbf{f}_{1}(\mathbf{Y})+\mathscr{O}\left(\varepsilon^{2}\right) \quad \text { with } \quad \mathbf{f}_{0}:=\tilde{\mathbf{f}}, \quad \mathbf{f}_{1}(\mathbf{Y}):=\sum_{j=1}^{L} \mathbf{F}_{j} Y_{j}, \quad \mathbf{F}_{j}:=\left(\nabla_{\mathbf{x}} \tilde{\mathbf{f}}\right) \boldsymbol{\varphi}_{j} \tag{3.54}
\end{equation*}
$$

The constant in the term of order $\varepsilon^{2}$ in (3.54) depends on the second derivatives of $\tilde{\mathbf{f}}$ and products $\boldsymbol{\varphi}_{i} \boldsymbol{\varphi}_{j}, i, j=1, \ldots, L$. Moreover, since $d=2$ we have

$$
\begin{equation*}
J_{\mathbf{x}}=\operatorname{det}\left(A^{-1}\right)=\operatorname{det}\left(I+\varepsilon A_{1}\right)=1+\varepsilon \operatorname{tr}\left(A_{1}\right)+\varepsilon^{2} \operatorname{det}\left(A_{1}\right) \quad \text { with } \quad \operatorname{det}\left(A_{1}\right) \leq \delta^{2} \tag{3.55}
\end{equation*}
$$

using assumption (3.50) to bound $\operatorname{det}\left(A_{1}\right)$ and

$$
\begin{equation*}
A=I-\varepsilon A_{1}+\sum_{k=2}^{\infty}(-1)^{k} \varepsilon^{k} A_{1}^{k} \quad \text { with } \quad\left\|\sum_{k=2}^{\infty}(-1)^{k} \varepsilon^{k} A_{1}^{k}\right\|_{2} \leq \frac{\varepsilon^{2} \delta^{2}}{1-\varepsilon \delta} \leq \frac{\varepsilon^{2} \delta^{2}}{\sigma_{\min }} \tag{3.56}
\end{equation*}
$$

where we have used a von Neumann series to expand $A=\left(I+\varepsilon A_{1}\right)^{-1}$. We use a perturbation approach expanding the solution $(\mathbf{u}, p)$ on the reference domain $D$ with respect to $\varepsilon$ up to a certain order as

$$
\begin{equation*}
(\mathbf{u}(\boldsymbol{\xi}, \mathbf{Y}(\omega)), p(\boldsymbol{\xi}, \mathbf{Y}(\omega)))=\left(\mathbf{u}_{0}(\boldsymbol{\xi}), p_{0}(\boldsymbol{\xi})\right)+\varepsilon\left(\mathbf{u}_{1}(\boldsymbol{\xi}, \mathbf{Y}(\omega)), p_{1}(\boldsymbol{\xi}, \mathbf{Y}(\omega))\right)+\ldots \tag{3.57}
\end{equation*}
$$

where $\left(\mathbf{u}_{0}, p_{0}\right)$ is the solution of the standard Navier-Stokes equations on $D$, i.e. it solves:
find $\mathbf{u}_{0}: D \rightarrow \mathbb{R}^{d}$ and $p_{0}: D \rightarrow \mathbb{R}$ such that:

$$
\left\{\begin{array}{rlrl}
-v \Delta \mathbf{u}_{0}+\left(\mathbf{u}_{0} \cdot \nabla\right) \mathbf{u}_{0}+\nabla p_{0} & =\mathbf{f}_{0}, & & \boldsymbol{\xi} \in D  \tag{3.58}\\
\nabla \cdot \mathbf{u}_{0} & =0, & & \boldsymbol{\xi} \in D \\
\mathbf{u}_{0} & =\mathbf{0}, & \boldsymbol{\xi} \in \partial D
\end{array}\right.
$$

Writing $\mathbf{u}_{1}=\sum_{j=1}^{L} \mathbf{U}_{j} Y_{j}$ and $p_{1}=\sum_{j=1}^{L} P_{j} Y_{j}$, it can be shown that the couple ( $\mathbf{u}_{1}, p_{1}$ ) is obtained by solving the $L$ (linear) problems:
for $j=1, \ldots, L$, find $\mathbf{U}_{j}: D \rightarrow \mathbb{R}^{d}$ and $P_{j}: D \rightarrow \mathbb{R}$ such that:

$$
\left\{\begin{align*}
-v \Delta \mathbf{U}_{j}+\left(\mathbf{u}_{0} \cdot \nabla\right) \mathbf{U}_{j}+\left(\mathbf{U}_{j} \cdot \nabla\right) \mathbf{u}_{0}+\nabla P_{j} & =g_{j}\left(\mathbf{u}_{0}, p_{0}\right), & & \boldsymbol{\xi} \in D  \tag{3.59}\\
\nabla \cdot \mathbf{U}_{j} & =h_{j}\left(\mathbf{u}_{0}\right), & & \boldsymbol{\xi} \in D \\
\mathbf{U}_{j} & =\mathbf{0}, & & \boldsymbol{\xi} \in \partial D
\end{align*}\right.
$$

where

$$
\begin{aligned}
g_{j}\left(\mathbf{u}_{0}, p_{0}\right) & =\left(\operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) \mathbf{f}_{0}+\mathbf{F}_{j}\right)+v \nabla \cdot\left[\left(\hat{B}_{j} \nabla\right) \mathbf{u}_{0}\right]-\left(\mathbf{u}_{0} \cdot B_{j} \nabla\right) \mathbf{u}_{0}-\left(B_{j} \nabla\right) p_{0}, \\
h_{j}\left(\mathbf{u}_{0}\right) & =-\left(B_{j} \nabla\right) \cdot \mathbf{u}_{0}
\end{aligned}
$$

with

$$
\begin{equation*}
B_{j}:=\operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) I-\nabla \boldsymbol{\varphi}_{j}^{T} \quad \text { and } \quad \hat{B}_{j}:=\operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) I-\left(\nabla \boldsymbol{\varphi}_{j}+\nabla \boldsymbol{\varphi}_{j}^{T}\right) . \tag{3.60}
\end{equation*}
$$

Some details about the derivation of problems (3.58) and (3.59) are given in Appendix 3.A. Here, we approximate the solution of the deterministic problem (3.58) using the finite element method to obtain an approximation ( $\mathbf{u}_{0, h}, p_{0, h}$ ) and we provide an a posteriori error estimate of ( $\mathbf{u}-\mathbf{u}_{0, h}, p-p_{0, h}$ ). For any $h>0$, let $\mathscr{T}_{h}$ be a family of shape regular partitions (see [49]) of $D$ into $d$-simplices $K$ of diameter $h_{K} \leq h$. Moreover, let ( $V_{h}, Q_{h}$ ) with $V_{h} \subset V$ and $Q_{h} \subset Q$ be a pair of inf-sup stable finite element spaces, such as mini-elements $\mathbb{P}_{1 b}-\mathbb{P}_{1}$ (see [5] or [69, p.175] for a proof of stability of these spaces) or Taylor-Hood $\mathbb{P}_{2}-\mathbb{P}_{1}$. We denote by $\left(\mathbf{u}_{0, h}, p_{0, h}\right)$ the FE approximation of the (weak) solution ( $\mathbf{u}_{0}, p_{0}$ ) of problem (3.58). Writing $\mathbf{y}_{0}=\mathbb{E}[\mathbf{Y}]=\mathbf{0}$, it is obtained by solving:
find $\left(\mathbf{u}_{0, h}, p_{0, h}\right) \in V_{h} \times Q_{h}$ such that

$$
\left\{\begin{align*}
a\left(\mathbf{u}_{0, h}, \mathbf{v}_{h} ; \mathbf{y}_{0}\right)+c\left(\mathbf{u}_{0, h} ; \mathbf{u}_{0, h}, \mathbf{v}_{h} ; \mathbf{y}_{0}\right)+b\left(\mathbf{v}_{h}, p_{0, h} ; \mathbf{y}_{0}\right) & =F\left(\mathbf{v}_{h} ; \mathbf{y}_{0}\right)  \tag{3.61}\\
b\left(\mathbf{u}_{0, h}, q_{h} ; \mathbf{y}_{0}\right) & =0
\end{align*}\right.
$$

for all $\left(\mathbf{v}_{h}, q_{h}\right) \in V_{h} \times Q_{h}$. The rest of this section is devoted to an a posteriori error analysis for the error $\left\|\left\|\left(\mathbf{u}-\mathbf{u}_{0, h}, p-p_{0, h}\right)\right\|\right.$, where the norm $\| \cdot\left\|\|\right.$ is defined for any $(\mathbf{v}, q) \in L_{P}^{2}(\Omega ; V) \times$ $L_{P}^{2}(\Omega ; Q)$ by

$$
\|\mathbf{v}, q\|:=\left(\mathbb{E}\left[v\|\nabla \mathbf{v}\|_{L^{2}(D)}^{2}+\frac{1}{v}\|q\|_{L^{2}(D)}^{2}\right)\right)^{\frac{1}{2}} .
$$

Remark 3.2.12. Notice that we obtain the same results if we use the norm $v^{2}\|\nabla \mathbf{v}\|^{2}+\|q\|^{2}$ or $\|\nabla \mathbf{v}\|^{2}+\frac{1}{v^{2}}\|q\|^{2}$ on $V \times Q$. This choice of scaling is guided by the dimension unit of $v, p$ and $\nabla \mathbf{u}$. This is moreover the natural scaling that arises when analysing the a priori estimates on the solution or when performing the a posteriori error analysis (see Appendix 3.B for more details).

As we will see in the following, the error estimate consists of two parts, namely a part due to the finite element approximation (in $h$ ) and another one due to the uncertainty (in $\varepsilon$ ). Let us define for any $\mathbf{y} \in \Gamma$ the residual $R(\cdot ; \mathbf{y}): V \times Q \rightarrow \mathbb{R}$, which depends on $\left(\mathbf{u}_{0, h}, p_{0, h}\right)$, by $R((\mathbf{v}, q) ; \mathbf{y})=R_{1}(\mathbf{v} ; \mathbf{y})+R_{2}(q ; \mathbf{y})$ with

$$
\begin{aligned}
& R_{1}(\mathbf{v} ; \mathbf{y}):=F(\mathbf{v} ; \mathbf{y})-a\left(\mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right)-b\left(\mathbf{v}, p_{0, h} ; \mathbf{y}\right)-c\left(\mathbf{u}_{0, h}, \mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right) \\
& R_{2}(q ; \mathbf{y}):=-b\left(\mathbf{u}_{0, h}, q ; \mathbf{y}\right) .
\end{aligned}
$$

The first step in the residual-based error estimation consists in linking the error to the residual. The norm of the residual is then bounded by a computable quantity (possibly up to a
multiplicative constant).

Proposition 3.2.13. Let $\sigma_{\text {min }}, \sigma_{\text {max }}, \beta$ and $\theta$ be defined in (3.36), (3.39) and (3.41), respectively. If $h$ is small enough, then there exists a constant $C>0$ depending only on $\theta, \sigma_{\text {min }}, \sigma_{\max }$ and $\beta$ such that a.s. in $\Omega$

$$
\begin{equation*}
v\left\|\nabla\left(\mathbf{u}(\mathbf{Y})-\mathbf{u}_{0, h}\right)\right\|_{L^{2}(D)}^{2}+\frac{1}{v}\left\|p(\mathbf{Y})-p_{0, h}\right\|_{L^{2}(D)}^{2} \leq C\left(\frac{1}{v}\left\|R_{1}(\cdot, \mathbf{Y})\right\|_{V^{\prime}}^{2}+v\left\|R_{2}(\cdot, \mathbf{Y})\right\|_{Q^{\prime}}^{2}\right) \tag{3.62}
\end{equation*}
$$

We mention that the closer $\theta$ to 1 , the larger $C$ in Proposition 3.2.13, see relation (3.71). Similarly, the closer $\sigma_{\min }$ to 0 , the larger $C$ will be. The proof of this proposition is inspired by what is done in [2] for the deterministic steady Navier-Stokes equations. In order to simplify the notation, we will write $\|\cdot\|$ instead of $\|\cdot\|_{L^{2}(D)}$ in the sequel.

Proof. In what follows, all equations depending on $\mathbf{y}$ hold $\rho$-a.e. in $\Gamma$, without specifically mentioning it. Moreover, the dependence of the functions with respect to $\mathbf{y} \in \Gamma$ will not necessarily be indicated. Let $\mathbf{e}(\mathbf{y}):=\mathbf{u}(\mathbf{y})-\mathbf{u}_{0, h}$ and $E(\mathbf{y}):=p(\mathbf{y})-p_{0, h}$. Then (3.51) yields

$$
\begin{equation*}
a(\mathbf{e}, \mathbf{v} ; \mathbf{y})+b(\mathbf{v}, E ; \mathbf{y})+b(\mathbf{e}, q ; \mathbf{y})+D\left(\mathbf{u}, \mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right)=R((\mathbf{v}, q) ; \mathbf{y}) \tag{3.63}
\end{equation*}
$$

for all $(\mathbf{v}, q) \in V \times Q$, where

$$
D\left(\mathbf{u}, \mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right):=c(\mathbf{u}, \mathbf{u}, \mathbf{v} ; \mathbf{y})-c\left(\mathbf{u}_{0, h}, \mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right) .
$$

We can show that

$$
\begin{equation*}
D\left(\mathbf{u}, \mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right) \leq\left(2 \theta v \alpha+\hat{C}\left\|\nabla \mathbf{e}_{0}\right\|\right)\|\nabla \mathbf{e}\|\|\nabla \mathbf{v}\| \tag{3.64}
\end{equation*}
$$

and

$$
\begin{equation*}
D\left(\mathbf{u}, \mathbf{u}_{0, h}, \mathbf{u}-\mathbf{u}_{0, h} ; \mathbf{y}\right) \leq\left(\theta v \alpha+\hat{C}\left\|\nabla \mathbf{e}_{0}\right\|\right)\|\nabla \mathbf{e}\|^{2} \tag{3.65}
\end{equation*}
$$

where $\mathbf{e}_{0}:=\mathbf{u}_{0}-\mathbf{u}_{0, h}$ and $M, \alpha$ and $\hat{C}$ are defined in Proposition 3.2.5. Indeed, for any $\mathbf{v} \in V$ we have

$$
\begin{aligned}
D\left(\mathbf{u}, \mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right) & =c\left(\mathbf{u}, \mathbf{u}-\mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right)+c\left(\mathbf{u}-\mathbf{u}_{0, h}, \mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right) \\
& \leq \hat{C}\left(\|\nabla \mathbf{u}\|+\left\|\nabla \mathbf{u}_{0}\right\|+\left\|\nabla \mathbf{e}_{0}\right\|\right)\|\nabla \mathbf{e}\|\|\nabla \mathbf{v}\| \\
& \leq \hat{C}\left(2 \theta \frac{\alpha v}{\hat{C}}+\left\|\nabla \mathbf{e}_{0}\right\|\right)\|\nabla \mathbf{e}\|\|\nabla \mathbf{v}\|
\end{aligned}
$$

thanks to (3.42), which proves relation (3.64). Relation (3.65) is proved analogously using the fact that $c(\mathbf{u}, \mathbf{v}, \mathbf{v} ; \mathbf{y})=0$ for any $\mathbf{v} \in V$. The rest of the proof consists of two steps, first the derivation of a bound on $\|E\|$ and then a bound on $\|\nabla \mathbf{e}\|$.

Using the inf-sup condition (3.39) for $b$, the bound (3.64) on $D$, the continuity of $a$ and the
relation (3.63) with $q=0$, we have

$$
\begin{align*}
\|E\| & \leq \frac{1}{\beta} \sup _{\mathbf{v} \in V} \frac{\left|b\left(\mathbf{v}, p-p_{0, h} ; \mathbf{y}\right)\right|}{\|\nabla \mathbf{v}\|}=\frac{1}{\beta} \sup _{\mathbf{v} \in V} \frac{\left|R_{1}(\mathbf{v} ; \mathbf{y})-a\left(\mathbf{u}-\mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right)-D\left(\mathbf{u}, \mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right)\right|}{\|\nabla \mathbf{v}\|} \\
& \leq \frac{1}{\beta}\left[\left\|R_{1}(\cdot ; \mathbf{y})\right\|_{V^{\prime}}+\left(v M+2 v \alpha+\hat{C}\left\|\nabla \mathbf{e}_{0}\right\|\right)\|\nabla \mathbf{e}\|\right] . \tag{3.66}
\end{align*}
$$

Therefore, using the relation $(a+b)^{2} \leq 2\left(a^{2}+b^{2}\right)$ we obtain

$$
\begin{equation*}
\frac{1}{v}\|E\|^{2} \leq \frac{2}{\beta^{2} v}\left\|R_{1}(\cdot ; \mathbf{y})\right\|_{V^{\prime}}^{2}+\frac{2\left(M+2 \alpha+\frac{\hat{C}}{v}\left\|\nabla \mathbf{e}_{0}\right\|\right)^{2}}{\beta^{2}} v\|\nabla \mathbf{e}\|^{2} \tag{3.67}
\end{equation*}
$$

We now give a bound on the error $\|\nabla \mathbf{e}\|$ for the velocity. Using the inequalities (3.65) and (3.66), the coercivity of the bilinear form $a$, Young's inequality several times and taking $\mathbf{v}=\mathbf{e}$ and $q=-E$ in (3.63), we get

$$
\begin{align*}
v \alpha\|\nabla \mathbf{e}\|^{2} \leq & a(\mathbf{e}, \mathbf{e} ; \mathbf{y})=R_{1}(\mathbf{e} ; \mathbf{y})-R_{2}(E ; \mathbf{y})-D\left(\mathbf{u}, \mathbf{u}_{0, h}, \mathbf{e}\right) \\
\leq & \left\|R_{1}(\cdot ; \mathbf{y})\right\|_{V^{\prime}}\|\nabla \mathbf{e}\|+\left\|R_{2}(; \mathbf{y})\right\|_{Q^{\prime}}\|E\|+\left(\theta v \alpha+\hat{C}\left\|\nabla \mathbf{e}_{0}\right\|\right)\|\nabla \mathbf{e}\|^{2} \\
\leq & \frac{1}{2 \gamma_{1} v}\left\|R_{1}(\cdot ; \mathbf{y})\right\|_{V^{\prime}}^{2}+\frac{v}{2 \beta^{2} \gamma_{2}}\left\|R_{2}(\cdot ; \mathbf{y})\right\|_{Q^{\prime}}^{2}+\frac{1}{\beta}\left\|R_{1}(\cdot ; \mathbf{y})\right\|_{V^{\prime}}\left\|R_{2}(\cdot ; \mathbf{y})\right\|_{Q^{\prime}} \\
& +\left(\frac{\gamma_{1}}{2}+\frac{\gamma_{2}\left(M+2 \alpha+\frac{\hat{C}}{v}\left\|\nabla \mathbf{e}_{0}\right\|\right)^{2}}{2}+\theta \alpha+\frac{\hat{C}}{v}\left\|\nabla \mathbf{e}_{0}\right\|\right) v\|\nabla \mathbf{e}\|^{2} \\
\leq & \frac{c_{1}}{v}\left\|R_{1}(\cdot ; \mathbf{y})\right\|_{V^{\prime}}^{2}+c_{2} v\left\|R_{2}(\cdot ; \mathbf{y})\right\|_{Q^{\prime}}^{2} \\
& +\left(\frac{\gamma_{1}}{2}+\frac{\gamma_{2}\left(M+2 \alpha+\frac{\hat{C}}{v}\left\|\nabla \mathbf{e}_{0}\right\|\right)^{2}}{2}+\theta \alpha+\frac{\hat{C}}{v}\left\|\nabla \mathbf{e}_{0}\right\|\right) v\|\nabla \mathbf{e}\|^{2}, \tag{3.68}
\end{align*}
$$

with

$$
c_{1}=\frac{1}{2 \gamma_{1}}+\frac{1}{2} \quad \text { and } \quad c_{2}=\frac{1}{2 \gamma_{2} \beta^{2}}+\frac{1}{2 \beta^{2}} .
$$

Recalling that $\theta \in\left[0,1\left[\right.\right.$ and using the convergence of $\mathbf{u}_{0, h}$ to $\mathbf{u}_{0}$ as $h$ tends to 0 , we can choose $h, \gamma_{1}$ and $\gamma_{2}$ small enough so that

$$
\begin{equation*}
\frac{\gamma_{1}}{2}+\frac{\gamma_{2}\left(M+2 \alpha+\frac{\hat{C}}{v}\left\|\nabla \mathbf{e}_{0}\right\|\right)^{2}}{2}+\theta \alpha+\frac{\hat{C}}{v}\left\|\nabla \mathbf{e}_{0}\right\| \leq \frac{1+\theta}{2} \alpha . \tag{3.69}
\end{equation*}
$$

For instance, we can choose $h$ small enough so that

$$
\begin{equation*}
\frac{\hat{C}}{v}\left\|\nabla \mathbf{e}_{0}\right\| \leq \frac{1-\theta}{6} \alpha \tag{3.70}
\end{equation*}
$$

and take

$$
\gamma_{1}=\frac{1-\theta}{3} \alpha \quad \text { and } \quad \gamma_{2}=\frac{1-\theta}{3\left(M+2 \alpha+\frac{1-\theta}{6} \alpha\right)^{2}} \alpha
$$

which depends only on $\theta, \sigma_{\min }$ and $\sigma_{\max }$. Therefore, the last term of the right-hand side of inequality (3.68) can be moved to the left and we get

$$
\begin{equation*}
v\|\nabla \mathbf{e}\|^{2} \leq \frac{2}{(1-\theta) \alpha}\left[\frac{c_{1}}{v}\left\|R_{1}(\cdot ; \mathbf{y})\right\|_{V^{\prime}}^{2}+c_{2} v\left\|R_{2}(\cdot ; \mathbf{y})\right\|_{Q^{\prime}}^{2}\right] . \tag{3.71}
\end{equation*}
$$

Using this bound in (3.67) together with (3.70) we get

$$
\frac{1}{v}\|E\|^{2} \leq\left(\frac{2}{\beta^{2}}+\frac{4 c_{1}}{3 \gamma_{2} \beta^{2}}\right) \frac{1}{v}\left\|R_{1}(\cdot ; \mathbf{y})\right\|_{V^{\prime}}^{2}+\frac{4 c_{2}}{3 \gamma_{2} \beta^{2}} v\left\|R_{2}(\cdot ; \mathbf{y})\right\|_{Q^{\prime}}^{2}
$$

Replacing finally $\mathbf{y}$ by $\mathbf{Y}(\omega)$, the combination of last two inequalities permits to conclude the proof since $c_{1}$ and $c_{2}$ depend only on $\beta$ as well as $\gamma_{1}$ and $\gamma_{2}$, which in turn depend only on $\theta$, $\sigma_{\text {min }}$ et $\sigma_{\text {max }}$.

From Proposition 3.2.13, we deduce the following bound on the error in the $\||\cdot|| |$ norm

$$
\begin{equation*}
\left\|\mathbf{u}-\mathbf{u}_{0, h}, p-p_{0, h} \mid\right\| \leq \sqrt{C}\left(\frac{1}{v} \mathbb{E}\left[\left\|R_{1}\right\|_{V^{\prime}}^{2}\right]+v \mathbb{E}\left[\left\|R_{2}\right\|_{Q^{\prime}}^{2}\right]\right)^{\frac{1}{2}} \tag{3.72}
\end{equation*}
$$

by simply taking first the expected value and then the square root on both sides of inequality (3.62). The goal is now to derive a computable (deterministic) error estimator by estimating the residuals that appear in the right-hand side of (3.72). We use a standard procedure to estimate the part due to the space discretization and proceed in two different ways for the part due to the uncertainty, more precisely the truncation in (3.57). The first one is straightforward and does not require the resolution of additional problems. However, it uses the triangle inequality as well as the Poincaré inequality (on the fixed domain $D$ ) to bound the terms due to the external forces and the convection. Even though the Poincaré constant is a uniform bound, the loss when using Poincaré's inequality can be different depending of the problem, affecting the sharpness of the error estimate from case to case. The second procedure consists in computing the dual norm of some functional, and therefore requires the resolution of additional (linear) problems. However, it has the advantage of requiring the use of Cauchy-Schwarz's inequality only and thus does not suffer from the drawback mentioned above.

## First error estimate

Let $[\cdot]_{\mathbf{n}_{e}}$ denotes the jump across an edge $e \in \mathscr{T}_{h}$ in the direction $\mathbf{n}_{e}$ defined by

$$
[\mathbf{g}]_{\mathbf{n}_{e}}(\boldsymbol{\xi}):=\lim _{t \rightarrow 0^{+}}\left[\mathbf{g}\left(\boldsymbol{\xi}+t \mathbf{n}_{e}\right)-\mathbf{g}\left(\boldsymbol{\xi}-t \mathbf{n}_{e}\right)\right]
$$

where $\mathbf{n}_{e}$ is a unit normal vector to $e$ of arbitrary (but fixed) direction for internal edges and the outward unit vector for boundary edges. Since we impose homogeneous Dirichlet conditions at the boundary, we set the jump to zero for boundary edges. We now have all the ingredients
necessary to derive our first error estimate.
Proposition 3.2.14. Let $(\mathbf{u}, p)$ be the (weak) solution of problem (3.35) and let $\left(\mathbf{u}_{0, h}, p_{0, h}\right)$ be the solution of problem (3.61). If the assumptions of Proposition 3.2.13 are satisfied, then there exist positive constants $C_{1}, C_{2}$ and $C_{3}$ independent of $h$ and $\varepsilon$ such that

$$
\begin{equation*}
\left\|\mathbf{u}-\mathbf{u}_{0, h}, p-p_{0, h}\right\| \leq \sqrt{2 C}\left(C_{1} \eta_{h}^{2}+C_{2} \eta_{\varepsilon}^{2}\right)^{\frac{1}{2}}+\sqrt{C} C_{3} \varepsilon^{2} \quad \text { with } \quad \eta_{h}^{2}=\sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2} \text { and } \eta_{\varepsilon}^{2}=\sum_{j=1}^{L} \eta_{j}^{2} \tag{3.73}
\end{equation*}
$$

where $C$ is the constant in Proposition 3.2.13 and

$$
\begin{equation*}
\eta_{K}^{2}:=\frac{1}{v} \eta_{K, 1}^{2}+v \eta_{K, 2}^{2} \quad \text { and } \quad \eta_{j}^{2}:=\frac{1}{v} \eta_{j, 1}^{2}+v \eta_{j, 2}^{2} \tag{3.74}
\end{equation*}
$$

with
$\eta_{K, 1}^{2}:=h_{K}^{2}\left\|\mathbf{f}_{0}+v \Delta \mathbf{u}_{0, h}-\left(\mathbf{u}_{0, h} \cdot \nabla\right) \mathbf{u}_{0, h}-\nabla p_{0, h}\right\|_{L^{2}(K)}^{2}+\sum_{e \subset K} h_{e}\left\|\frac{1}{2}\left[v\left(\nabla \mathbf{u}_{0, h}\right) \mathbf{n}_{e}-p_{0, h} \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2}$
$\eta_{K, 2}^{2}:=\left\|\nabla \cdot \mathbf{u}_{0, h}\right\|_{L^{2}(K)}^{2}$
$\eta_{j, 1}^{2}:=\varepsilon^{2}\left(\left\|\operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) \mathbf{f}_{0}+\mathbf{F}_{j}\right\|^{2}+v^{2}\left\|\left(\hat{B}_{j} \nabla\right) \mathbf{u}_{0, h}\right\|^{2}+\left\|p_{0, h} B_{j}\right\|^{2}+\left\|\left(\mathbf{u}_{0, h} \cdot B_{j} \nabla\right) \mathbf{u}_{0, h}\right\|^{2}\right)$
$\eta_{j, 2}^{2}:=\varepsilon^{2}\left\|\left(B_{j} \nabla\right) \cdot \mathbf{u}_{0, h}\right\|^{2}$,
$B_{j}$ and $\hat{B}_{j}$ being defined in (3.60), $\mathbf{f}_{0}$ and $\mathbf{F}_{j}$ in (3.54). Moreover, $C_{1}$ depends only on the mesh aspect ratio while $C_{2}$ depends only on the Poincaré constant on $D$.

Remark 3.2.15. Notice that if $\varepsilon_{\max } \delta$ is close to 1 , or in other words $\sigma_{\text {min }}$ is close to 0 , then the constant $C_{3}$ in Proposition 3.2.14 might be large, see (3.56). Therefore, in order for the last term of (3.73) to be negligible, we need to assume small perturbations of the domain, for instance by imposing $\varepsilon_{\max } \leq \frac{1}{2 \delta}$.

Proof. Similarly to the proof of Proposition 3.2.13, it is understood that all equations depending on $\mathbf{y}$ hold $\rho$-a.e. in $\Gamma$ unless explicitly stated. Thanks to (3.72), we only need to bound the expectation of $\frac{1}{v}\left\|R_{1}(\cdot ; \mathbf{Y})\right\|_{V^{\prime}}^{2}$ and $v\left\|R_{2}(\cdot ; \mathbf{Y})\right\|_{Q^{\prime}}^{2}$, that is

$$
\int_{\Gamma} \frac{1}{v}\left\|R_{1}(\cdot ; \mathbf{y})\right\|_{V^{\prime}}^{2} \rho(\mathbf{y}) d \mathbf{y} \quad \text { and } \quad \int_{\Gamma} v\left\|R_{2}(\cdot ; \mathbf{y})\right\|_{Q^{\prime}}^{2} \rho(\mathbf{y}) d \mathbf{y}
$$

by computable quantities. We decompose each term $R_{1}$ and $R_{2}$ into two parts which control the FE error and the error due to truncation in the expansion (3.57), respectively. For $\mathbf{y}_{0}=$ $\mathbb{E}[\mathbf{Y}]=\mathbf{0}$ and for all $\mathbf{y} \in \Gamma, \mathbf{v} \in V$ and $q \in Q$ we write

$$
R_{1}(\mathbf{v} ; \mathbf{y})=R_{1}\left(\mathbf{v} ; \mathbf{y}_{0}\right)+\left[R_{1}(\mathbf{v} ; \mathbf{y})-R_{1}\left(\mathbf{v} ; \mathbf{y}_{0}\right)\right]
$$

and

$$
R_{2}(q ; \mathbf{y})=R_{2}\left(q ; \mathbf{y}_{0}\right)+\left[R_{2}(q ; \mathbf{y})-R_{2}\left(q ; \mathbf{y}_{0}\right)\right] .
$$

Using standard procedure (Galerkin orthogonality, Clément interpolation [50]), see for instance [118], and taking the contribution of the constant $v$ into account, the deterministic quantities can be bounded by

$$
\frac{1}{v}\left\|R_{1}\left(\cdot ; \mathbf{y}_{0}\right)\right\|_{V^{\prime}}^{2}+v\left\|R_{2}\left(\cdot ; \mathbf{y}_{0}\right)\right\|_{Q^{\prime}}^{2} \leq C_{1} \sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2}
$$

where $C_{1}$ depends only on the Clément interpolation constant and the regularity of the mesh and the local error estimator $\eta_{K}$ is defined in (3.74). We now bound the terms due to the uncertainty. We have

$$
R_{1}(\mathbf{v} ; \mathbf{y})-R_{1}\left(\mathbf{v} ; \mathbf{y}_{0}\right)=\mathrm{II}_{1}+\mathrm{II}_{2}+\mathrm{II}_{3}+\mathrm{II}_{4} \quad \text { and } \quad R_{2}(q ; \mathbf{y})-R_{2}\left(q ; \mathbf{y}_{0}\right)=\mathrm{II}_{5}
$$

with

$$
\begin{aligned}
& \mathrm{II}_{1}:=F(\mathbf{v} ; \mathbf{y})-F\left(\mathbf{v} ; \mathbf{y}_{0}\right) \leq C_{P}\left\|J_{\mathbf{x}} \mathbf{f}-\mathbf{f}_{0}\right\|\|\nabla \mathbf{v}\| \\
& \mathrm{I}_{2}:=a\left(\mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}_{0}\right)-a\left(\mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right) \leq v\left\|\left[\left(J_{\mathbf{x}} A A^{T}-I\right) \nabla\right] \mathbf{u}_{0, h}\right\|\|\nabla \mathbf{v}\| \\
& \mathrm{I}_{3}:=b\left(\mathbf{v}, p_{0, h} ; \mathbf{y}_{0}\right)-b\left(\mathbf{v}, p_{0, h} ; \mathbf{y}\right) \leq\left\|\left(J_{\mathbf{x}} A^{T}-I\right) p_{0, h}\right\|\|\nabla \mathbf{v}\| \\
& \mathrm{II}_{4}:=c\left(\mathbf{u}_{0, h}, \mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}_{0}\right)-c\left(\mathbf{u}_{0, h}, \mathbf{u}_{0, h}, \mathbf{v} ; \mathbf{y}\right) \leq C_{P}\left\|\left[\mathbf{u}_{0, h} \cdot\left(J_{\mathbf{x}} A^{T}-I\right) \nabla\right] \mathbf{u}_{0, h}\right\|\|\nabla \mathbf{v}\| \\
& \mathrm{I}_{5}:=b\left(\mathbf{u}_{0, h}, q ; \mathbf{y}_{0}\right)-b\left(\mathbf{u}_{0, h}, q ; \mathbf{y}\right) \leq\left\|\left[\left(J_{\mathbf{x}} A^{T}-I\right) \nabla\right] \cdot \mathbf{u}_{0, h}\right\|\|q\| .
\end{aligned}
$$

The bound for each term is straightforward, except the one for the term $\mathrm{II}_{3}$ which can be obtained by writing it in component form, see Appendix 3.C for details. Therefore, we obtain

$$
\frac{1}{v}\left\|R_{1}(\cdot ; \mathbf{y})\right\|_{V^{\prime}}^{2}+v\left\|R_{2}(\cdot ; \mathbf{y})\right\|_{Q^{\prime}}^{2} \leq C_{1} \eta_{h}^{2}+C_{2} \kappa_{\varepsilon}(\mathbf{y})^{2}
$$

where $C_{2}$ is a (deterministic) constant that depends only on $C_{P}$ and

$$
\begin{aligned}
\kappa_{\varepsilon}^{2}:= & \frac{1}{v}\left\|J_{\mathbf{x}} \mathbf{f}-\mathbf{f}_{0}\right\|^{2}+v\left\|\left[\left(J_{\mathbf{x}} A A^{T}-I\right) \nabla\right] \mathbf{u}_{0, h}\right\|^{2}+\frac{1}{v}\left\|\left(J_{\mathbf{x}} A^{T}-I\right) p_{0, h}\right\|^{2} \\
& +\frac{1}{v}\left\|\left[\mathbf{u}_{0, h} \cdot\left(J_{\mathbf{x}} A^{T}-I\right) \nabla\right] \mathbf{u}_{0, h}\right\|^{2}+v\left\|\left[\left(J_{\mathbf{x}} A^{T}-I\right) \nabla\right] \cdot \mathbf{u}_{0, h}\right\|^{2} .
\end{aligned}
$$

Since the independent random variables $\left\{Y_{j}\right\}$ are assumed to be of zero mean and unit variance, we have $\mathbb{E}\left[Y_{j}\right]=0$ and $\mathbb{E}\left[Y_{i} Y_{j}\right]=\delta_{i j}$ for $i, j=1, \ldots, L$ and thus, using Young's inequality and the
relations (3.54), (3.55) and (3.56), among others, we get

$$
\begin{aligned}
\mathbb{E}\left[\left\|J_{\mathbf{x}} \mathbf{f}-\mathbf{f}_{0}\right\|^{2}\right] & =\varepsilon^{2} \sum_{j=1}^{L}\left\|\operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) \mathbf{f}_{0}+\mathbf{F}_{j}\right\|^{2}+\mathscr{O}\left(\varepsilon^{3}\right) \\
\mathbb{E}\left[\left\|\left[\left(J_{\mathbf{x}} A A^{T}-I\right) \nabla\right] \mathbf{u}_{0, h}\right\|^{2}\right] & =\varepsilon^{2} \sum_{j=1}^{L}\left\|\left(\hat{B}_{j} \nabla\right) \mathbf{u}_{0, h}\right\|^{2}+\mathscr{O}\left(\varepsilon^{3}\right) \\
\mathbb{E}\left[\left\|\left(J_{\mathbf{x}} A^{T}-I\right) p_{0, h}\right\|^{2}\right] & =\varepsilon^{2} \sum_{j=1}^{L}\left\|p_{0, h} B_{j}\right\|^{2}+\mathscr{O}\left(\varepsilon^{3}\right) \\
\mathbb{E}\left[\left\|\left[\mathbf{u}_{0, h} \cdot\left(J_{\mathbf{x}} A^{T}-I\right) \nabla\right] \mathbf{u}_{0, h}\right\|^{2}\right] & =\varepsilon^{2} \sum_{j=1}^{L}\left\|\left(\mathbf{u}_{0, h} \cdot B_{j} \nabla\right) \mathbf{u}_{0, h}\right\|^{2}+\mathscr{O}\left(\varepsilon^{3}\right) \\
\mathbb{E}\left[\left\|\left(J_{\mathbf{x}} A^{T}-I\right) \nabla \cdot \mathbf{u}_{0, h}\right\|^{2}\right] & =\varepsilon^{2} \sum_{j=1}^{L}\left\|\left(B_{j} \nabla\right) \cdot \mathbf{u}_{0, h}\right\|^{2}+\mathscr{O}\left(\varepsilon^{3}\right)
\end{aligned}
$$

with $B_{j}$ and $\hat{B}_{j}$ defined in (3.60). Therefore, for some constant $c_{3}>0$ independent of $\varepsilon$ and $h$ we get

$$
\begin{equation*}
\frac{1}{v} \mathbb{E}\left[\left\|R_{1}\right\|_{V^{\prime}}^{2}\right]+v \mathbb{E}\left[\left\|R_{2}\right\|_{Q^{\prime}}^{2}\right] \leq C_{1} \sum_{K \in \mathscr{T}_{h}} \eta_{K}^{2}+C_{2} \sum_{j=1}^{L} \eta_{j}^{2}+c_{3} \varepsilon^{3} \tag{3.76}
\end{equation*}
$$

where $\eta_{j}$ is defined in (3.74). To conclude the proof, it only remains to take the square root on both sides of inequality (3.76). Indeed, using the notation $\eta_{h}$ and $\eta_{\varepsilon}$ introduced in (3.73), we have

$$
\left(\frac{1}{v} \mathbb{E}\left[\left\|R_{1}\right\|_{V^{\prime}}^{2}\right]+v \mathbb{E}\left[\left\|R_{2}\right\|_{Q^{\prime}}^{2}\right]\right)^{\frac{1}{2}} \leq\left(C_{1} \eta_{h}^{2}+C_{2} \eta_{\varepsilon}^{2}+c_{3} \varepsilon^{3}\right)^{\frac{1}{2}} \leq \sqrt{C_{1}} \eta_{h}+\left(C_{2} \eta_{\varepsilon}^{2}+c_{3} \varepsilon^{3}\right)^{\frac{1}{2}}
$$

thanks to the inequality $\sqrt{a^{2}+b^{2}} \leq a+b$ for any $a, b \geq 0$. Moreover, since $\eta_{\varepsilon}=\mathscr{O}(\varepsilon)$ we get for some constant $C_{3}>0$ independent of $\varepsilon$ and $h$

$$
\left(C_{2} \eta_{\varepsilon}^{2}+c_{3} \varepsilon^{3}\right)^{\frac{1}{2}}=\sqrt{C_{2}} \eta_{\varepsilon}\left(1+\frac{c_{3} \varepsilon^{3}}{C_{2} \eta_{\varepsilon}^{2}}\right)^{\frac{1}{2}}=\sqrt{C_{2}} \eta_{\varepsilon}\left(1+\frac{1}{2} \frac{c_{3} \varepsilon^{3}}{C_{2} \eta_{\varepsilon}^{2}}-\frac{1}{8}\left(\frac{c_{3} \varepsilon^{3}}{C_{2} \eta_{\varepsilon}^{2}}\right)^{2}+\ldots\right) \leq \sqrt{C_{2}} \eta_{\varepsilon}+C_{3} \varepsilon^{2}
$$

Finally, using the inequality $a+b \leq \sqrt{2}\left(a^{2}+b^{2}\right)^{\frac{1}{2}}$ we obtain

$$
\left(\frac{1}{v} \mathbb{E}\left[\left\|R_{1}\right\|_{V^{\prime}}^{2}\right]+v \mathbb{E}\left[\left\|R_{2}\right\|_{Q^{\prime}}^{2}\right]\right)^{\frac{1}{2}} \leq \sqrt{C_{1}} \eta_{h}+\sqrt{C_{2}} \eta_{\varepsilon}+C_{3} \varepsilon^{2} \leq \sqrt{2}\left(C_{1} \eta_{h}^{2}+C_{2} \eta_{\varepsilon}^{2}\right)^{\frac{1}{2}}+C_{3} \varepsilon^{2}
$$

which yields (3.73) thanks to (3.72).

## Second error estimate

As mentioned above, the use of the triangle inequality to bound each term linked to $R_{1}$ separately, plus the Poincaré inequality for some of them (namely $\mathrm{II}_{1}$ and $\mathrm{II}_{4}$ ), in the derivation of the error estimate controlling the randomness of the problem can affect the sharpness of the error estimator $\eta_{\varepsilon}$. However, it has the advantage to require the resolution of only one (nonlinear) problem, namely the problem for $\left(\mathbf{u}_{0, h}, p_{0, h}\right)$. We propose in this section a second error estimate for which the use of these inequalities is not required. It is obtained by computing, approximately, the dual norm of the residual $R_{1}(\mathbf{v} ; \mathbf{y})-R_{1}\left(\mathbf{v} ; \mathbf{y}_{0}\right)$. Similarly to the error estimate of Proposition 3.2.14, the terms of higher order are neglected.

Proposition 3.2.16. Under the assumptions of Proposition 3.2.14, there exist constants $C_{1}, C_{3}$ and $C_{4}$ independent of $h$ and $\varepsilon$ and $s \in(0,1]$ such that

$$
\begin{equation*}
\left\|\mathbf{u}-\mathbf{u}_{0, h}, p-p_{0, h}\right\| \| \leq \sqrt{2 C}\left(C_{1} \eta_{h}^{2}+\hat{\eta}_{\varepsilon}^{2}\right)^{\frac{1}{2}}+\sqrt{C}\left(C_{3} \varepsilon^{2}+C_{4} h^{s} \varepsilon\right) \quad \text { with } \quad \hat{\eta}_{\varepsilon}^{2}=\sum_{j=1}^{L} \hat{\eta}_{j}^{2}, \tag{3.77}
\end{equation*}
$$

where $\eta_{h}$ is as in (3.73) and

$$
\hat{\eta}_{j}^{2}:=\frac{1}{v} \hat{\eta}_{j, 1}^{2}+v \eta_{j, 2}^{2}
$$

with $\eta_{j, 2}$ given in (3.75) and $\hat{\eta}_{j, 1}^{2}:=\varepsilon^{2}\left\|\nabla \mathbf{w}_{j, h}\right\|_{L^{2}(D)}^{2}$ for $j=1, \ldots, L$, and $\mathbf{w}_{j, h} \in V_{h}$ is the solution of

$$
\begin{align*}
\int_{D} \nabla \mathbf{w}_{j, h}: \nabla \mathbf{v}_{h} d \boldsymbol{\xi}= & \int_{D}\left(\operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) \mathbf{f}_{0}+\mathbf{F}_{j}\right) \cdot \mathbf{v}_{h} d \boldsymbol{\xi}-v \int_{D}\left(\hat{B}_{j} \nabla\right) \mathbf{u}_{0, h}: \nabla \mathbf{v}_{h} d \boldsymbol{\xi}+\int_{D} p_{0, h}\left(B_{j} \nabla\right) \cdot \mathbf{v}_{h} d \boldsymbol{\xi} \\
& -\int_{D}\left[\left(\mathbf{u}_{0, h} \cdot B_{j} \nabla\right) \mathbf{u}_{0, h}\right] \cdot \mathbf{v}_{h} d \boldsymbol{\xi} \tag{3.78}
\end{align*}
$$

for all $\mathbf{v}_{h} \in V_{h}$. Moreover, the constant $C_{1}$ depends only on the mesh aspect ratio.

Notice that contrary to the error estimate of Proposition 3.2.14, there is no internal constant multiplying $\hat{\eta}_{\varepsilon}$ in (3.77), the constant $C_{2}=C_{2}\left(C_{P}\right)$ appearing in (3.73) being indeed no longer present.

Proof. The proof is similar to the one of Proposition 3.2.14. The only difference is the estimation of the term $r(\mathbf{v} ; \mathbf{y}):=R_{1}(\mathbf{v} ; \mathbf{y})-R_{1}\left(\mathbf{v} ; \mathbf{y}_{0}\right)$ in the $V^{\prime}$ norm. We have $\|r(\cdot ; \mathbf{y})\|_{V^{\prime}}=\|\nabla \mathbf{w}(\mathbf{y})\|_{L^{2}(D)}$, where $\mathbf{w}$ denotes the Riesz representant of $r$, i.e. $\mathbf{w}(\mathbf{y}) \in V$ is such that $\int_{D} \nabla \mathbf{w}(\mathbf{y}): \nabla \mathbf{v}=r(\mathbf{v} ; \mathbf{y})$ for all $\mathbf{v} \in V$ and $\rho$-a.e. in $\Gamma$. If we keep only the terms of order $\varepsilon$ and use the properties of the random variables $Y_{j}, j=1, \ldots, L$, taking the expected value of $\|r(\cdot ; \mathbf{Y})\|_{V^{\prime}}^{2}$, we get

$$
\mathbb{E}\left[\|r\|_{V^{\prime}}^{2}\right] \leq \varepsilon^{2} \sum_{j=1}^{L}\left\|\nabla \mathbf{w}_{j}\right\|_{L^{2}(D)}^{2}+\mathscr{O}\left(\varepsilon^{3}\right)
$$

where $\mathbf{w}_{j}$ is the solution of

$$
\int_{D} \nabla \mathbf{w}_{j}: \nabla \mathbf{v} d \boldsymbol{\xi}=r_{j}\left(\mathbf{v} ; \mathbf{u}_{0, h}, p_{0, h}\right) \quad \forall \mathbf{v} \in V
$$

with

$$
\begin{aligned}
r_{j}\left(\mathbf{v} ; \mathbf{u}_{0, h}, p_{0, h}\right):= & \int_{D}\left(\operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) \mathbf{f}_{0}+\mathbf{F}_{j}\right) \cdot \mathbf{v} d \boldsymbol{\xi}-v \int_{D}\left(\hat{B}_{j} \nabla\right) \mathbf{u}_{0, h}: \nabla \mathbf{v} d \boldsymbol{\xi}+\int_{D} p_{0, h}\left(B_{j} \nabla\right) \cdot \mathbf{v} d \boldsymbol{\xi} \\
& -\int_{D}\left[\left(\mathbf{u}_{0, h} \cdot B_{j} \nabla\right) \mathbf{u}_{0, h}\right] \cdot \mathbf{v} d \boldsymbol{\xi} .
\end{aligned}
$$

Obviously, the solution $\mathbf{w}_{j}$ cannot be computed exactly. However, replacing $\mathbf{w}_{j}$ by its finite element approximation $\mathbf{w}_{j, h} \in V_{h}$ introduces an error of higher order, namely an error of order $\varepsilon h^{s}$ with $s \in(0,1]$. Indeed, introducing for $j=1, \ldots, L$ the solution $\boldsymbol{\psi}_{j} \in V$ of

$$
\int_{D} \nabla \boldsymbol{\psi}_{j}: \nabla \mathbf{v}=r_{j}\left(\mathbf{v} ; u_{0}, p_{0}\right) \quad \mathbf{v} \in V
$$

and its finite element approximation $\boldsymbol{\psi}_{j, h} \in V_{h}$, we have thanks to triangle's inequality

$$
\begin{aligned}
\left\|\nabla \mathbf{w}_{j}\right\|_{L^{2}(D)} & \leq\left\|\nabla\left(\boldsymbol{\psi}_{j}-\mathbf{w}_{j}\right)\right\|_{L^{2}(D)}+\left\|\nabla\left(\boldsymbol{\psi}_{j}-\boldsymbol{\psi}_{j, h}\right)\right\|_{L^{2}(D)}+\left\|\nabla\left(\boldsymbol{\psi}_{j, h}-\mathbf{w}_{j, h}\right)\right\|_{L^{2}(D)}+\left\|\nabla \mathbf{w}_{j, h}\right\|_{L^{2}(D)} \\
& \leq\left\|r_{j}\left(\cdot ; \mathbf{u}_{0}, p_{0}\right)-r_{j}\left(\cdot ; \mathbf{u}_{0, h}, p_{0, h}\right)\right\|_{V^{\prime}}+\left\|\nabla\left(\boldsymbol{\psi}_{j}-\boldsymbol{\psi}_{j, h}\right)\right\|_{L^{2}(D)}+\left\|\nabla \mathbf{w}_{j, h}\right\|_{L^{2}(D)} \\
& \leq C_{4} h^{s}+\left\|\nabla \mathbf{w}_{j, h}\right\|_{L^{2}(D)}
\end{aligned}
$$

where $s \in(0,1]$ depends only on the regularity of $\mathbf{u}_{0}, p_{0}, \boldsymbol{\psi}_{j}, j=1, \ldots, L$, and the domain $D[53,72]$ and $C_{4}$ is independent of $h$ and $\varepsilon$ but depends on the mesh aspect ratio, $\left|\mathbf{u}_{0}\right|_{H^{1+s}(D)}$, $\left|p_{0}\right|_{H^{s}(D)}$ and $\left|\boldsymbol{\psi}_{j}\right|_{H^{1+s}(D)}, j=1, \ldots, L$.

Based on Propositions 3.2.14 and 3.2.16, we can define two computable error estimators $\eta=\left(\eta_{h}^{2}+\eta_{\varepsilon}^{2}\right)^{\frac{1}{2}}$ and $\hat{\eta}=\left(\eta_{h}^{2}+\hat{\eta}_{\varepsilon}^{2}\right)^{\frac{1}{2}}$, where $\eta_{h}$ and $\eta_{\varepsilon}$ are defined in (3.73) and $\hat{\eta}_{\varepsilon}$ is defined in (3.77). From a computational point of view, the computation of $\hat{\eta}$ requires the solution of $L$ additional (linear) problems compared to the cost of getting the error estimator $\eta$. However, the gain of the second error estimator is twofold: it does not use the triangle inequality to bound each term of $r(\mathbf{v} ; \mathbf{y})$ separately and it does not require the use of the Poincaré inequality. The numerical tests of the next section provide an illustration of the theoretical results obtained so far.

### 3.2.6 Numerical results

We present now two numerical examples to test the error estimators derived in the previous section. We consider the problem of a flow past a cylinder and consider two different types of perturbation of the domain, namely a perturbation along the vertical axis of the position of the cylinder and a perturbation of its shape. The true error $\left\|\left|\mathbf{u}-\mathbf{u}_{0, h}, p-p_{0, h}\right|\right\|$ is approximated
with the standard Monte Carlo method using

$$
\|\mid \mathbf{v}, q\| \| \approx\left(\frac{1}{K} \sum_{k=1}^{K}\left\{v\left\|\nabla \mathbf{v}\left(\mathbf{y}_{k}\right)\right\|_{L^{2}(D)}^{2}+\frac{1}{v}\left\|q\left(\mathbf{y}_{k}\right)\right\|_{L^{2}(D)}^{2}\right\}\right)^{\frac{1}{2}}
$$

where $\left\{\mathbf{y}_{k}\right\} \in \Gamma$ are i.i.d. realizations of the random vector $\mathbf{Y}$. We choose a sample size of $K=1000$ in which case the variance of the estimation of the error is at least a factor $2 \cdot 10^{-4}$ smaller than the estimated error in all considered test cases. In what follows, whenever we refer to error it should be understood that the true error has been computed by the Monte Carlo procedure. Finally, the approximate solution $\left(\mathbf{u}_{0, h}, p_{0, h}\right)$ is computed using $\mathbb{P}_{1 b}-\mathbb{P}_{1}$ finite elements and, since the exact solution ( $\mathbf{u}, p$ ) of the problem is not known, we compute a reference solution using $\mathbb{P}_{2}-\mathbb{P}_{1}$ finite elements on the finest mesh considered.

## First example

For this first problem, based on a well-known benchmark problem described in [110], we consider the geometry presented in Figure 3.5 and assume that it corresponds to the reference domain $D$. More precisely, $D$ consists of the rectangle $\left[a_{1}, b_{1}\right] \times\left[a_{2}, b_{2}\right]$ with a hole of radius $R$ located at $\mathbf{c}=\left(c_{1}, c_{2}\right)$. We assume that the rectangle is fixed and that the center $\mathbf{c}$ of the cylinder is randomly moved along the vertical axis, namely that it is given by $\left(c_{1}, c_{2}+\varepsilon Y\right)$ in $D_{\omega}$ with $Y$ a uniform random variable in $[-1,1]$. We take $\tilde{\mathbf{f}}=\mathbf{0}$ and we prescribe the following inflow and


Figure 3.5: Geometry with prescribed boundary conditions for the first example.
outflow (parabolic) velocity profile on the inlet and outlet part of $\partial D_{\omega}$

$$
\tilde{\mathbf{u}}\left(a_{1}, x_{2}\right)=\tilde{\mathbf{u}}\left(b_{1}, x_{2}\right)=\left(4 U_{\max }\left(x_{2}-a_{2}\right)\left(b_{2}-x_{2}\right) /\left(b_{2}-a_{2}\right)^{2}, 0\right)^{T} \quad \text { for } a_{2} \leq x_{2} \leq b_{2}
$$

with a maximum velocity $U_{\max }=0.3$ achieved at $x_{2}=\frac{a_{2}+b_{2}}{2}$. We impose homogeneous Dirichlet boundary conditions on the remaining parts of the boundary. The Reynolds number is then given by $\frac{2}{3} U_{\max }(2 R) v^{-1}$, where $\frac{2}{3} U_{\max }$ corresponds to the mean velocity.

We choose a mapping $\mathbf{x}_{\omega}$, consistent with the perturbation mentioned above, such that all the boundary nodes are fixed. In such a case, the boundary conditions for the equivalent problem on $D$ are the same than the ones on $D_{\omega}$. More precisely, we consider the mapping $\mathbf{x}_{\omega}: D \rightarrow D_{\omega}$
given componentwise by:

$$
\left[\begin{array}{rl}
x_{1} & =\xi_{1} \\
x_{2} & =\xi_{2}+\varepsilon \varphi_{1}\left(\xi_{1}\right) \varphi_{2}\left(\xi_{2}\right) Y(\omega)
\end{array}\right.
$$

where for $i=1,2$

$$
\varphi_{i}\left(\xi_{i}\right)= \begin{cases}\frac{\xi_{i}-a_{i}}{c_{i}-R-a_{i}}-\tau \frac{\left(\xi_{i}-a_{i}\right)\left(\xi_{i}-c_{i}+R\right)}{\left(c_{i}-R-a_{i}\right)^{2}} & \text { if } \xi_{i} \in\left[a_{i}, c_{i}-R[ \right.  \tag{3.79}\\ 1 & \text { if } \xi_{i} \in\left[c_{i}-R, c_{i}+R\right] \\ \frac{\xi_{i}-b_{i}}{c_{i}+R-b_{i}}-\tau \frac{\left(\xi_{i}-b_{i}\right)\left(\xi_{i}-c_{i}-R\right)}{\left(c_{i}+R-b_{i}\right)^{2}} & \text { if } \left.\left.\xi_{i} \in\right] c_{i}+R, b_{i}\right]\end{cases}
$$

which can be written under the form (3.48) as $\mathbf{x}(\boldsymbol{\xi}, \omega)=\boldsymbol{\xi}+\varepsilon \boldsymbol{\varphi}(\boldsymbol{\xi}) Y(\omega) / \sqrt{3}$ with $Y$ a uniform random variable in $[-\sqrt{3}, \sqrt{3}]$ and $\boldsymbol{\varphi}(\boldsymbol{\xi})=\left(0, \varphi_{1}\left(\xi_{1}\right) \varphi_{2}\left(\xi_{2}\right)\right)^{T}$. The function $\varphi_{2}$ alone fits the required perturbation of the domain but we use the function $\varphi_{1}$ to fix the nodes on the inlet and outlet boundaries. Moreover, the parameter $\tau \in\{0,1\}$ is used to control the regularity of the mapping. Indeed, choosing $\tau=1$ implies that all the functions appearing in the Jacobian matrix $A^{-1}$ of the mapping $\mathbf{x}_{\omega}$ are continuous. From now on, according to [110], we fix the value of the various geometry parameters to $a_{1}=a_{2}=0, b_{1}=2.2, b_{2}=0.41, c_{1}=c_{2}=0.2$ and $R=0.05$, and we choose $\tau=1$. The functions $\varphi_{1}$ and $\varphi_{2}$ for these values of the various geometrical parameters are given in Figure 3.6.


Figure 3.6: Functions $\varphi_{1}\left(\xi_{1}\right), \xi_{1} \in[0,2.2]$ (left) and $\varphi_{2}\left(\xi_{2}\right), \xi_{2} \in[0,0.41]$ (right) defined in (3.79).

The numerical tests are performed using FreeFem++ 3.19.1-1 [78]. The mesh is constructed with a Delaunay triangulation using $n$ equispaced points on the left and right boundaries, $5 n$ on the upper and lower boundaries and $2 n$ along the hole. The mesh size is then given by $h \approx(\sqrt{2} n)^{-1}$ while the number of elements and vertices are about $12 n^{2}$ and $7 n^{2}$, respectively. Notice that we are using piecewise linear triangular elements to mesh the physical domain $D$ whose boundary has a curved part, namely the hole modelling the cylinder. We are not accounting this error here and we refer to [31, Chapter 10] or [48, Chapter VI] for an analysis of such variational crime, introducing for instance isoparametric finite elements. Finally, we recall that the error estimates derived in Sections 3.2.5 and 3.2.5 are valid for homogeneous Dirichlet boundary conditions. In the case of inhomogeneous conditions, as considered here, an additional term due to the approximation of the Dirichlet data should be included.

However, thanks to the fact that the later is not affected by the mapping, it is a higher order term in $h$ (see for instance [16]) and thus we do not take it into account in the numerical results.

## Deterministic case

We first consider the deterministic case, namely when $\varepsilon$ is set to zero. The reference values in [110] include the drag ( $c_{D}$ ) and lift ( $c_{L}$ ) coefficients and the pressure difference $\Delta p=$ $p(0.15,0.2)-p(0.25,0.2)$ between the value at the front and the end point of the cylinder. Using $\mathbb{P}_{2}-\mathbb{P}_{1}$ FE on a mesh with $n=80$, we obtain the values $c_{D}=5.57469, c_{L}=0.0104584$ and $\Delta p=0.117525$ which are consistent with the bounds given in [110].

We give in Figure 3.7 the velocity magnitude, the two components $u_{1}$ and $u_{2}$ and the pressure obtained using $\mathbb{P}_{2}-\mathbb{P}_{1}$ finite elements on the finest mesh, i.e. $n=64$.


Figure 3.7: Velocity magnitude, components $u_{1}$ and $u_{2}$ and pressure for the first problem in the case $\varepsilon=0$ and $v=0.001$.

In Table 3.4, we give the results obtained for various values of $n$ and $v$, where err, $\eta$ and e.i. denote respectively the error, the error estimator $\left(\eta_{h}^{2}+\eta_{\varepsilon}^{2} \frac{1}{2}\right.$ with $\eta_{h}$ and $\eta_{\varepsilon}$ defined in (3.73) and the effectivity index, namely the ratio between the error estimator and the error. Notice that $\eta_{\varepsilon}=0$ here since $\varepsilon=0$. We can see that in all cases, for $h$ small enough, the effectivity index is about 2.8. This value is consistent with the one obtained in Appendix 1.C, see Table 1.15.

Chapter 3. PDEs in random domains

|  | $v=0.001$ |  |  | $v=0.01$ |  |  | $v=0.1$ |  |  | $v=1$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $n$ | err | $\eta$ | e.i. | err | $\eta$ | e.i. | err | $\eta$ | e.i. | err | $\eta$ | e.i. |
| 4 | 0.136 | 0.566 | 4.17 | 0.158 | 0.310 | 1.96 | 0.514 | 0.963 | 1.87 | 1.628 | 3.052 | 1.87 |
| 8 | 0.039 | 0.150 | 3.87 | 0.060 | 0.135 | 2.27 | 0.188 | 0.415 | 2.20 | 0.596 | 1.312 | 2.20 |
| 16 | 0.015 | 0.044 | 2.87 | 0.028 | 0.070 | 2.55 | 0.086 | 0.216 | 2.52 | 0.271 | 0.684 | 2.52 |
| 32 | 0.007 | 0.019 | 2.73 | 0.013 | 0.034 | 2.70 | 0.039 | 0.105 | 2.69 | 0.124 | 0.333 | 2.69 |
| 64 | 0.003 | 0.009 | 2.75 | 0.006 | 0.017 | 2.78 | 0.019 | 0.052 | 2.78 | 0.060 | 0.166 | 2.78 |

Table 3.4: Error, error estimator and effectivity index for the deterministic case $(\varepsilon=0)$ and various viscosities for the first example.

## Random case

We treat now the random case by considering values of $\varepsilon$ between 0 and 0.05 . With $\varepsilon=0.05$, the random position of the cylinder on the vertical axis lies between 0.15 and 0.25 with nominal value in 0.2 , which is quite a large perturbation considering that the height of the rectangle is equal to 0.41 .

The velocity magnitude for the case $v=0.001$ when the cylinder is moved from 0.2 to 0.25 is given in Figure 3.8. We plot the solution obtained when performing the computation on the physical domain and on the reference domain, with the appropriate modification of the coefficients in the equations for the latter case. The solution for the case $\varepsilon=0$ is again given for comparison.


Figure 3.8: Velocity magnitude for $v=0.001$ in the case $\varepsilon=0$ (top) and $\varepsilon=0.05$ with $Y=1$ computed on $D_{\omega}$ (middle) and on $D$ (bottom) for the first example.

We give in Table 3.5 the numerical results obtained for $v=0.001$ and $v=1$ and various values of $n$ and $\varepsilon$.

|  |  | $v=0.001$ |  |  |  |  | $v=1$ |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $n$ | $\varepsilon$ | err | $\eta_{h}$ | $\eta_{\varepsilon}$ | e.i. | err | $\eta_{h}$ | $\eta_{\varepsilon}$ | e.i. |  |
| 4 | 0.05 | 0.1389 | 0.5656 | 1.0649 | 8.68 | 1.8881 | 3.0521 | 2.4890 | 2.09 |  |
| 8 | 0.05 | 0.0591 | 0.1503 | 0.6797 | 11.78 | 1.0157 | 1.3124 | 2.3458 | 2.65 |  |
| 16 | 0.05 | 0.0452 | 0.0440 | 0.5487 | 12.19 | 0.8110 | 0.6839 | 2.3018 | 2.96 |  |
| 32 | 0.05 | 0.0429 | 0.0190 | 0.5288 | 12.32 | 0.7713 | 0.3333 | 2.2887 | 3.00 |  |
| 64 | 0.05 | 0.0428 | 0.0091 | 0.5246 | 12.25 | 0.7526 | 0.1655 | 2.2856 | 3.05 |  |
| 4 | 0.025 | 0.1361 | 0.5656 | 0.5325 | 5.71 | 1.6989 | 3.0521 | 1.2445 | 1.94 |  |
| 8 | 0.025 | 0.0436 | 0.1503 | 0.3399 | 8.52 | 0.7159 | 1.3124 | 1.1729 | 2.46 |  |
| 16 | 0.025 | 0.0249 | 0.0440 | 0.2743 | 11.15 | 0.4701 | 0.6839 | 1.1509 | 2.85 |  |
| 32 | 0.025 | 0.0205 | 0.0190 | 0.2644 | 12.96 | 0.3916 | 0.3333 | 1.1444 | 3.04 |  |
| 64 | 0.025 | 0.0194 | 0.0091 | 0.2623 | 13.51 | 0.3831 | 0.1655 | 1.1428 | 3.01 |  |
| 4 | 0.0125 | 0.1356 | 0.5656 | 0.2662 | 4.61 | 1.6458 | 3.0521 | 0.6223 | 1.89 |  |
| 8 | 0.0125 | 0.0401 | 0.1503 | 0.1699 | 5.66 | 0.6291 | 1.3124 | 0.5865 | 2.29 |  |
| 16 | 0.0125 | 0.0181 | 0.0440 | 0.1372 | 7.98 | 0.3310 | 0.6839 | 0.5755 | 2.70 |  |
| 32 | 0.0125 | 0.0119 | 0.0190 | 0.1322 | 11.25 | 0.2264 | 0.3333 | 0.5722 | 2.92 |  |
| 64 | 0.0125 | 0.0100 | 0.0091 | 0.1311 | 13.13 | 0.2056 | 0.1655 | 0.5714 | 2.89 |  |
| 4 | 0.00625 | 0.1356 | 0.5656 | 0.1331 | 4.29 | 1.6324 | 3.0521 | 0.3111 | 1.88 |  |
| 8 | 0.00625 | 0.0392 | 0.1503 | 0.0850 | 4.41 | 0.6043 | 1.3124 | 0.2932 | 2.23 |  |
| 16 | 0.00625 | 0.0160 | 0.0440 | 0.0686 | 5.08 | 0.2872 | 0.6839 | 0.2877 | 2.58 |  |
| 32 | 0.00625 | 0.0084 | 0.0190 | 0.0661 | 8.17 | 0.1559 | 0.3333 | 0.2861 | 2.82 |  |
| 64 | 0.00625 | 0.0058 | 0.0091 | 0.0656 | 11.45 | 0.1117 | 0.1655 | 0.2857 | 2.96 |  |
| 4 | 0.003125 | 0.1355 | 0.5656 | 0.0666 | 4.20 | 1.6324 | 3.0521 | 0.1556 | 1.88 |  |
| 8 | 0.003125 | 0.0389 | 0.1503 | 0.0425 | 4.01 | 0.6043 | 1.3124 | 0.1466 | 2.23 |  |
| 16 | 0.003125 | 0.0155 | 0.0440 | 0.0343 | 3.60 | 0.2872 | 0.6839 | 0.1439 | 2.58 |  |
| 32 | 0.003125 | 0.0074 | 0.0190 | 0.0330 | 5.18 | 0.1328 | 0.3333 | 0.1430 | 2.73 |  |
| 64 | 0.003125 | 0.0041 | 0.0091 | 0.0328 | 8.32 | 0.0760 | 0.1655 | 0.1429 | 2.88 |  |

Table 3.5: The error, the two contributions $\eta_{h}$ and $\eta_{\varepsilon}$ of the error estimator in (3.73) and the effectivity index for $v=0.001$ and $v=1$ for the first example.

We recall that we use different FE spaces for the reference and the approximate solution and thus, even in the case where the same mesh is used for both solutions, there is still an error due to space discretization. We can see in Table 3.5 that the effectivity index tends to the one obtained in Table 3.4 when the spatial error is dominating while when the statistical error dominates, it is about 13 and 3 for $v=0.001$ and $v=1$, respectively. This highlights the dependence of the error estimate given in Section 3.2 .5 with respect to the input data. However, we can see that when both $h$ and $\varepsilon$ are divided by 2 then the effectivity index remains constant, this observation being tempered by the fact that the effectivity index for $\varepsilon=0$ is not constant for the various meshes considered (see Table 3.4). For instance, in the case $v=0.001$ and $\varepsilon=(5 n)^{-1}$, which corresponds to $h \approx 3.5 \varepsilon$, the effectivity index is about 8 . We study now the efficiency of the second error estimate with respect to the viscosity. In Figure 3.9, we give the effectivity index with respect to $v$ for both error estimators $\eta$ and $\hat{\eta}=\left(\eta_{h}^{2}+\hat{\eta}_{\varepsilon}^{2}\right)^{\frac{1}{2}}$, where $\hat{\eta}_{\varepsilon}$ is given in (3.77), in the case $\varepsilon=0.025, n=64$ and $n_{\text {ref }}=64$, which corresponds to a statistical error dominant regime.


Figure 3.9: Effectivity index with respect to the viscosity $v$ for the two error estimators $\eta$ and $\hat{\eta}$ defined in (3.73) and (3.77) for the first example.

We can see that the effectivity index of the first error estimator $\eta$ remains constant for viscosities greater than 0.01 while below this value, it starts increasing as $v$ decreases. The situation is different for the second estimator $\hat{\eta}$ of Section 3.2.5, whose efficiency is not sensitive to the value of $v$.

Remark 3.2.17. In order to have the correct balance of the two terms appearing in the error estimator $\eta$ or $\hat{\eta}$, we could estimate numerically the constants in front of each term $\eta_{h}$ and $\eta_{\varepsilon}$ or $\hat{\eta}_{\varepsilon}$. The estimation of these constants can also be used to construct a sharp error estimator, namely an error estimator with effectivity index close to 1 . According to the results in Table 3.4, the term $\eta_{h}$ should be multiplied by a factor $1 / 2.8$. For the term due to uncertainty, we obtain that $\hat{\eta}_{\varepsilon}$ should be multiplied by about 1.5, considering for instance same FE spaces and fine mesh for both the reference and approximate solutions, whereas the constant in front of $\eta_{\varepsilon}$
depends on the viscosity as seen in Table 3.5 or Figure 3.9 (for instance $1 / 13$ for $v=0.001$ or $1 / 3$ for $v \geq 0.01$ ).

To conclude the analysis of this first example, we mention that similar results are obtained if we use homogeneous Neumann boundary conditions on the outlet part of the boundary. Notice that in this case, the jump term should be modified appropriately since it is no longer zero on the boundary edges belonging to the outlet.

## Second example

For this second example, the reference geometry $D$ consists in a square $[-H, H]^{2}$ with $H=0.5$ and a circular hole of radius $R=0.15$ centred at the origin, as depicted in Figure 3.10 where the prescribed boundary conditions are also indicated. The shape of the hole is given on $D$


Figure 3.10: Geometry with prescribed boundary conditions for the second example.
by $\left(\xi_{1}, \xi_{2}\right)=(R \cos (\theta), R \sin (\theta))$ with $\theta \in[0,2 \pi]$. We perturb this hole by modifying its radius with respect to the angle by the formula $R+\varepsilon d_{\theta}$, where $d_{\theta}=\sum_{j=1}^{L} \alpha_{j} \cos \left(k_{j} \theta\right) Y_{j}$ and $Y_{j}$ are i.i.d uniform random variables in $[-1,1]$. The coefficients $k_{j}$ and $\alpha_{j}$ control the frequency and the amplitude of each term, respectively. We mention that a similar perturbation is considered in [125], where the mapping is not constructed explicitly but computed through solutions of Laplace equations. We consider here the following mapping $\mathbf{x}_{\omega}$ from $D$ to $D_{\omega}$ which fits the above perturbation: denoting $r=\sqrt{\xi_{1}^{2}+\xi_{2}^{2}}$ and $\theta=\arctan \left(\frac{\xi_{2}}{\xi_{1}}\right)$ the polar coordinates of any point $\boldsymbol{\xi}=\left(\xi_{1}, \xi_{2}\right)$ of $D$, we take

$$
\mathbf{x}=\boldsymbol{\xi}+\varepsilon \sum_{j=1}^{L} \boldsymbol{\varphi}_{j}(\boldsymbol{\xi}) Y_{j}(\omega), \quad \boldsymbol{\varphi}_{j}(\boldsymbol{\xi})=\alpha_{j} \cos \left(k_{j} \theta\right) g(\boldsymbol{\xi})\left[\begin{array}{c}
\cos (\theta)  \tag{3.80}\\
\sin (\theta)
\end{array}\right]
$$

where the cutoff function $g$ is such that it vanishes at the boundary of the domain and is equal to 1 in the hole, namely we use

$$
g(\boldsymbol{\xi})= \begin{cases}1 & \text { if } r \in[0, R]  \tag{3.81}\\ \frac{\left(\xi_{1}^{2}-H^{2}\right)\left(\xi_{2}^{2}-H^{2}\right)}{\left(R^{2} \xi_{1}^{2} r^{-2}-H^{2}\right)\left(R^{2} \xi_{2}^{2} r^{-2}-H^{2}\right)} & \text { otherwise }\end{cases}
$$

The graph of this function is depicted in Figure 3.11


Figure 3.11: Function $g=g\left(\xi_{1}, \xi_{2}\right)$ defined in (3.81).

The mesh is again built with a Delaunay triangulation using $n$ equispaced points on the boundaries of the square and $2 n$ on the hole for various values of $n$ with corresponding mesh size $h \approx 1.5 n^{-1}$ and number of elements and vertices of about $3.5 n^{2}$ and $2 n^{2}$, respectively.

Remark 3.2.18. Contrary to the previous example, the choice of the boundary conditions on the outlet has an impact on the solution of this problem, due to the fact that the outlet is close to the cylinder. This is especially true for small viscosities, in which case some flow is re-entering the domain when homogeneous Neumann conditions are used while the solution presents a boundary layer when Dirichlet conditions are enforced.

## Deterministic case

We consider first the deterministic case taking $\varepsilon=0$. The plot of the velocity magnitude, the two components $u_{1}$ and $u_{2}$ and the pressure obtained using $\mathbb{P}_{2}-\mathbb{P}_{1} \mathrm{FE}$ and the finest mesh ( $n=160$ ) is given in Figure 3.12.

Moreover, we give in Table 3.6 the results we get for various values of $n$ and $v$. Similarly to the previous example, the effectivity index is about 2.8 in all cases, when $h$ is small enough.


Figure 3.12: From left to right: velocity magnitude, components $u_{1}$ and $u_{2}$ and pressure for the second problem in the case $\varepsilon=0$ and $v=0.05$.

|  | $v=0.05$ |  |  |  | $v=0.1$ |  |  |  | $v=0.5$ |  |  | $v=1$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $n$ | err | $\eta$ | e.i. | err | $\eta$ | e.i. | err | $\eta$ | e.i. | err | $\eta$ | e.i. |  |  |
| 10 | 0.477 | 1.149 | 2.41 | 0.621 | 1.405 | 2.26 | 1.364 | 2.988 | 2.19 | 1.930 | 4.221 | 2.19 |  |  |
| 20 | 0.230 | 0.579 | 2.51 | 0.278 | 0.697 | 2.51 | 0.590 | 1.470 | 2.49 | 0.833 | 2.074 | 2.49 |  |  |
| 40 | 0.112 | 0.294 | 2.63 | 0.132 | 0.353 | 2.67 | 0.279 | 0.745 | 2.68 | 0.393 | 1.052 | 2.68 |  |  |
| 80 | 0.055 | 0.148 | 2.71 | 0.064 | 0.176 | 2.75 | 0.134 | 0.371 | 2.76 | 0.190 | 0.523 | 2.76 |  |  |
| 160 | 0.026 | 0.073 | 2.77 | 0.031 | 0.087 | 2.80 | 0.066 | 0.184 | 2.80 | 0.096 | 0.259 | 2.80 |  |  |

Table 3.6: Error, error estimator and effectivity index for the deterministic case $(\varepsilon=0)$ and various viscosities for the second example.

## Random case

We consider first $L=1$ random variable, we fix $\alpha_{1}=1$ and $k_{1}=6$ in the definition of $d_{\theta}$ and we let $0 \leq \varepsilon \leq 0.01$. The vorticity of the velocity $\mathbf{u}$ and the pressure $p$ in the case $\varepsilon=0.01, v=0.05$ and $Y=1$ is given in Figure 3.13, where the solution obtained by solving the problem defined on $D_{\omega}$ as well as the solution for the case $\varepsilon=0$ are also given for comparison.

We give in Table 3.7 the numerical results obtained for $v=0.05$ and $v=1$ and various values of $n$ and $\varepsilon$.

Similarly to the previous example, we observe that the effectivity index tends to the one obtained for the deterministic case $(\varepsilon=0)$ when the error in $h$ is dominating, while it is about 6 and 1.5 for $v=0.05$ and $v=1$, respectively, when the statistical error dominates. This shows again the sensitivity of the efficiency of the first error estimator with respect to the input data but, as before, the effectivity index remains about constant when both $h$ and $\varepsilon$ are divided by 2. Indeed, for instance for $v=0.05$ and $\varepsilon=(10 n)^{-1}$, corresponding to $h \approx 15 \varepsilon$, it stays between 3.81 and 4.05. Finally, the same behaviour than in the previous example is observed for the efficiency of the second error estimator $\hat{\eta}$ with respect to the viscosity, as can be seen in Figure 3.14 where the results are given for the case $\varepsilon=0.005, n=160$ and $n_{\text {ref }}=160$.

The results are similar when we consider other kinds of perturbation. For instance, let consider


Figure 3.13: Vorticity of the velocity and pressure for $v=0.05$ in the case $\varepsilon=0$ (left) and $\varepsilon=0.01$ with $Y=1$ computed $D_{\omega}$ (middle) and on $D$ (right) for the second example.


Figure 3.14: Effectivity index with respect to the viscosity $v$ for the two error estimators $\eta$ and $\hat{\eta}$ defined in (3.73) and (3.77) for the second example.

|  |  | $v=0.05$ |  |  |  |  | $v=1$ |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $n$ | $\varepsilon$ | err | $\eta_{h}$ | $\eta_{\varepsilon}$ | e.i. | err | $\eta_{h}$ | $\eta_{\varepsilon}$ | e.i. |  |
| 10 | 0.01 | 0.5125 | 1.1492 | 1.6181 | 3.87 | 2.0403 | 4.2209 | 1.4479 | 2.19 |  |
| 20 | 0.01 | 0.3251 | 0.5785 | 1.5682 | 5.14 | 1.2200 | 2.0741 | 1.3862 | 2.04 |  |
| 40 | 0.01 | 0.2625 | 0.2937 | 1.5552 | 6.03 | 1.0216 | 1.0524 | 1.3730 | 1.69 |  |
| 80 | 0.01 | 0.2486 | 0.1478 | 1.5519 | 6.27 | 1.0040 | 0.5233 | 1.3696 | 1.46 |  |
| 160 | 0.01 | 0.2431 | 0.07279 | 1.5511 | 6.39 | 0.9630 | 0.2594 | 1.3687 | 1.45 |  |
| 10 | 0.005 | 0.4859 | 1.1492 | 0.8090 | 2.89 | 1.9575 | 4.2209 | 0.7240 | 2.19 |  |
| 20 | 0.005 | 0.2556 | 0.5785 | 0.7841 | 3.81 | 0.9477 | 2.0741 | 0.6931 | 2.31 |  |
| 40 | 0.005 | 0.1628 | 0.2937 | 0.7776 | 5.11 | 0.6163 | 1.0524 | 0.6865 | 2.04 |  |
| 80 | 0.005 | 0.1340 | 0.1478 | 0.7759 | 5.91 | 0.5149 | 0.5233 | 0.6848 | 1.67 |  |
| 160 | 0.005 | 0.1238 | 0.0728 | 0.7755 | 6.29 | 0.4891 | 0.2594 | 0.6843 | 1.50 |  |
| 10 | 0.0025 | 0.4792 | 1.1492 | 0.4045 | 2.54 | 1.9363 | 4.2209 | 0.3620 | 2.19 |  |
| 20 | 0.0025 | 0.2370 | 0.5785 | 0.3921 | 2.95 | 0.8602 | 2.0741 | 0.3465 | 2.44 |  |
| 40 | 0.0025 | 0.1263 | 0.2937 | 0.3888 | 3.86 | 0.4538 | 1.0524 | 0.3433 | 2.44 |  |
| 80 | 0.0025 | 0.0808 | 0.1478 | 0.3880 | 5.14 | 0.3085 | 0.5233 | 0.3424 | 2.03 |  |
| 160 | 0.0025 | 0.0662 | 0.0728 | 0.3878 | 5.96 | 0.2584 | 0.2594 | 0.3422 | 1.66 |  |
| 10 | 0.00125 | 0.4776 | 1.1492 | 0.2023 | 2.44 | 1.9317 | 4.2209 | 0.1810 | 2.19 |  |
| 20 | 0.00125 | 0.2319 | 0.5785 | 0.1960 | 2.63 | 0.8399 | 2.0741 | 0.1733 | 2.48 |  |
| 40 | 0.00125 | 0.1154 | 0.2937 | 0.1944 | 3.05 | 0.4098 | 1.0524 | 0.1716 | 2.60 |  |
| 80 | 0.00125 | 0.0624 | 0.1478 | 0.1940 | 3.91 | 0.2237 | 0.5233 | 0.1712 | 2.46 |  |
| 160 | 0.00125 | 0.0405 | 0.0728 | 0.1939 | 5.12 | 0.1517 | 0.2594 | 0.1711 | 2.05 |  |
| 10 | 0.000625 | 0.4772 | 1.1492 | 0.1011 | 2.42 | 1.9304 | 4.2209 | 0.0905 | 2.19 |  |
| 20 | 0.000625 | 0.2306 | 0.5785 | 0.0980 | 2.54 | 0.8347 | 2.0741 | 0.0866 | 2.49 |  |
| 40 | 0.000625 | 0.1125 | 0.2937 | 0.0972 | 2.75 | 0.3977 | 1.0524 | 0.0858 | 2.66 |  |
| 80 | 0.000625 | 0.0565 | 0.1479 | 0.0970 | 3.13 | 0.1987 | 0.5233 | 0.0856 | 2.67 |  |
| 160 | 0.000625 | 0.0304 | 0.0728 | 0.0970 | 3.99 | 0.1101 | 0.2594 | 0.0855 | 2.48 |  |

Table 3.7: The error, the two contributions $\eta_{h}$ and $\eta_{\varepsilon}$ of the estimator in (3.73) and the effectivity index for $v=0.05$ and $v=1$.

## Chapter 3. PDEs in random domains

(3.80) with $L=2$ with $k_{1}=6, k_{2}=11, \alpha_{1}=1$ and $\alpha_{2}=0.8$. The results we obtained, given in Figure 3.15 and in Table 3.8, are very similar to those presented in Table 3.7. The results for the second error estimator $\hat{\eta}$ with the estimated constant, see Remark 3.2.17, are also provided. We can see that for $h$ small enough, namely when the effectivity index for the spatial error estimator is about 2.8 (see Table 3.6), the error estimator is sharp.


Figure 3.15: Vorticity of the velocity and pressure for $v=0.05$ in the case $\varepsilon=0$ (left) and $\varepsilon=0.01$ with $Y=1$ computed $D_{\omega}$ (middle) and on $D$ (right) for the second example with $L=2$.

## Conclusions

In this chapter, we have considered steady-state nonlinear PDEs on random domains, namely the one-dimensional viscous Burgers' equation and the incompressible Navier-Stokes equations. We have used the domain mapping method to transform them into PDEs on a fixed reference domain with random coefficients.

We have first studied the deterministic Burgers' equation with mixed Dirichlet-Neumann boundary conditions. We have shown the well-posedness of the problem under suitable assumptions on the input data and we have derived an a posteriori error estimate. Then, the case of random intervals has been considered, performing all the analysis on the fixed reference domain. Finally, we have presented two numerical examples both in the deterministic and random cases.

For the Navier-Stokes equations, we started the analysis by showing the well-posedness of the

| $n$ | $\varepsilon$ | err | $\eta_{h}$ | $\eta_{\varepsilon}$ | $\eta / \mathrm{err}$ | $\eta_{h} / 2.8$ | $1.5 \hat{\eta}_{\varepsilon}$ | $\hat{\eta} / \mathrm{err}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 10 | 0.005 | 0.4994 | 1.1492 | 1.4301 | 3.67 | 0.4104 | 0.1849 | 0.90 |
| 20 | 0.005 | 0.2924 | 0.5785 | 1.3884 | 5.14 | 0.2066 | 0.1992 | 0.98 |
| 40 | 0.005 | 0.22061 | 0.2937 | 1.3768 | 6.38 | 0.1049 | 0.2054 | 1.05 |
| 80 | 0.005 | 0.1983 | 0.1478 | 1.3739 | 6.97 | 0.0528 | 0.2072 | 1.08 |
| 160 | 0.005 | 0.1928 | 0.0728 | 1.3732 | 7.13 | 0.0260 | 0.2077 | 1.09 |
| 10 | 0.0025 | 0.4826 | 1.1492 | 0.7151 | 2.80 | 0.4104 | 0.0924 | 0.87 |
| 20 | 0.0025 | 0.2477 | 0.5785 | 0.6942 | 3.65 | 0.2066 | 0.0996 | 0.93 |
| 40 | 0.0025 | 0.1464 | 0.2937 | 0.6884 | 5.11 | 0.1049 | 0.1027 | 1.00 |
| 80 | 0.0025 | 0.1080 | 0.1478 | 0.6869 | 6.51 | 0.0528 | 0.1036 | 1.08 |
| 160 | 0.0025 | 0.0988 | 0.0728 | 0.6866 | 6.99 | 0.0260 | 0.1038 | 1.08 |
| 10 | 0.00125 | 0.4784 | 1.1492 | 0.3575 | 2.52 | 0.4104 | 0.0462 | 0.86 |
| 20 | 0.00125 | 0.2345 | 0.5785 | 0.3471 | 2.88 | 0.2066 | 0.0498 | 0.91 |
| 40 | 0.00125 | 0.1212 | 0.2937 | 0.3442 | 3.73 | 0.1049 | 0.0513 | 0.96 |
| 80 | 0.00125 | 0.0731 | 0.1478 | 0.3435 | 5.12 | 0.0528 | 0.0518 | 1.01 |
| 160 | 0.00125 | 0.0545 | 0.0728 | 0.3433 | 6.44 | 0.0260 | 0.0519 | 1.06 |

Table 3.8: Effectivity index of the two error estimators in the case $v=0.05$ for the second example with $L=2$.
problem under suitable assumptions on the input data and the mapping, before performing an a posteriori error analysis. Using a perturbation method, we obtained two error estimates for the first order approximation $(\mathbf{u}, p) \approx\left(\mathbf{u}_{0, h}, p_{0, h}\right)$. Both estimates are constituted of two parts, namely one part due to space discretization in $h$ and one due to the uncertainty in $\varepsilon$. They already give useful information, especially when the problem contains small uncertainties. They can indeed be used to adaptively find a spatial mesh that balances the two sources of error. Further mesh refinement should then be avoided since it would not decrease the total error, the statistical error being dominant. The latter can only be decreased by adding more terms in the expansion of the solution. Notice that if we want to analyse higher order approximations in $\varepsilon$, then we should impose additional regularity assumptions on $f$ and on the random mapping, namely that the Jacobian matrix $\nabla \boldsymbol{\varphi}_{j}$ belongs to $\left[W^{1, \infty}(D)\right]^{d \times d}$ for $j=0,1, \ldots, L$ and not only for $j=0$. Indeed, we have that the residual for the FE approximation ( $U_{j, h}, P_{j, h}$ ) of ( $U_{j}, P_{j}$ ) belongs to $L^{2}(D)$ for $j=1, \ldots, L$, where $\left(U_{j}, P_{j}\right)$ is the solution of (3.59) and appears in the second term of the expansion of the solution. The same holds for the residual of the higher order terms.

Each of the two error estimators $\eta$ and $\hat{\eta}$ that we obtained presents its advantages and drawbacks. The first one can be computed by solving only one nonlinear problem, namely the standard Navier-Stokes equations on the reference domain. We have seen however that the sharpness of this estimator might be affected when changing the input data, as predicted by the theory. In the two numerical examples considered here, the effectivity index remains constant for moderate Reynolds numbers but then starts to increase as the viscosity diminishes. The second error estimator shows promising results, its efficiency being indeed independent
of the input data for all the cases we have considered. The extra cost to pay is the resolution of $L$ additional linear problems. Finally, as mentioned in Remark 3.2.17, the constant in front of the two terms in $h$ and $\varepsilon$ can be estimated numerically (once for all for the second estimator) to get a sharp error estimator, that is an estimator with effectivity index close to 1 .

## 3.A Derivation of problems (3.58) and (3.59)

We give here some details about the derivation of the problems (3.58) and (3.59) that we need to solve to obtain the first two terms in the expansion of the solution ( $\mathbf{u}, p$ ), namely $\left(\mathbf{u}_{0}, p_{0}\right)$ and $\left(\mathbf{u}_{1}, p_{1}\right)$. These problems are obtained by replacing each term in (3.35), the problem in strong form for $(\mathbf{u}, p)$, by its expansion with respect to $\varepsilon$ and keeping only the appropriate terms. Using relations (3.55) and (3.56), we can write

$$
\begin{aligned}
J_{\mathbf{x}} A A^{T} & =\left(1+\varepsilon \operatorname{tr}\left(A_{1}\right)+\mathscr{O}\left(\varepsilon^{2}\right)\right)\left(I-\varepsilon A_{1}+\mathscr{O}\left(\varepsilon^{2}\right)\right)\left(I-\varepsilon A_{1}^{T}+\mathscr{O}\left(\varepsilon^{2}\right)\right) \\
& =I+\varepsilon\left(\operatorname{tr}\left(A_{1}\right) I-A_{1}-A_{1}^{T}\right)+\mathscr{O}\left(\varepsilon^{2}\right)
\end{aligned}
$$

and similarly

$$
J_{\mathbf{x}} A^{T}=I+\varepsilon\left(\operatorname{tr}\left(A_{1}\right) I-A_{1}^{T}\right)+\mathscr{O}\left(\varepsilon^{2}\right) .
$$

Therefore, considering for instance the convection term, we get

$$
\begin{aligned}
\left(\mathbf{u} \cdot J_{\mathbf{x}} A^{T} \nabla\right) \mathbf{u} & =\left(\left(\mathbf{u}_{0}+\varepsilon \mathbf{u}_{1}+\mathscr{O}\left(\varepsilon^{2}\right)\right) \cdot\left(I+\varepsilon\left(\operatorname{tr}\left(A_{1}\right) I-A_{1}^{T}\right)+\mathscr{O}\left(\varepsilon^{2}\right)\right) \nabla\right)\left(\mathbf{u}_{0}+\varepsilon \mathbf{u}_{1}+\mathscr{O}\left(\varepsilon^{2}\right)\right) \\
& =\left(\mathbf{u}_{0} \cdot \nabla\right) \mathbf{u}_{0}+\varepsilon\left[\left(\mathbf{u}_{1} \cdot \nabla\right) \mathbf{u}_{0}+\left(\mathbf{u}_{0} \cdot \nabla\right) \mathbf{u}_{1}+\left(\mathbf{u}_{0} \cdot\left(\operatorname{tr}\left(A_{1}\right) I-A_{1}^{T}\right) \nabla\right) \mathbf{u}_{0}\right]+\mathscr{O}\left(\varepsilon^{2}\right) .
\end{aligned}
$$

Proceeding similarly for all the terms involved in the first equation of (3.35) and keeping the $\mathscr{O}(1)$ terms with respect to $\varepsilon$ we obtain

$$
-v \Delta \mathbf{u}_{0}+\left(\mathbf{u}_{0} \cdot \nabla\right) \mathbf{u}_{0}+\nabla p_{0}=\mathbf{f}_{0}
$$

which is the first equation of (3.58). If we collect now the terms of order $\mathscr{O}(\varepsilon)$ we get

$$
\begin{align*}
-v \Delta \mathbf{u}_{1}+\left(\mathbf{u}_{0} \cdot \nabla\right) \mathbf{u}_{1}+\left(\mathbf{u}_{1} \cdot \nabla\right) \mathbf{u}_{0}+\nabla p_{1}= & \operatorname{tr}\left(A_{1}\right) \mathbf{f}_{0}+\mathbf{f}_{1}+v \nabla \cdot\left[\left(\left(\operatorname{tr}\left(A_{1}\right) I-A_{1}-A_{1}^{T}\right) \nabla\right) \mathbf{u}_{0}\right] \\
& -\left(\mathbf{u}_{0} \cdot\left(\operatorname{tr}\left(A_{1}\right) I-A_{1}^{T}\right) \nabla\right) \mathbf{u}_{0}-\left(\left(\operatorname{tr}\left(A_{1}\right) I-A_{1}^{T}\right) \nabla\right) p_{0} . \tag{3.82}
\end{align*}
$$

Finally, since

$$
A_{1}=\sum_{j=1}^{L} \nabla \boldsymbol{\varphi}_{j} y_{j}, \quad \mathbf{f}_{1}=\sum_{j=1}^{L} \mathbf{F}_{j} y_{j}, \quad \mathbf{u}_{1}=\sum_{j=1}^{L} \mathbf{U}_{j} y_{j} \quad \text { and } \quad p_{1}=\sum_{j=1}^{L} P_{j} y_{j}
$$

equation (3.82) is satisfied if

$$
\begin{align*}
-v \Delta U_{j}+\left(\mathbf{u}_{0} \cdot \nabla\right) \mathbf{U}_{j}+\left(\mathbf{U}_{j} \cdot \nabla\right) \mathbf{u}_{0}+\nabla P_{j}= & \operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) \mathbf{f}_{0}+\mathbf{F}_{j} \\
& +v \nabla \cdot\left[\left(\left(\operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) I-\nabla \boldsymbol{\varphi}_{j}-\nabla \boldsymbol{\varphi}_{j}^{T}\right) \nabla\right) \mathbf{u}_{0}\right] \\
& -\left(\mathbf{u}_{0} \cdot\left(\operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) I-\nabla \boldsymbol{\varphi}_{j}^{T}\right) \nabla\right) \mathbf{u}_{0} \\
& -\left(\left(\operatorname{tr}\left(\nabla \boldsymbol{\varphi}_{j}\right) I-\nabla \boldsymbol{\varphi}_{j}^{T}\right) \nabla\right) p_{0} \tag{3.83}
\end{align*}
$$

for $j=1, \ldots, L$, which is the second equation of problem (3.59). In fact, relations (3.82) and (3.83) are equivalent since the random variables $\left\{Y_{j}\right\}$ are independent, with zero mean and unit variance and thus form an orthonormal set. The second equation of (3.35), corresponding to the incompressibility constraint, is treated analogously.

## 3.B Choice of the norm

We give here three justifications about the choice of the norm on the space $V \times Q$ for the couple $(\mathbf{u}, p)$, more precisely about the scaling with respect to the kinematic viscosity $v$. We claim that the appropriate scaling is given by

$$
\begin{equation*}
\|\mid \mathbf{v}, q\|_{k}^{2}:=v^{k}\|\nabla \mathbf{v}\|^{2}+v^{k-2}\|q\|^{2} \quad \text { for any choice } k=0,1,2 \tag{3.84}
\end{equation*}
$$

First of all, we can perform a dimensional analysis. The dimension unit of the kinematic viscosity is $[v]=\frac{m^{2}}{s}$ while we have, recall that $p$ corresponds to the pressure divided by the density of the fluid,

$$
\left[|\nabla \mathbf{u}|^{2}\right]=\left(\frac{1}{m} \cdot \frac{m}{s}\right)^{2}=\frac{1}{s^{2}} \quad \text { and } \quad\left[p^{2}\right]=\left(\frac{N}{m^{2}} \cdot \frac{m^{3}}{k g}\right)^{2}=\frac{m^{4}}{s^{4}}
$$

from which we deduce that $\left[v^{k}|\nabla \mathbf{u}|\right]=\left[v^{k-2} p^{2}\right]$ for all $k$. This is also the natural choice of scaling that arises when looking at the a priori estimates on the solution ( $\mathbf{u}, p$ ) or when performing a posteriori error estimation. For simplicity, let us consider the (deterministic) Stokes problem given under the weak form by:
find $(\mathbf{u}, p) \in V \times Q$ such that

$$
\begin{aligned}
a(\mathbf{u}, \mathbf{v})+b(\mathbf{v}, p) & =F(\mathbf{v}) & & \forall \mathbf{v} \in V \\
b(\mathbf{u}, q) & =0 & & \forall q \in Q
\end{aligned}
$$

with $V=\left[H_{0}^{1}(D)\right]^{d}, Q=L_{0}^{2}(D), a(\mathbf{u}, \mathbf{v})=v \int_{D} \nabla \mathbf{u}: \nabla \mathbf{v}, b(\mathbf{v}, q)=-\int_{D} q \nabla \cdot \mathbf{v}$ and $F(\mathbf{v})=\int_{D} \mathbf{f} \cdot \mathbf{v}$. The bilinear form $a$ is continuous and coercive on $V$ with constant $v$ and $b$ is continuous on $V$ with constant 1 and satisfies the inf-sup condition with constant $\beta=\beta(D)$. The problem is
thus well-posed (see [32]) and the following a priori estimates are satisfied

$$
\|\nabla \mathbf{u}\| \leq \frac{1}{v}\|F\|_{V^{\prime}} \quad \text { and } \quad\|p\| \leq \frac{1}{\beta}\left(\|F\|_{V^{\prime}}+v\|\nabla \mathbf{u}\|\right) \leq \frac{2}{\beta}\|F\|_{V^{\prime}} .
$$

Therefore, we have

$$
v^{k / 2}\|\nabla \mathbf{u}\|+v^{k / 2-1}\|p\| \leq C v^{k / 2-1}\|\mathbf{f}\|_{V^{\prime}} \quad \forall k
$$

where $C=(1+2 / \beta)$ is independent of $v$, which is consistent with the scaling (3.84). Finally, for the $a$ posteriori error analysis, denoting $e=\mathbf{u}-\mathbf{u}_{h}$ and $E=p-p_{h}$ with $\mathbf{u}_{h}$ and $p_{h}$ the finite element approximation of $\mathbf{u}$ and $p$, respectively, we have for any $(\mathbf{v}, q) \in V \times Q$

$$
\begin{equation*}
a(\mathbf{e}, \mathbf{v})+b(\mathbf{v}, E)+b(\mathbf{e}, q)=R_{1}(\mathbf{v})+R_{2}(q) \tag{3.85}
\end{equation*}
$$

with

$$
R_{1}(\mathbf{v}):=F(\mathbf{v})-a\left(\mathbf{u}_{h}, \mathbf{v}\right)-b\left(\mathbf{v}, p_{h}\right) \quad \text { and } \quad R_{2}(q):=-b\left(\mathbf{u}_{h}, q\right) .
$$

Using relation (3.85), Young's inequality and the properties of $a$ and $b$, we can easily show that

$$
\begin{equation*}
\|E\| \leq \frac{1}{\beta}\left\|R_{1}\right\|_{V^{\prime}}+\frac{v}{\beta}\|\nabla \mathbf{e}\| \tag{3.86}
\end{equation*}
$$

and

$$
\begin{equation*}
v\|\nabla \mathbf{e}\|^{2} \leq \frac{c_{1}}{v}\left\|R_{1}\right\|_{V^{\prime}}^{2}+\frac{c_{2} v}{\beta^{2}}\left\|R_{2}\right\|_{Q^{\prime}}^{2} \tag{3.87}
\end{equation*}
$$

with for instance $c_{1}=c_{2}=3$, the value of these constants depending only on how we use Young's inequality. From the last two inequalities, we deduce that the scaling (3.84) should be used to get

$$
v^{k}\|\nabla \mathbf{e}\|^{2}+v^{k-2}\|E\|^{2} \leq C\left(v^{k-2}\left\|R_{1}\right\|_{V^{\prime}}^{2}+v^{k}\left\|R_{2}\right\|_{Q^{\prime}}^{2}\right),
$$

where $C$ is a constant independent of $v$ (but which depends on the inf-sup constant $\beta$ ).
We mention that in a diffusion-dominating regime, the choice $k=0$ yields a total error $\|\mathbf{e}, E\|_{0}$ which remains constant when $v$ varies. Indeed, in such a case the velocity error $\|\nabla \mathbf{e}\|$ is constant while the pressure error $\|E\|$ behaves as $v$, i.e. $\frac{1}{v}\|E\|$ is constant.

## 3.C Proof of some properties

Proposition 3.C.1. Let $A, B, C \in \mathbb{R}^{n \times n}$ be square matrices with coefficients denoted respectively by $a_{i j}, b_{i j}$ and $c_{i j}$ for $1 \leq i, j \leq n$, and let $\mathbf{w}$ be any smooth function with value in $\mathbb{R}^{n}$. We then have

$$
\begin{equation*}
A B: C B=A B B^{T}: C \tag{3.88}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(B^{T} \nabla\right) \mathbf{w}=\nabla \mathbf{w} B . \tag{3.89}
\end{equation*}
$$

Proof. We first show (3.88). For the term on the left-hand side, we have

$$
A B: C B=\sum_{i, j=1}^{n}(A B)_{i j}(C B)_{i j}=\sum_{i, j=1}^{n}\left(\sum_{l=1}^{n} a_{i l} b_{l j}\right)\left(\sum_{k=1}^{n} c_{i k} b_{k j}\right)=\sum_{i, j, k, l=1}^{n} a_{i l} b_{l j} c_{i k} b_{k j}
$$

while for the right-hand side, we get

$$
A B B^{T}: C=\sum_{i, k=1}^{n}\left(A B B^{T}\right)_{i k}(C)_{i k}=\sum_{i, k=1}^{n} \sum_{j=1}^{n}(A B)_{i j}\left(B^{T}\right)_{j k}(C)_{i k}=\sum_{i, j, k, l=1}^{n} a_{i l} b_{l j} b_{k j} c_{i k}
$$

We now prove (3.89). From the definition of the gradient operator applied to a vector field, we have

$$
\left(B^{T} \nabla\right) \mathbf{w}=\left(\begin{array}{c}
\left(\left(B^{T} \nabla\right) w_{1}\right)^{T} \\
\vdots \\
\left(\left(B^{T} \nabla\right) w_{n}\right)^{T}
\end{array}\right)=\left(\begin{array}{ccc}
\left(B^{T} \nabla\right)_{1} w_{1} & \cdots & \left(B^{T} \nabla\right)_{n} w_{1} \\
\vdots & \ddots & \vdots \\
\left(B^{T} \nabla\right)_{1} w_{n} & \cdots & \left(B^{T} \nabla\right)_{n} w_{n}
\end{array}\right)
$$

where $w_{i}$ denotes the $i^{\text {th }}$ component of $\mathbf{w}$, and thus

$$
\left[\left(B^{T} \nabla\right) \mathbf{w}\right]_{i j}=\left(B^{T} \nabla\right)_{j}(\mathbf{w})_{i}
$$

Therefore, the coefficient of the $i^{t h}$-row and $j^{t h}$-column of the $n \times n$ matrix $\left(B^{T} \nabla\right) \mathbf{w}$ is given by

$$
\left[\left(B^{T} \nabla\right) \mathbf{w}\right]_{i j}=\sum_{k=1}^{n}\left(B^{T}\right)_{j k}(\nabla)_{k} w_{i}=\sum_{k, l=1}^{n} b_{k j} \frac{\partial w_{i}}{\partial \xi_{k}}=\sum_{k=1}^{n}(\nabla \mathbf{w})_{i k}(B)_{k j}=(\nabla \mathbf{w} B)_{i j}
$$

We now show the relation (3.34) used in Section 3.2.2 to write the strong formulation of the problem on $D$. It can be proven by an integration by part back on the random domain $D_{\omega}$ or using the Piola identity $\nabla \cdot\left(J_{\mathbf{x}} A^{T}\right)=\mathbf{0}$ (see [101] for instance). Indeed, we have

$$
\int_{D} q\left|J_{\mathbf{x}}\right|\left(A^{T} \nabla\right) \cdot \mathbf{v} d \boldsymbol{\xi}=\int_{D_{\omega}} \tilde{q} \nabla_{\mathbf{x}} \cdot \tilde{\mathbf{v}} d \mathbf{x}=-\int_{D_{\omega}} \nabla_{\mathbf{x}} \tilde{q} \cdot \tilde{\mathbf{v}} d \mathbf{x}=-\int_{D}\left|J_{\mathbf{x}}\right|\left(A^{T} \nabla q\right) \cdot \mathbf{v} d \boldsymbol{\xi}
$$

which yields (3.34) since $J_{\mathbf{x}}$ is either positive or negative, depending if the orientation is preserved or not by the mapping. Using the second alternative, since $\nabla \cdot\left(J_{\mathbf{x}} A \mathbf{v}\right)=\left(\nabla \cdot\left(J_{\mathbf{x}} A^{T}\right)\right)$. $\mathbf{v}+\left(J_{\mathbf{x}} A^{T} \nabla\right) \cdot \mathbf{v}$ we have

$$
\begin{aligned}
\int_{D} q J_{\mathbf{x}}\left(A^{T} \nabla\right) \cdot \mathbf{v} d \boldsymbol{\xi} & =\int_{D} q \nabla \cdot\left(J_{\mathbf{x}} A \mathbf{v}\right) d \boldsymbol{\xi}-\int_{D}(\underbrace{\nabla \cdot\left(J_{\mathbf{x}} A^{T}\right)}_{=\mathbf{0}}) \cdot(q \mathbf{v}) d \boldsymbol{\xi} \\
& =-\int_{D} J_{\mathbf{x}}\left(A^{T} \nabla q\right) \cdot \mathbf{v} d \boldsymbol{\xi}
\end{aligned}
$$

Be aware that in [101], the divergence operator applied to a tensor field is defined as the divergence applied to its transposed according to the definition used here. Recall that here we
defined $\left[\nabla \cdot\left(J_{\mathbf{x}} A^{T}\right)\right]_{i}=\sum_{j=1}^{d} \frac{\partial}{\partial \xi_{j}}\left(J_{\mathbf{x}}\left(A^{T}\right)_{i j}\right)=\sum_{j=1}^{d} \frac{\partial}{\partial \xi_{j}}\left(J_{\mathbf{x}} \frac{\partial\left(\xi_{\omega}\right)_{j}}{\partial x_{i}} \circ \mathbf{X}_{\omega}\right)$ for $i=1, \ldots, d$. Moreover, we mention that the Piola identity, which is easily obtained for smooth functions, say $C^{2}$ functions, is still valid (in a weak sense) for less regular functions such as $H^{1}$ functions (see for instance [12, 47]).

Finally, we derive the bound for the term $\mathrm{II}_{3}=b\left(\mathbf{v}, p_{0, h} ; \mathbf{y}_{0}\right)-b\left(\mathbf{v}, p_{0, h} ; \mathbf{y}\right)$ that appear in the proof of Proposition 3.2.14. Writing $\boldsymbol{\xi}=\left(\xi_{1}, \xi_{2}\right)$ and $\mathbf{v}=\left(\nu_{1}, \nu_{2}\right)^{T}$, the two terms in component form read

$$
b\left(\mathbf{v}, p_{0, h} ; \mathbf{y}_{0}\right)=-\int_{D} p_{0, h} \nabla \cdot \mathbf{v} d \boldsymbol{\xi}=-\int_{D} p_{0, h}\left(\frac{\partial \nu_{1}}{\partial \xi_{1}}+\frac{\partial \nu_{2}}{\partial \xi_{2}}\right) d \boldsymbol{\xi}
$$

and

$$
\begin{aligned}
b\left(\mathbf{v}, p_{0, h} ; \mathbf{y}\right) & =-\int_{D} p_{0, h} J_{\mathbf{x}}\left(A^{T} \nabla\right) \cdot \mathbf{v} d \boldsymbol{\xi} \\
& =-\int_{D} p_{0, h} J_{\mathbf{x}}\left(A_{11} \frac{\partial v_{1}}{\partial \xi_{1}}+A_{21} \frac{\partial \nu_{1}}{\partial \xi_{2}}+A_{12} \frac{\partial v_{2}}{\partial \xi_{1}}+A_{22} \frac{\partial v_{2}}{\partial \xi_{2}}\right) d \boldsymbol{\xi} .
\end{aligned}
$$

Subtracting these two terms and using (both continuous and discrete version of) CauchySchwarz's inequality we finally obtain

$$
\begin{aligned}
\mathrm{II}_{3}= & \int_{D}\left(J_{\mathbf{x}} A_{11}-1\right) p_{0, h} \frac{\partial v_{1}}{\partial \xi_{1}} d \boldsymbol{\xi}+\int_{D} J_{\mathbf{x}} A_{21} p_{0, h} \frac{\partial \nu_{1}}{\partial \xi_{2}} d \boldsymbol{\xi}+\int_{D} J_{\mathbf{x}} A_{12} p_{0, h} \frac{\partial v_{2}}{\partial \xi_{1}} d \boldsymbol{\xi} \\
& +\int_{D}\left(J_{\mathbf{x}} A_{22}-1\right) p_{0, h} \frac{\partial v_{2}}{\partial \xi_{2}} d \boldsymbol{\xi} \\
\leq & \left\|\left(J_{\mathbf{x}} A_{11}-1\right) p_{0, h}\right\|\left\|\frac{\partial v_{1}}{\partial \xi_{1}}\right\|+\left\|J_{\mathbf{x}} A_{21} p_{0, h}\right\|\left\|\frac{\partial v_{1}}{\partial \xi_{2}}\right\|+\left\|J_{\mathbf{x}} A_{12} p_{0, h}\right\|\left\|\frac{\partial v_{2}}{\partial \xi_{1}}\right\| \\
& +\left\|\left(J_{\mathbf{x}} A_{22}-1\right) p_{0, h}\right\|\left\|\frac{\partial v_{2}}{\partial \xi_{2}}\right\| \\
\leq & \left(\left\|\left(J_{\mathbf{x}} A_{11}-1\right) p_{0, h}\right\|^{2}+\left\|J_{\mathbf{x}} A_{21} p_{0, h}\right\|^{2}+\left\|J_{\mathbf{x}} A_{12} p_{0, h}\right\|^{2}\right. \\
& \left.+\left\|\left(J_{\mathbf{x}} A_{22}-1\right) p_{0, h}\right\|^{2}\right)^{\frac{1}{2}}\left(\sum_{i, j=1}^{2}\left\|\frac{\partial v_{i}}{\partial \xi_{j}}\right\|\right)^{\frac{1}{2}} \\
= & \left\|\left(J_{\mathbf{x}} A^{T}-I\right) p_{0, h}\right\|\|\nabla \mathbf{v}\| .
\end{aligned}
$$

We could also proceed as follows:

$$
\begin{aligned}
\mathrm{II}_{3} & =b\left(\mathbf{v}, p_{0, h} ; \mathbf{y}_{0}\right)-b\left(\mathbf{v}, p_{0, h} ; \mathbf{y}\right)=-\int_{D} p_{0, h} \nabla \cdot \mathbf{v} d \boldsymbol{\xi}+\int_{D} p_{0, h} J_{\mathbf{x}}\left(A^{T} \nabla\right) \cdot \mathbf{v} d \boldsymbol{\xi} \\
& =\sum_{i=1}^{d}\left[-\int_{D} p_{0, h}(I \nabla)_{i} v_{i} d \boldsymbol{\xi}+\int_{D} p_{0, h} J_{\mathbf{x}}\left(A^{T} \nabla\right)_{i} v_{i} d \boldsymbol{\xi}\right] \\
& =\sum_{i, j=1}^{d}\left[-\int_{D} p_{0, h} \delta_{i j} \frac{\partial v_{i}}{\partial \xi_{j}} d \boldsymbol{\xi}+\int_{D} p_{0, h} J_{\mathbf{x}}\left(A^{T}\right)_{i j} \frac{\partial v_{i}}{\partial \xi_{j}} d \boldsymbol{\xi}\right] \\
& =\sum_{i, j=1}^{d}\left[\int_{D} p_{0, h}\left(J_{\mathbf{x}}\left(A^{T}\right)_{i j}-\delta_{i j}\right) \frac{\partial v_{i}}{\partial \xi_{j}} d \boldsymbol{\xi}\right]=\int_{D} p_{0, h}\left(J_{\mathbf{x}} A^{T}-I\right): \nabla \mathbf{v} d \boldsymbol{\xi} \\
& \leq \int_{D}\left\|p_{0, h}\left(J_{\mathbf{x}} A^{T}-I\right)\right\|_{F}\|\nabla \mathbf{v}\|_{F} \\
& \leq\left\|p_{0, h}\left(J_{\mathbf{x}} A^{T}-I\right)\right\|_{L^{2}(D)}\|\nabla \mathbf{v}\|_{L^{2}(D)},
\end{aligned}
$$

where $\|\cdot\|_{F}$ denotes the Froebenius norm.

## 4 Time-dependent heat equation with random Robin boundary conditions

## Introduction

In this chapter, we perform an a posteriori error analysis for a time-dependent PDE with random input data, namely the heat equation with random Robin boundary conditions. The analysis is very similar to what has been done in the previous chapters, except that we have to take into account the error due to time discretization. For instance, for the approximation $u \approx u_{0, h \tau}$, where $u_{0, h \tau}$ is a space-time approximation of the deterministic part $u_{0}$ in the expansion of the solution $u$, the a posteriori error estimate is constituted of three parts, see Proposition 4.3.1. Each part controls a different source of error, namely the error due to space discretization, time discretization and uncertainty (truncation in the expansion of $u$ ).

### 4.1 Problem statement

Let $D \subset \mathbb{R}^{d}, d=2,3$, be an open bounded domain with Lipschitz continuous boundary $\partial D$ and let $(\Omega, \mathscr{F}, P)$ be a compete probability space. We consider the following heat problem with random Robin boundary conditions:
find $u:(0, T) \times D \times \Omega \rightarrow \mathbb{R}$ such that a.s. in $\Omega$ the following equations hold

$$
\left\{\begin{array}{rll}
\frac{\partial u(t, \mathbf{x}, \omega)}{\partial t}-\nabla \cdot(k(\mathbf{x}) \nabla u(t, \mathbf{x}, \omega)) & =f(t, \mathbf{x}) &  \tag{4.1}\\
\mathbf{x} \in D, t \in(0, T) \\
u(t, \mathbf{x}, \omega) & =0 & \\
\mathbf{x} \in \Gamma_{D}, t \in(0, T) \\
k(\mathbf{x}) \frac{\partial u(t, \mathbf{x}, \omega)}{\partial \mathbf{n}}+\alpha(\mathbf{x}, \omega) u(t, \mathbf{x}, \omega) & =g(t, \mathbf{x}) & \mathbf{x} \in \Gamma_{R}, t \in(0, T) \\
u(t, \mathbf{x}, \omega) & =\varphi(\mathbf{x}) & \\
\mathbf{x} \in D, t=0
\end{array}\right.
$$

with $\Gamma_{D}$ and $\Gamma_{R}$ the Dirichlet and Robin boundary parts such that $\Gamma_{D} \cup \Gamma_{R}=\partial D$ and $\Gamma_{D} \cap \Gamma_{R}=\varnothing$ and $\mathbf{n}$ is the outward unit normal vector on $\Gamma_{R}$. Notice that the subsequent analysis can be quite easily extended to the cases $f=f(t, \mathbf{x}, \omega), g=g(t, \mathbf{x}, \omega), \varphi=\varphi(\mathbf{x}, \omega)$ or $k=k(\mathbf{x}, \omega)$. From a physical point of view, the Robin boundary conditions for the heat problem are used to model the Newton's law of cooling [123], namely that the rate of change of temperature is
proportional to the temperature difference between the solid surface $\Gamma_{R}$ and its surroundings. Mathematically, this results in imposing a linear combination of Dirichlet (impose the temperature) and Neumann (impose the heat flux) boundary conditions. The parameter $\alpha$ is the heat transfer coefficient and depends on the material, the geometry, the environment, etc. In practise, this coefficient is often determined from experiments and is therefore subject to uncertainty. Another similar problem arises for instance in glaciology, when modelling the motion of glaciers, see for instance [80,104] and references therein. The boundary conditions prescribed on the sliding basal part are indeed affected by uncertainty, for instance due to a lack of knowledge of the shape of the mountain or the difficulty to get measurements of the velocity of the ice on the base of the glacier.

We make the following assumptions on the input data

$$
f \in L^{2}\left(0, T ; L^{2}(D)\right), g \in L^{2}\left(0, T ; L^{2}\left(\Gamma_{R}\right)\right), k \in L^{\infty}\left(D ; \mathbb{R}^{d \times d}\right), \varphi \in L^{2}(D), \alpha(\cdot, \omega) \in L^{\infty}\left(\Gamma_{R}\right) \text { a.s. }
$$

and

$$
\begin{equation*}
\exists k_{\min }>0 \quad \text { such that } \quad \forall \boldsymbol{\xi} \in \mathbb{R}^{d}, \quad k(\mathbf{x}) \boldsymbol{\xi} \cdot \boldsymbol{\xi} \geq k_{\min }|\boldsymbol{\xi}|^{2} \quad \text { a.e. in } D . \tag{4.2}
\end{equation*}
$$

Moreover, we assume that the random field $\alpha$ depends on a finite number of random variables $\left\{Y_{j}\right\}_{j=1}^{L}$, namely

$$
\alpha(\mathbf{x}, \omega)=\alpha(\mathbf{x}, \mathbf{Y}(\omega))=\alpha\left(\mathbf{x}, Y_{1}(\omega), \ldots, Y_{L}(\omega)\right)
$$

Let $\Gamma=\Gamma_{1} \times \ldots \times \Gamma_{L}$, where $\Gamma_{j}=Y_{j}(\Omega)$, and let $\rho: \Gamma \rightarrow \mathbb{R}^{+}$be the joint density function of the random vector $\mathbf{Y}$. Let

$$
V=H_{\Gamma_{D}}^{1}=\left\{v \in H^{1}(D): \quad v=0 \text { on } \Gamma_{D}\right\}
$$

endowed with the norm

$$
\|v\|_{V}:= \begin{cases}|\nu|_{H^{1}(D)}=\|\nabla v\|_{L^{2}(D)} & \text { if } \Gamma_{D} \neq \varnothing \\ \|v\|_{H^{1}(D)}=\sqrt{\|v\|_{L^{2}(D)}^{2}+\|\nabla v\|_{L^{2}(D)}^{2}} & \text { if } \Gamma_{D}=\varnothing\end{cases}
$$

The parametric (pointwise in $\mathbf{y}$ and $t$ ) weak formulation of problem (4.1) reads:
find $u \in L_{\rho}^{2}\left(\Gamma ; L^{2}(0, T ; V) \cap C^{0}\left([0, T] ; L^{2}(D)\right)\right)$ such that

$$
\left\{\begin{align*}
u(0, \mathbf{x}, \mathbf{y}) & =\varphi(\mathbf{x}) \quad \mathbf{x} \in D, \rho \text {-a.e. } \mathbf{y} \in \Gamma  \tag{4.3}\\
\frac{d}{d t} \int_{D} u v+a(u, v ; \mathbf{y}) & =F(v) \quad \forall v \in V, \text { a.e. } t \in(0, T), \rho \text {-a.e. } \mathbf{y} \in \Gamma
\end{align*}\right.
$$

with

$$
\begin{align*}
a(u, v ; \mathbf{y}) & :=\int_{D} k \nabla u \cdot \nabla v+\int_{\Gamma_{R}} \alpha(\mathbf{y}) u v  \tag{4.4}\\
F(v) & :=\int_{D} f v+\int_{\Gamma_{R}} g v . \tag{4.5}
\end{align*}
$$

We can easily show that problem (4.3) is well-posed under the assumption

$$
\begin{equation*}
\alpha(\mathbf{x}, \mathbf{y}) \geq \alpha_{\min }>0 \quad \text { a.e. } \mathbf{x} \in \Gamma_{R}, \rho \text {-a.e. } \mathbf{y} \in \Gamma . \tag{4.6}
\end{equation*}
$$

Indeed, the condition (4.6) ensures the (uniform) coercivity of the bilinear form $a$ defined in (4.4), that is there exists a constant $C_{a}>0$ such that

$$
\begin{equation*}
C_{a}\|v\|_{V}^{2} \leq a(v, v ; \mathbf{y}) \quad \forall v \in V \text { and } \rho \text {-a.e. } \mathbf{y} \in \Gamma \tag{4.7}
\end{equation*}
$$

It is obvious that (4.6) implies (4.7) for the case $\Gamma_{D} \neq \varnothing$, i.e. when $V$ is endowed with the gradient norm, while it can be proved proceeding ab absurdo for the case $\Gamma_{D}=\varnothing$.

Remark 4.1.1. In the case $\Gamma_{D} \neq \varnothing$, the assumption (4.6) can be relaxed since the bilinear form $a$ is also coercive under the condition

$$
\begin{equation*}
\|\alpha(\cdot, \mathbf{y})\|_{L^{\infty}\left(\Gamma_{R}\right)}<\frac{k_{\min }}{C_{T}^{2}\left(1+C_{F}^{2}\right)} \quad \rho-\text { a.e. } \mathbf{y} \in \Gamma \tag{4.8}
\end{equation*}
$$

where $C_{F}$ and $C_{T}$ denote the Friedrich-Poincaré and trace constants in (2.5) and (2.10), respectively. In particular, it is not necessary that $\alpha$ remains positive. Indeed, thanks to (4.2) and using

$$
-\int_{\Gamma_{R}} \alpha v^{2} \leq\|\alpha\|_{L^{\infty}\left(\Gamma_{R}\right)}\|v\|_{L^{2}\left(\Gamma_{R}\right)}^{2} \leq C_{T}^{2}\|\alpha\|_{L^{\infty}\left(\Gamma_{R}\right)}\|\nu\|_{H^{1}(D)}^{2} \leq C_{T}^{2}\left(1+C_{F}^{2}\right)\|\alpha\|_{L^{\infty}\left(\Gamma_{R}\right)}\|\nabla \nu\|_{L^{2}(D)}^{2}
$$

we have

$$
a(v, v ; \mathbf{y})=\int_{D} k|\nabla v|^{2}+\int_{\Gamma_{R}} \alpha v^{2} \geq\left(k_{\min }-C_{T}^{2}\left(1+C_{F}^{2}\right)\|\alpha\|_{L^{\infty}\left(\Gamma_{R}\right)}\right)\|\nabla v\|_{L^{2}(D)}^{2}
$$

for any $v \in V$ and $\rho$-a.e. in $\Gamma$. The coercivity constant $C_{a}>0$ is then given by

$$
C_{a}= \begin{cases}k_{\min } & \text { if (4.6) holds } \\ k_{\min }-\|\alpha\|_{L^{\infty}\left(\Gamma_{R}\right)} C_{T}^{2}\left(1+C_{F}^{2}\right) & \text { if (4.8) holds }\end{cases}
$$

## Specific form of $\alpha$

We assume that the random coefficient $\alpha$, which appears in the Robin boundary condition, depends in an affine way on the random variables, namely that it can be written

$$
\alpha(\mathbf{x}, \mathbf{Y}(\omega))=\alpha_{0}(\mathbf{x})+\varepsilon \sum_{j=1}^{L} \alpha_{j}(\mathbf{x}) Y_{j}(\omega)
$$

where $\left\{Y_{j}\right\}_{j=1}^{L}$ are independent random variables with zero mean and unit variance.
Example 4.1.2. Let $D=(0,1)^{2}$ with boundary $\Gamma_{D}$ and $\Gamma_{R}=\Gamma_{R_{1}} \cup \Gamma_{R_{2}} \cup \Gamma_{R_{3}}$ as shown in Figure 4.1.


Figure 4.1: Geometry with label for each part of the boundary.

We take then $\alpha(\mathbf{x}, \mathbf{Y}(\omega))=\alpha_{0}(\mathbf{x})+\varepsilon \sum_{j=1}^{3} \alpha_{j}(\mathbf{x}) Y_{j}(\omega)$ with

$$
\alpha_{0}=\left\{\begin{array}{ll}
\alpha_{0,1} & \text { if } \mathbf{x} \in \Gamma_{R_{1}} \\
\alpha_{0,2} & \text { if } \mathbf{x} \in \Gamma_{R_{2}} \\
\alpha_{0,3} & \text { if } \mathbf{x} \in \Gamma_{R_{3}}
\end{array} \quad, \quad \alpha_{j}=\left\{\begin{array}{ll}
a_{j} & \text { if } \mathbf{x} \in \Gamma_{R_{j}} \\
0 & \text { if } \mathbf{x} \in \Gamma_{R} \backslash \Gamma_{R_{j}}
\end{array}, \quad g= \begin{cases}g_{1} & \text { if } \mathbf{x} \in \Gamma_{R_{1}} \\
g_{2} & \text { if } \mathbf{x} \in \Gamma_{R_{2}} \\
g_{3} & \text { if } \mathbf{x} \in \Gamma_{R_{3}}\end{cases}\right.\right.
$$

and $\alpha_{0, j}, a_{j} \in L^{\infty}\left(\Gamma_{R_{j}}\right), j=1,2,3$, such that (4.6) holds, i.e. $\alpha \geq \alpha_{\min }>0$. For instance, in the case $\Gamma=[-1,1]^{3}$, it is then required that $\varepsilon\left|a_{j}\right|<\alpha_{0, j}$ for $j=1,2,3$.

## Methodology

As in the previous chapters, we use a perturbation technique expanding the (random) solution $u$ with respect to $\varepsilon$ as:

$$
u(t, \mathbf{x}, \mathbf{Y}(\omega))=u_{0}(t, \mathbf{x})+\varepsilon u_{1}(t, \mathbf{x}, \mathbf{Y}(\omega))+\varepsilon^{2} u_{2}(t, \mathbf{x}, \mathbf{Y}(\omega))+\ldots
$$

The problem for the first term $u_{0}$ in the expansion simply reads:
find $u_{0}:(0, T) \times D \rightarrow \mathbb{R}$ such that

$$
\left\{\begin{array}{rlll}
\frac{\partial u_{0}(t, \mathbf{x})}{\partial t}-\nabla \cdot\left(k(\mathbf{x}) \nabla u_{0}(t, \mathbf{x})\right) & =f(t, \mathbf{x}) & \mathbf{x} \in D, t \in(0, T)  \tag{4.9}\\
u_{0}(t, \mathbf{x}) & =0 & \mathbf{x} \in \Gamma_{D}, t \in(0, T) \\
k(\mathbf{x}) \frac{\partial u_{0}(t, \mathbf{x})}{\partial \mathbf{n}}+\alpha_{0}(\mathbf{x}) u_{0}(t, \mathbf{x}) & =g(t, \mathbf{x}) & \mathbf{x} \in \Gamma_{R}, t \in(0, T) \\
u_{0}(t, \mathbf{x}) & =\varphi(\mathbf{x}) & & \mathbf{x} \in D, t=0
\end{array}\right.
$$

whose weak formulation can be written:
find $u_{0} \in L^{2}(0, T ; V) \cap C^{0}\left([0, T] ; L^{2}(D)\right)$ such that

$$
\left\{\begin{align*}
u_{0}(0, \mathbf{x}) & =\varphi(\mathbf{x}) & \mathbf{x} \in D  \tag{4.10}\\
\frac{d}{d t} \int_{D} u_{0} v+\int_{D} k \nabla u_{0} \cdot \nabla v+\int_{\Gamma_{R}} \alpha_{0} u_{0} v & =\int_{D} f v+\int_{\Gamma_{R}} g v & \forall v \in V, \text { a.e. } t \in(0, T) .
\end{align*}\right.
$$

Notice that problem (4.10) is nothing else than problem (4.3) with $\mathbf{y}=\mathbb{E}[\mathbf{Y}]=\mathbf{0}$. Writing $u_{1}(t, \mathbf{x}, \mathbf{Y}(\omega))=\sum_{j=1}^{L} U_{j}(t, \mathbf{x}) Y_{j}(\omega)$, the second term in the expansion can be obtained by solving the $L$ problems:
find $U_{j}:(0, T) \times D \rightarrow \mathbb{R}$ such that

$$
\left\{\begin{align*}
\frac{\partial U_{j}(t, \mathbf{x})}{\partial t}-\nabla \cdot\left(k(\mathbf{x}) \nabla U_{j}(t, \mathbf{x})\right) & =0 & & \mathbf{x} \in D, t \in(0, T)  \tag{4.11}\\
U_{j}(t, \mathbf{x}) & =0 & & \mathbf{x} \in \Gamma_{D}, t \in(0, T) \\
k(\mathbf{x}) \frac{\partial U_{j}(t, \mathbf{x})}{\partial \mathbf{n}}+\alpha_{0}(\mathbf{x}) U_{j}(t, \mathbf{x}) & =-\alpha_{j}(\mathbf{x}) u_{0}(\mathbf{x}) & & \mathbf{x} \in \Gamma_{R}, t \in(0, T) \\
U_{j}(t, \mathbf{x}) & =0 & & \mathbf{x} \in D, t=0
\end{align*}\right.
$$

### 4.2 Numerical approximation

We assume from now on that $f \in C^{0}\left([0, T] ; L^{2}(D)\right), g \in C^{0}\left([0, T] ; L^{2}\left(\Gamma_{R}\right)\right)$ and $\varphi \in C^{0}(\bar{D})$.
We approximate the solution $u_{0}$ of problem (4.10) using the (implicit) Backward Euler scheme in time and $\left(\mathbb{P}_{k}\right)$ finite elements in space. For any $\tau>0$, let $0=t_{0}<t_{1}<\ldots<t_{M}=T$ be a discretization of the time interval $[0, T]$ into $M$ subintervals $I_{n}=\left[t_{n-1}, t_{n}\right]$ of length $\tau_{n}=$ $t_{n}-t_{n-1} \leq \tau, n=1, \ldots, M$. Moreover, for any $h>0$, let $\mathscr{T}_{h}$ be a shape regular (in the sense of [49]) partition of $D$ into $d$-simplices $K$ of diameter $h_{K} \leq h$ and let

$$
V_{h}=\left\{v \in C^{0}(\bar{D}): \quad v_{\mid K} \in \mathbb{P}_{k}, \forall K \in \mathscr{T}_{h}\right\} \cap V
$$

be the subspace of $V$ constituted of continuous, piecewise polynomial functions on $\mathscr{T}_{h}$.
Remark 4.2.1. Notice that a different mesh could be used for each time step, see e.g. [103], in which case we would write $\mathscr{T}_{h}^{n}$ and $V_{h}^{n}$ the mesh and FE space at time $t_{n}$. This functionality would be needed for instance when using adaptive algorithms, to allow the spatial meshes to vary in time. The introduction of an (interpolant) operator between two successive meshes is then required.

The fully discretized problem reads:

1. Initialization: $u_{0, h}^{0}=r_{h} \varphi$
2. For $n=1, \ldots, M$ : find $u_{0, h}^{n} \in V_{h}$ such that:

$$
\begin{equation*}
\int_{D} \frac{u_{0, h}^{n}-u_{0, h}^{n-1}}{\tau_{n}} v_{h}+\int_{D} k \nabla u_{0, h}^{n} \cdot \nabla v_{h}+\int_{\Gamma_{R}} \alpha_{0} u_{0, h}^{n} v_{h}=\int_{D} f^{n} v_{h}+\int_{\Gamma_{R}} g^{n} v_{h} \quad \forall v_{h} \in V_{h} \tag{4.12}
\end{equation*}
$$

where $f^{n}=f\left(\cdot, t_{n}\right)$ and $g^{n}=g\left(\cdot, t_{n}\right)$. Finally, we define the global approximation $u_{0, h \tau}$, linear
on each subinterval $I_{n}$, by

$$
\begin{equation*}
u_{0, h \tau}(t, \mathbf{x}):=\frac{t-t_{n-1}}{\tau_{n}} u_{0, h}^{n}(\mathbf{x})+\frac{t_{n}-t}{\tau_{n}} u_{0, h}^{n-1}(\mathbf{x}) \quad \text { for } t \in\left[t_{n-1}, t_{n}\right], n=1, \ldots, M . \tag{4.13}
\end{equation*}
$$

### 4.3 A posteriori error analysis

For ease of notation, we introduce the element and edge residuals $R$ and $J$ defined on each element $K$ and each edge $e$ by, respectively,

$$
\begin{equation*}
\left.R\left(u_{0, h \tau}\right)\right|_{K}:=f-\frac{\partial u_{0, h \tau}}{\partial t}+\nabla \cdot\left(k \nabla u_{0, h \tau}\right) \tag{4.14}
\end{equation*}
$$

and

$$
\left.J\left(u_{0, h \tau}\right)\right|_{e}:= \begin{cases}\frac{1}{2}\left[k \nabla u_{0, h \tau} \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}} & \text { if } e \subset D  \tag{4.15}\\ g-\alpha_{0} u_{0, h \tau}-k \nabla u_{0, h \tau} \cdot \mathbf{n}_{e} & \text { if } e \subset \Gamma_{R} \\ 0 & \text { if } e \subset \Gamma_{D}\end{cases}
$$

We have denoted by $[\cdot]_{\mathbf{n}_{e}}$ the jump across an interior edge $e$, defined by

$$
[\varphi]_{\mathbf{n}_{e}}(\mathbf{x}):=\lim _{t \rightarrow 0^{+}}\left(\varphi\left(\mathbf{x}+t \mathbf{n}_{e}\right)-\varphi\left(\mathbf{x}-t \mathbf{n}_{e}\right)\right)
$$

Here, $\mathbf{n}_{e}$ is the outer unit normal vector to the edge $e$ if $e \subset \Gamma_{R}$ while for interior edges $e \subset D$, it is a unit normal vector to $e$ of arbitrary (but fixed) direction. Notice that the choice of direction is irrelevant since quantities of the type $\left[\nabla \varphi \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}$ is not affected by this choice, while $[\varphi]_{\mathbf{n}_{e}}$ is.

We have now introduced all the ingredients necessary to derive our a posteriori error estimate for the error $e:=u-u_{0, h \tau}$ given in the following proposition.

Proposition 4.3.1. Let $u$ be the weak solution of problem (4.1) and let $u_{0, h \tau}$ be defined in (4.13). Then there exists a constant $C>0$ depending only on the trace constant and the mesh aspect ratio such that

$$
\begin{align*}
\mathbb{E}\left[\left\|\left(u-u_{0, h \tau}\right)(T)\right\|_{L^{2}(D)}^{2}\right]+ & C_{a} \int_{0}^{T} \mathbb{E}\left[\left\|u-u_{0, h \tau}\right\|_{V}^{2}\right] d t \leq \\
& \left\|\varphi-r_{h} \varphi\right\|_{L^{2}(D)}^{2}+\frac{C}{C_{a}} \sum_{n=1}^{M} \sum_{K \in \mathscr{T}_{h}}\left[\int_{t_{n-1}}^{t_{n}}\left(\left(\eta_{K}^{n}\right)^{2}+\left(\gamma_{K}^{n}\right)^{2}+\left(\theta_{K}^{n}\right)^{2}\right) d t\right], \tag{4.16}
\end{align*}
$$

where $C_{a}$ is the constant in (4.7) and

$$
\begin{align*}
& \left(\eta_{K}^{n}\right)^{2}:=h_{K}^{2}\left\|R\left(u_{0, h \tau}\right)\right\|_{L^{2}(K)}^{2}+\sum_{e \subset \partial K} h_{e}\left\|J\left(u_{0, h \tau}\right)\right\|_{L^{2}(e)}^{2}  \tag{4.17}\\
& \left(\gamma_{K}^{n}\right)^{2}:=\left\|k \nabla\left(u_{0, h \tau}-u_{0, h}^{n}\right)\right\|_{L^{2}(K)}^{2}+\left\|f-f^{n}\right\|_{L^{2}(K)}^{2}+\sum_{e \subset \partial K \cap \Gamma_{R}}\left\|g-g^{n}-\alpha_{0}\left(u_{0, h \tau}-u_{0, h}^{n}\right)\right\|_{L^{2}(e)}^{2} \\
& \left(\theta_{K}^{n}\right)^{2}:=\varepsilon^{2} \sum_{j=1}^{L}\left\|\alpha_{j} u_{0, h \tau}\right\|_{L^{2}\left(\partial K \cap \Gamma_{R}\right)}^{2} . \tag{4.18}
\end{align*}
$$

Proof. Let us write $e=u-u_{0, h \tau}$. In what follows, all equations are valid for a.e. $t$ and a.s. in $\Omega$ without necessarily mentioning it. Moreover, $C$ will denote a generic constant, whose value might change from one occurrence to another, that depends only on the interpolation constants in (1.26), (1.27) and (1.28), the trace constant in (2.10) and, if $\Gamma_{D} \neq \varnothing$, the FriedrichPoincaré constant in (2.5). Thanks to equations (4.12) and (4.13), we have for each $v_{h} \in V_{h}$ and each $n \in\{1, \ldots, M\}$

$$
\begin{align*}
\int_{D} \frac{\partial u_{0, h \tau}}{\partial t} v_{h}+\int_{D} k \nabla u_{0, h \tau} \cdot \nabla v_{h}+\int_{\Gamma_{R}} \alpha_{0} u_{0, h \tau} v_{h}= & \int_{D} f v_{h}+\int_{\Gamma_{R}} g v_{h}+\int_{D} k \nabla\left(u_{0, h \tau}-u_{0, h}^{n}\right) \cdot \nabla v_{h} \\
& +\int_{\Gamma_{R}} \alpha_{0}\left(u_{0, h \tau}-u_{0, h}^{n}\right) v_{h}+\int_{D}\left(f^{n}-f\right) v_{h} \\
& +\int_{\Gamma_{R}}\left(g^{n}-g\right) v_{h} \tag{4.20}
\end{align*}
$$

using the fact that $\frac{\partial u_{0, h \tau}}{\partial t}=\frac{u_{0, h}^{n}-u_{0, h}^{n-1}}{\tau_{n}}$ on each time subinterval $I_{n}, n=1, \ldots, M$. Thanks to the coercivity of $a$, see (4.7), we have

$$
C_{a}\|e\|_{V}^{2} \leq \int_{D} k|\nabla e|^{2}+\int_{\Gamma_{R}} \alpha e^{2}
$$

We now let $n$ be any value in $\{1, \ldots, M\}$. Then, for all $v \in V$ we have

$$
\begin{aligned}
\frac{d}{d t} \int_{D} e v+\int_{D} k \nabla e \cdot \nabla v+\int_{\Gamma_{R}} \alpha e v= & \int_{D} f v+\int_{\Gamma_{R}} g v-\int_{D} \frac{\partial u_{0, h \tau}}{\partial t} v-\int_{D} k \nabla u_{0, h \tau} \cdot \nabla v \\
& -\int_{\Gamma_{R}} \alpha_{0} u_{0, h \tau} v-\int_{\Gamma_{R}}(\alpha-20) \\
& \left.\int_{D} f\left(v-I_{h}\right) u_{0, h \tau} v\right)+\int_{\Gamma_{R}} g\left(v-I_{h} v\right)-\int_{D} \frac{\partial u_{0, h \tau}}{\partial t}\left(v-I_{h} v\right) \\
& -\int_{D} k \nabla u_{0, h \tau} \cdot \nabla\left(v-I_{h} v\right)-\int_{\Gamma_{R}} \alpha_{0} u_{0, h \tau}\left(v-I_{h} v\right) \\
& -\int_{D} k \nabla\left(u_{0, h \tau}-u_{0, h}^{n}\right) \cdot \nabla I_{h} v-\int_{\Gamma_{R}} \alpha_{0}\left(u_{0, h \tau}-u_{0, h}^{n}\right) I_{h} v \\
& -\int_{D}\left(f^{n}-f\right) I_{h} v-\int_{\Gamma_{R}}\left(g^{n}-g\right) I_{h} v-\int_{\Gamma_{R}}\left(\alpha-\alpha_{0}\right) u_{0, h \tau} v,
\end{aligned}
$$

where $I_{h}$ denotes the Clément interpolant of $v$. Taking then $v=e(t, \cdot, \mathbf{Y}(\omega))$ a.e. $t \in I_{n}$ and a.s.
in $\Omega$ in the last inequality, we get

$$
\begin{equation*}
\frac{1}{2} \frac{d}{d t}\|e\|_{L^{2}(D)}^{2}+C_{a}\|e\|_{V}^{2} \leq \mathrm{I}+\mathrm{II}+\mathrm{III} \tag{4.21}
\end{equation*}
$$

with

$$
\begin{aligned}
\mathrm{I} & :=\sum_{K \in \mathscr{T}_{h}}\left\{\int_{K} R\left(u_{0, h \tau}\right)\left(e-I_{h} e\right)+\int_{\partial K} J\left(u_{0, h \tau}\right)\left(e-I_{h} e\right)\right\} \\
\mathrm{II} & :=-\int_{D} k \nabla\left(u_{0, h \tau}-u_{0, h}^{n}\right) \cdot \nabla I_{h} e-\int_{\Gamma_{R}} \alpha_{0}\left(u_{0, h \tau}-u_{0, h}^{n}\right) I_{h} e+\int_{D}\left(f-f^{n}\right) I_{h} e+\int_{\Gamma_{R}}\left(g-g^{n}\right) I_{h} e \\
\text { III } & :=-\int_{\Gamma_{R}}\left(\alpha-\alpha_{0}\right) u_{0, h \tau} e
\end{aligned}
$$

and $R$ and $J$ defined in (4.14) and (4.15), respectively. Notice that the terms I, II and III control the error due to space discretization, time discretization and truncation in the expansion of $u$, respectively. We now bound each of these terms separately.
bound for I: recalling the definition of $\eta_{K}^{n}$ in (4.17), we obtain using a standard procedure the bound

$$
\begin{equation*}
\mathrm{I} \leq C_{1}\left(\sum_{K \in \mathscr{T}_{h}}\left(\eta_{k}^{n}\right)^{2}\right)^{\frac{1}{2}}\|e\|_{V} \tag{4.22}
\end{equation*}
$$

where $C_{1}$ is a positive constant that depends only on the interpolation constants in (1.26) and (1.28).
bound for II: thanks to the triangle inequality, the interpolation error bounds (1.27) and (1.28) and the trace inequality (2.10), the following inequalities hold true

$$
\begin{aligned}
\left\|\nabla I_{h} e\right\|_{L^{2}(K)} & \leq\|\nabla e\|_{L^{2}(K)}+\left\|\nabla\left(e-I_{h} e\right)\right\|_{L^{2}(K)} \leq C|e|_{H^{1}(N(K))} \\
\left\|I_{h} e\right\|_{L^{2}(K)} & \leq\|e\|_{L^{2}(K)}+\left\|e-I_{h} e\right\|_{L^{2}(K)} \leq C\left(1+h_{K}\right)\|e\|_{H^{1}(N(K))} \leq C\|e\|_{H^{1}(N(K))} \\
\left\|I_{h} e\right\|_{L^{2}\left(\Gamma_{R}\right)} & \leq C_{T}\left\|I_{h} e\right\|_{H^{1}(D)}=C_{T}\left[\sum_{K \in \mathscr{T}_{h}}\left(\left\|I_{h} e\right\|_{L^{2}(K)}^{2}+\left\|\nabla I_{h} e\right\|_{L^{2}(K)}^{2}\right)\right]^{\frac{1}{2}} \leq C\|e\|_{V} .
\end{aligned}
$$

Therefore, regrouping the integrals over the boundary $\Gamma_{R}$, we obtain the bound

$$
\begin{align*}
\mathrm{II} \leq & C_{2}\left[\sum_{K \in \mathscr{T}_{h}}\left(\left\|k \nabla\left(u_{0, h \tau}-u_{0, h}^{n}\right)\right\|_{L^{2}(K)}^{2}+\left\|f-f^{n}\right\|_{L^{2}(K)}^{2}\right)\right. \\
& \left.+\sum_{e \subset \Gamma_{R}}\left\|g-g^{n}-\alpha_{0}\left(u_{0, h \tau}-u_{0, h}^{n}\right)\right\|_{L^{2}(e)}^{2}\right]^{\frac{1}{2}}\|e\|_{V} \\
= & C_{2}\left(\sum_{K \in \mathscr{T}_{h}}\left(\gamma_{K}^{n}\right)^{2}\right)^{\frac{1}{2}}\|e\|_{V} \tag{4.23}
\end{align*}
$$

with $\gamma_{K}^{n}$ given in (4.18) and where $C_{2}$ is a positive constant that depends only on the constants in (1.27), (1.28) and (2.10). It additionally depends on the Friedrich-Poincaré constant in (2.5)
in the case $\Gamma_{D} \neq \varnothing$.
bound for III: for the last term, we easily get

$$
\begin{equation*}
\mathrm{III} \leq\left\|\left(\alpha-\alpha_{0}\right) u_{0, h \tau}\right\|_{L^{2}\left(\Gamma_{R}\right)}\|e\|_{L^{2}\left(\Gamma_{R}\right)} \leq C_{3}\left\|\left(\alpha-\alpha_{0}\right) u_{0, h \tau}\right\|_{L^{2}\left(\Gamma_{R}\right)}\|e\|_{V} \tag{4.24}
\end{equation*}
$$

where $C_{3}=C_{T}$ if $\Gamma_{D}=\varnothing$ and $C_{3}=C_{T} \sqrt{1+C_{F}^{2}}$ otherwise, with $C_{T}$ and $C_{F}$ given in (2.10) and (2.5), respectively.

Using the bounds (4.22), (4.23) and (4.24) in (4.21) yields

$$
\begin{aligned}
\frac{1}{2} \frac{d}{d t}\|e\|_{L^{2}(D)}^{2}+C_{a}\|e\|_{V}^{2} & \leq C\left[\sum_{K \in \mathscr{T}_{h}}\left(\left(\eta_{K}^{n}\right)^{2}+\left(\gamma_{K}^{n}\right)^{2}\right)+\left\|\left(\alpha-\alpha_{0}\right) u_{0, h \tau}\right\|_{L^{2}\left(\Gamma_{R}\right)}^{2}\right]^{\frac{1}{2}}\|e\|_{V} \\
& \leq \frac{C}{2 C_{a}}\left[\sum_{K \in \mathscr{T}_{h}}\left(\left(\eta_{K}^{n}\right)^{2}+\left(\gamma_{K}^{n}\right)^{2}\right)+\left\|\left(\alpha-\alpha_{0}\right) u_{0, h \tau}\right\|_{L^{2}\left(\Gamma_{R}\right)}^{2}\right]+\frac{C_{a}}{2}\|e\|_{V}^{2}
\end{aligned}
$$

and thus, splitting the integral of the last term of the right-hand side over the elements $K$ we get

$$
\frac{d}{d t}\|e\|_{L^{2}(D)}^{2}+C_{a}\|e\|_{V}^{2} \leq \frac{C}{C_{a}} \sum_{K \in \mathscr{T}_{h}}\left[\left(\eta_{K}^{n}\right)^{2}+\left(\gamma_{K}^{n}\right)^{2}+\left\|\left(\alpha-\alpha_{0}\right) u_{0, h \tau}\right\|_{L^{2}\left(\partial K \cap \Gamma_{R}\right)}^{2}\right]
$$

To conclude the proof, we integrate the last inequality over the time subinterval $I_{n}$, we sum then over $n$ ranging from 1 to $M$ and finally, we take the expected value on both sides recalling that $\mathbb{E}\left[Y_{i} Y_{j}\right]=\delta_{i j}$.

### 4.4 Numerical results

We give here two numerical examples to test the a posteriori error estimate derived in Section 4.3, see Proposition 4.3.1. We use $\mathbb{P}_{1}$ finite elements for the physical space approximation. In both examples, we set $k=I$ and we consider the case $\Gamma_{D} \neq \varnothing$. Therefore, the error $e=u-u_{0, h \tau}$ with $u_{0, h \tau}$ defined in (4.13) is measured with the norm

$$
\begin{equation*}
\operatorname{err}:=\mathbb{E}\left[\int_{0}^{T}\|\nabla e(t, \cdot, \cdot)\|_{L^{2}(D)}^{2} d t\right]^{\frac{1}{2}}=\left(\int_{\Omega} \int_{0}^{T} \int_{D}|\nabla e(t, \mathbf{x}, \omega)|^{2} d \mathbf{x} d t d P(\omega)\right)^{\frac{1}{2}} \tag{4.25}
\end{equation*}
$$

Similarly to [103], we define then the error estimator

$$
\begin{equation*}
\text { est }:=\left(w_{\eta}^{2} \eta^{2}+w_{\gamma}^{2} \gamma^{2}+w_{\theta}^{2} \theta^{2}\right)^{\frac{1}{2}} \tag{4.26}
\end{equation*}
$$

with weights $w_{\eta}, w_{\gamma}$ and $w_{\theta}$ to be defined and

$$
\eta^{2}=\sum_{n=1}^{M} \sum_{K \in \mathscr{T}_{h}} \int_{t_{n-1}}^{t_{n}}\left(\eta_{K}^{n}\right)^{2} d t, \quad \gamma^{2}=\sum_{n=1}^{M} \sum_{K \in \mathscr{T}_{h}} \int_{t_{n-1}}^{t_{n}}\left(\gamma_{K}^{n}\right)^{2} d t, \quad \theta^{2}=\sum_{n=1}^{M} \sum_{K \in \mathscr{T}_{h}} \int_{t_{n-1}}^{t_{n}}\left(\theta_{K}^{n}\right)^{2} d t
$$

where $\eta_{K}^{n}, \gamma_{K}^{n}$ and $\theta_{K}^{n}$ are defined in (4.17), (4.18) and (4.19), respectively. Notice that $\eta$ controls the space discretization, $\gamma$ the time discretization and $\theta$ the truncation in the expansion of $u$ with respect to $\varepsilon$.

Let $D=(0,1)^{2}$ with boundary $\partial D=\Gamma_{D} \cup \Gamma_{R}$ as in Figure 4.1, let $T=1$ and let $Y_{j}, j=1,2,3$, be independent uniform random variables in $[-\sqrt{3}, \sqrt{3}]$. For the first case ${ }^{1}$, we consider

$$
\begin{equation*}
u_{0}\left(t, x_{1}, x_{2}\right)=\sin \left(\frac{10 \pi t}{2}\right) \sin \left(\frac{\pi x_{1}}{2}\right) \sin \left(\frac{\pi x_{2}}{2}\right) \quad \text { and } \quad \alpha(\mathbf{x}, \mathbf{Y}(\omega))=\alpha_{0}(\mathbf{x})+\varepsilon \sum_{j=1}^{3} \alpha_{j}(\mathbf{x}) Y_{j}(\omega) \tag{4.27}
\end{equation*}
$$

with $\alpha_{0}(\mathbf{x})=1$ and $\alpha_{j}(\mathbf{x})=\chi_{\Gamma_{R_{j}}}(\mathbf{x}), \chi$ being the indicator function. We plug then $u_{0}$ and $\alpha_{0}$ in (4.9) and compute the corresponding (deterministic) right-hand side $f$, boundary data $g$ and initial condition $\varphi$. For the second case, using the same notation as in Example 4.1.2, we choose

$$
f=\sin \left(2 \pi x_{1}\right) t, \varphi=0, g_{1}=g_{2}=g_{3}=0, \alpha_{0}=\left\{\begin{align*}
1 & \text { if } \mathbf{x} \in \Gamma_{R_{1}}  \tag{4.28}\\
2 & \text { if } \mathbf{x} \in \Gamma_{R_{2}} \\
1.4 & \text { if } \mathbf{x} \in \Gamma_{R_{3}}
\end{align*} \text { and } a_{1}=0.9, a_{2}=1.2, a_{3}=1 .\right.
$$

We use a Delaunay triangulation with $N$ equispaced vertices on each side of $D$ for the space discretization and a uniform time step $\tau$ for the time discretization.

## Deterministic case

We start considering the case $\varepsilon=0$. For the first problem, the error is computed with respect to the exact solution $u_{0}$ in (4.27) while for the second case (4.28), we use a reference solution computed with $N_{r e f}=80$ and $\tau_{r e f}=2^{-7}$. Moreover, the constants $w_{\eta}$ and $w_{\gamma}$ in (4.26) have been tuned considering two test problems with exact solutions for (4.9), namely $u_{0}=$ $\sin \left(\pi x_{1} / 2\right) \sin \left(\pi x_{2} / 2\right)$ (mainly space error) and $u_{0}=\sin (\pi t / 2)$ (mainly time error), leading to $w_{\eta}=1 / 5$ and $w_{\gamma}=1 / 13$.

We give in Table 4.1 the results we get for the first case described in (4.27), considering various meshes with $N=10,20,30,40$ and various time steps $\tau=2^{-4}, 2^{-5}, 2^{-6}, 2^{-7}$. The results obtained when computing the error with respect to a reference solution obtained with $N_{\text {ref }}=80$ and $\tau_{r e f}=2^{-9}$ are also provided, for comparison with the random case below where such reference discretization parameters are used. The results for the case (4.28) with $N=10,20,40$ and $\tau=2^{-4}, 2^{-5}, 2^{-6}$ are provided in Table 4.2.

We see that for the first case (4.27), the error due to time discretization dominates the one due to the space approximation. The contrary holds for the second case (4.28) where the FE error is dominant. In both cases, the error estimator that contains the weights $w_{\eta}$ and $w_{\gamma}$ provides

[^5]| $N$ | $\tau$ | $e r r$ | $w_{\eta} \eta$ | $w_{\gamma} \gamma$ | $e s t$ | e.i. | $e r r$ ref | e.i. ref |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 10 | $2^{-4}$ | $3.0665 \mathrm{e}-1$ | $1.0463 \mathrm{e}-1$ | $3.0351 \mathrm{e}-1$ | $3.2104 \mathrm{e}-1$ | 1.0469 | $2.9537 \mathrm{e}-1$ | 1.0869 |
| 10 | $2^{-5}$ | $1.6313 \mathrm{e}-1$ | $7.5454 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $1.7373 \mathrm{e}-1$ | 1.0650 | $1.5339 \mathrm{e}-1$ | 1.1326 |
| 10 | $2^{-6}$ | $9.2626 \mathrm{e}-2$ | $6.0711 \mathrm{e}-2$ | $7.8871 \mathrm{e}-2$ | $9.9531 \mathrm{e}-2$ | 1.0745 | $8.415 \mathrm{e}-2$ | 1.1827 |
| 10 | $2^{-7}$ | $6.2591 \mathrm{e}-2$ | $5.4026 \mathrm{e}-2$ | $3.9521 \mathrm{e}-2$ | $6.6938 \mathrm{e}-2$ | 1.0695 | $5.6666 \mathrm{e}-2$ | 1.1813 |
|  |  |  |  |  |  |  |  |  |
| 20 | $2^{-4}$ | $3.0436 \mathrm{e}-1$ | $5.2280 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $3.0798 \mathrm{e}-1$ | 1.0119 | $2.9298 \mathrm{e}-1$ | 1.0512 |
| 20 | $2^{-5}$ | $1.5801 \mathrm{e}-1$ | $3.7632 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $1.6095 \mathrm{e}-1$ | 1.0186 | $1.4795 \mathrm{e}-1$ | 1.0879 |
| 20 | $2^{-6}$ | $8.2734 \mathrm{e}-2$ | $3.0209 \mathrm{e}-2$ | $7.8869 \mathrm{e}-2$ | $8.4457 \mathrm{e}-2$ | 1.0208 | $7.3232 \mathrm{e}-2$ | 1.1533 |
| 20 | $2^{-7}$ | $4.6377 \mathrm{e}-2$ | $2.6833 \mathrm{e}-2$ | $3.9520 \mathrm{e}-2$ | $4.7768 \mathrm{e}-2$ | 1.0300 | $3.8267 \mathrm{e}-2$ | 1.2483 |
|  |  |  |  |  |  |  |  |  |
| 40 | $2^{-4}$ | $3.0383 \mathrm{e}-1$ | $2.5771 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $3.0460 \mathrm{e}-1$ | 1.0025 | $2.9242 \mathrm{e}-1$ | 1.0416 |
| 40 | $2^{-5}$ | $1.5679 \mathrm{e}-1$ | $1.8571 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $1.5759 \mathrm{e}-1$ | 1.0051 | $1.4665 \mathrm{e}-1$ | 1.0746 |
| 40 | $2^{-6}$ | $8.0254 \mathrm{e}-2$ | $1.4930 \mathrm{e}-2$ | $7.8869 \mathrm{e}-2$ | $8.0270 \mathrm{e}-2$ | 1.0002 | $7.0421 \mathrm{e}-2$ | 1.1399 |
| 40 | $2^{-7}$ | $4.1706 \mathrm{e}-2$ | $1.3278 \mathrm{e}-2$ | $3.9520 \mathrm{e}-2$ | $4.1691 \mathrm{e}-2$ | 0.9996 | $3.2461 \mathrm{e}-2$ | 1.2843 |
|  |  |  |  |  |  |  |  |  |
| 80 | $2^{-4}$ | $3.0369 \mathrm{e}-1$ | $1.2951 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $3.0378 \mathrm{e}-1$ | 1.0003 | $2.9222 \mathrm{e}-1$ | 1.0396 |
| 80 | $2^{-5}$ | $1.5648 \mathrm{e}-1$ | $9.3271 \mathrm{e}-3$ | $1.5649 \mathrm{e}-1$ | $1.5677 \mathrm{e}-1$ | 1.0019 | $1.4616 \mathrm{e}-1$ | 1.0726 |
| 80 | $2^{-6}$ | $7.9614 \mathrm{e}-2$ | $7.4936 \mathrm{e}-3$ | $7.8869 \mathrm{e}-2$ | $7.9224 \mathrm{e}-2$ | 0.9951 | $6.9327 \mathrm{e}-2$ | 1.1428 |
| 80 | $2^{-7}$ | $4.0438 \mathrm{e}-2$ | $6.6615 \mathrm{e}-3$ | $3.9520 \mathrm{e}-2$ | $4.0077 \mathrm{e}-2$ | 0.9911 | $2.9953 \mathrm{e}-2$ | 1.3380 |

Table 4.1: Error, estimators and effectivity index for the first case (4.27) with $\varepsilon=0$.

| $N$ | $\tau$ | $e r r$ | $w_{\eta} \eta$ | $w_{\gamma} \gamma$ | est | e.i. |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 10 | $2^{-4}$ | $9.8673 \mathrm{e}-3$ | $1.0500 \mathrm{e}-2$ | $2.4507 \mathrm{e}-3$ | $1.0782 \mathrm{e}-2$ | 1.0928 |
| 10 | $2^{-5}$ | $9.8634 \mathrm{e}-3$ | $1.0491 \mathrm{e}-2$ | $1.2254 \mathrm{e}-3$ | $1.0562 \mathrm{e}-2$ | 1.0708 |
| 10 | $2^{-6}$ | $9.8624 \mathrm{e}-3$ | $1.0488 \mathrm{e}-2$ | $6.1275 \mathrm{e}-4$ | $1.0506 \mathrm{e}-2$ | 1.0653 |
|  |  |  |  |  |  |  |
| 20 | $2^{-4}$ | $5.1306 \mathrm{e}-3$ | $5.2838 \mathrm{e}-3$ | $2.4512 \mathrm{e}-3$ | $5.8247 \mathrm{e}-3$ | 1.1353 |
| 20 | $2^{-5}$ | $5.1233 \mathrm{e}-3$ | $5.2790 \mathrm{e}-3$ | $1.2257 \mathrm{e}-3$ | $5.4194 \mathrm{e}-3$ | 1.0578 |
| 20 | $2^{-6}$ | $5.1217 \mathrm{e}-3$ | $5.2777 \mathrm{e}-3$ | $6.1287 \mathrm{e}-4$ | $5.3131 \mathrm{e}-3$ | 1.0374 |
|  |  |  |  |  |  |  |
| 40 | $2^{-4}$ | $2.7265 \mathrm{e}-3$ | $2.6335 \mathrm{e}-3$ | $2.4513 \mathrm{e}-3$ | $3.5978 \mathrm{e}-3$ | 1.3196 |
| 40 | $2^{-5}$ | $2.7129 \mathrm{e}-3$ | $2.6311 \mathrm{e}-3$ | $1.2257 \mathrm{e}-3$ | $2.9026 \mathrm{e}-3$ | 1.0699 |
| 40 | $2^{-6}$ | $2.7099 \mathrm{e}-3$ | $2.6304 \mathrm{e}-3$ | $6.1290 \mathrm{e}-4$ | $2.7009 \mathrm{e}-3$ | 0.9967 |

Table 4.2: Error, estimators and effectivity index for the second case (4.28) with $\varepsilon=0$.
an efficient estimation of the error, the effectivity index being close to 1.

## Random case

Let us now analyse the random case. The true error err in (4.25) is computed using the standard Monte-Carlo method with sample size $K=100$. Moreover, for the first case (4.27), the reference solution is computed using $N_{r e f}=80$ and $\tau_{r e f}=2^{-9}$ while we use again $N_{r e f}=80$ and $\tau_{r e f}=2^{-7}$ for the second case (4.28). We choose $w_{\theta}=1 / 3$ in (4.26), value obtained by considering either case with the same mesh for the approximation and the reference solution, for instance with the coarsest mesh parameters $N=10$ and $\tau=2^{-4}$. Notice that we get similar value for the case $N=N_{r e f}$ and $\tau=\tau_{r e f}$. We report in Tables 4.3, 4.4 and 4.5 the results we get for the first example (4.27) with $\varepsilon=0.4, \varepsilon=0.2$ and $\varepsilon=0.1$, respectively.

| $N$ | $\tau$ | $e r r$ | $w_{\eta} \eta$ | $w_{\gamma} \gamma$ | $w_{\theta} \theta$ | est | e.i. |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 10 | $2^{-4}$ | $3.0838 \mathrm{e}-1$ | $1.0463 \mathrm{e}-1$ | $3.0351 \mathrm{e}-1$ | $8.1729 \mathrm{e}-2$ | $3.3128 \mathrm{e}-1$ | 1.0742 |
| 10 | $2^{-5}$ | $1.8362 \mathrm{e}-1$ | $7.5454 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $8.6368 \mathrm{e}-2$ | $1.9402 \mathrm{e}-1$ | 1.0566 |
| 10 | $2^{-6}$ | $1.3418 \mathrm{e}-1$ | $6.0711 \mathrm{e}-2$ | $7.8871 \mathrm{e}-2$ | $8.9909 \mathrm{e}-2$ | $1.3413 \mathrm{e}-1$ | 0.9996 |
| 10 | $2^{-7}$ | $1.1464 \mathrm{e}-1$ | $5.4026 \mathrm{e}-2$ | $3.9521 \mathrm{e}-2$ | $9.2003 \mathrm{e}-2$ | $1.1378 \mathrm{e}-1$ | 0.9925 |
|  |  |  |  |  |  |  |  |
| 20 | $2^{-4}$ | $3.1287 \mathrm{e}-1$ | $5.2280 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $8.1727 \mathrm{e}-2$ | $3.1864 \mathrm{e}-1$ | 1.0184 |
| 20 | $2^{-5}$ | $1.8145 \mathrm{e}-1$ | $3.7632 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $8.6356 \mathrm{e}-2$ | $1.8265 \mathrm{e}-1$ | 1.0067 |
| 20 | $2^{-6}$ | $1.2883 \mathrm{e}-1$ | $3.0209 \mathrm{e}-2$ | $7.8869 \mathrm{e}-2$ | $8.9889 \mathrm{e}-2$ | $1.2334 \mathrm{e}-1$ | 0.9574 |
| 20 | $2^{-7}$ | $1.0510 \mathrm{e}-1$ | $2.6833 \mathrm{e}-2$ | $3.9520 \mathrm{e}-2$ | $9.1978 \mathrm{e}-2$ | $1.0364 \mathrm{e}-1$ | 0.9861 |
|  |  |  |  |  |  |  |  |
| 40 | $2^{-4}$ | $3.1236 \mathrm{e}-1$ | $2.5771 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $8.1726 \mathrm{e}-2$ | $3.1537 \mathrm{e}-1$ | 1.0097 |
| 40 | $2^{-5}$ | $1.7917 \mathrm{e}-1$ | $1.8571 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $8.6352 \mathrm{e}-2$ | $1.7970 \mathrm{e}-1$ | 1.0029 |
| 40 | $2^{-6}$ | $1.2198 \mathrm{e}-1$ | $1.4930 \mathrm{e}-2$ | $7.8869 \mathrm{e}-2$ | $8.9884 \mathrm{e}-2$ | $1.2051 \mathrm{e}-1$ | 0.9880 |
| 40 | $2^{-7}$ | $1.0494 \mathrm{e}-1$ | $1.3278 \mathrm{e}-2$ | $3.9520 \mathrm{e}-2$ | $9.1971 \mathrm{e}-2$ | $1.0098 \mathrm{e}-1$ | 0.9622 |
|  |  |  |  |  |  |  |  |
| 80 | $2^{-4}$ | $3.0781 \mathrm{e}-1$ | $1.2951 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $8.1726 \mathrm{e}-2$ | $3.1458 \mathrm{e}-1$ | 1.0220 |
| 80 | $2^{-5}$ | $1.8436 \mathrm{e}-1$ | $9.3271 \mathrm{e}-3$ | $1.5649 \mathrm{e}-1$ | $8.6352 \mathrm{e}-2$ | $1.7898 \mathrm{e}-1$ | 0.9708 |
| 80 | $2^{-6}$ | $1.1655 \mathrm{e}-1$ | $7.4936 \mathrm{e}-3$ | $7.8869 \mathrm{e}-2$ | $8.9882 \mathrm{e}-2$ | $1.1981 \mathrm{e}-1$ | 1.0280 |
| 80 | $2^{-7}$ | $1.0339 \mathrm{e}-1$ | $6.6615 \mathrm{e}-3$ | $3.9520 \mathrm{e}-2$ | $9.1970 \mathrm{e}-2$ | $1.0032 \mathrm{e}-1$ | 0.9703 |

Table 4.3: Error, estimators and effectivity index for the first case (4.27) with $\varepsilon=0.4$.

By analysing the results for this first case, we see that the (weighted) error estimator defined in (4.26) provides a good control of the error. Indeed, the effectivity index remains close to one for any value of $N, \tau$ and $\varepsilon$. Moreover, examining the behaviour of the error into more details, we see that each contribution $w_{\eta} \eta, w_{\gamma} \gamma$ and $w_{\theta} \theta$ efficiently controls the error. For instance, let us consider the case $N=80$ for which the FE error is negligible. When $\varepsilon=0.1$, the time estimator is dominant for any value of $\tau$ and the error is indeed divided by two when $\tau$ is halved. On the contrary, for $\varepsilon=0.4$, the stochastic estimator is dominant for $\tau=2^{-6}$ and $\tau=2^{-7}$ and we can indeed observe it on the error: for the various time steps, the error decreases by a factor 1.67,

| $N$ | $\tau$ | $e r r$ | $w_{\eta} \eta$ | $w_{\gamma} \gamma$ | $w_{\theta} \theta$ | $e s t$ | e.i. |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 10 | $2^{-4}$ | $3.0097 \mathrm{e}-1$ | $1.0463 \mathrm{e}-1$ | $3.0351 \mathrm{e}-1$ | $4.0865 \mathrm{e}-2$ | $3.2363 \mathrm{e}-1$ | 1.0753 |
| 10 | $2^{-5}$ | $1.6461 \mathrm{e}-1$ | $7.5454 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $4.3184 \mathrm{e}-2$ | $1.7902 \mathrm{e}-1$ | 1.0875 |
| 10 | $2^{-6}$ | $9.8203 \mathrm{e}-2$ | $6.0711 \mathrm{e}-2$ | $7.8871 \mathrm{e}-2$ | $4.4955 \mathrm{e}-2$ | $1.0921 \mathrm{e}-1$ | 1.1121 |
| 10 | $2^{-7}$ | $7.3308 \mathrm{e}-2$ | $5.4026 \mathrm{e}-2$ | $3.9521 \mathrm{e}-2$ | $4.6002 \mathrm{e}-2$ | $8.1221 \mathrm{e}-2$ | 1.1079 |
|  |  |  |  |  |  |  |  |
| 20 | $2^{-4}$ | $2.9975 \mathrm{e}-1$ | $5.2280 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $4.0863 \mathrm{e}-2$ | $3.1068 \mathrm{e}-1$ | 1.0365 |
| 20 | $2^{-5}$ | $1.5843 \mathrm{e}-1$ | $3.7632 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $4.3178 \mathrm{e}-2$ | $1.6664 \mathrm{e}-1$ | 1.0518 |
| 20 | $2^{-6}$ | $8.6561 \mathrm{e}-2$ | $3.0209 \mathrm{e}-2$ | $7.8869 \mathrm{e}-2$ | $4.4945 \mathrm{e}-2$ | $9.5671 \mathrm{e}-2$ | 1.1052 |
| 20 | $2^{-7}$ | $6.2790 \mathrm{e}-2$ | $2.6833 \mathrm{e}-2$ | $3.9520 \mathrm{e}-2$ | $4.5989 \mathrm{e}-2$ | $6.6308 \mathrm{e}-2$ | 1.0560 |
|  |  |  |  |  |  |  |  |
| 40 | $2^{-4}$ | $2.9500 \mathrm{e}-1$ | $2.5771 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $4.0863 \mathrm{e}-2$ | $3.0733 \mathrm{e}-1$ | 1.0418 |
| 40 | $2^{-5}$ | $1.5450 \mathrm{e}-1$ | $1.8571 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $4.3176 \mathrm{e}-2$ | $1.6340 \mathrm{e}-1$ | 1.0576 |
| 40 | $2^{-6}$ | $8.8589 \mathrm{e}-2$ | $1.4930 \mathrm{e}-2$ | $7.8869 \mathrm{e}-2$ | $4.4942 \mathrm{e}-2$ | $9.1995 \mathrm{e}-2$ | 1.0384 |
| 40 | $2^{-7}$ | $5.9959 \mathrm{e}-2$ | $1.3278 \mathrm{e}-2$ | $3.9520 \mathrm{e}-2$ | $4.5986 \mathrm{e}-2$ | $6.2071 \mathrm{e}-2$ | 1.0352 |
|  |  |  |  |  |  |  |  |
| 80 | $2^{-4}$ | $2.9687 \mathrm{e}-1$ | $1.2951 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $4.0863 \mathrm{e}-2$ | $3.0652 \mathrm{e}-1$ | 1.0325 |
| 80 | $2^{-5}$ | $1.5454 \mathrm{e}-1$ | $9.3271 \mathrm{e}-3$ | $1.5649 \mathrm{e}-1$ | $4.3176 \mathrm{e}-2$ | $1.6260 \mathrm{e}-1$ | 1.0522 |
| 80 | $2^{-6}$ | $8.6499 \mathrm{e}-2$ | $7.4936 \mathrm{e}-3$ | $7.8869 \mathrm{e}-2$ | $4.4941 \mathrm{e}-2$ | $9.1084 \mathrm{e}-2$ | 1.0530 |
| 80 | $2^{-7}$ | $5.5422 \mathrm{e}-2$ | $6.6615 \mathrm{e}-3$ | $3.9520 \mathrm{e}-2$ | $4.5985 \mathrm{e}-2$ | $6.0998 \mathrm{e}-2$ | 1.1006 |

Table 4.4: Error, estimators and effectivity index for the first case (4.27) with $\varepsilon=0.2$.

| $N$ | $\tau$ | $e r r$ | $w_{\eta} \eta$ | $w_{\gamma} \gamma$ | $w_{\theta} \theta$ | $e s t$ | e.i. |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 10 | $2^{-4}$ | $2.9570 \mathrm{e}-1$ | $1.0463 \mathrm{e}-1$ | $3.0351 \mathrm{e}-1$ | $2.0432 \mathrm{e}-2$ | $3.2169 \mathrm{e}-1$ | 1.0879 |
| 10 | $2^{-5}$ | $1.5506 \mathrm{e}-1$ | $7.5454 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $2.1592 \mathrm{e}-2$ | $1.7507 \mathrm{e}-1$ | 1.1291 |
| 10 | $2^{-6}$ | $8.7940 \mathrm{e}-2$ | $6.0711 \mathrm{e}-2$ | $7.8871 \mathrm{e}-2$ | $2.2477 \mathrm{e}-2$ | $1.0204 \mathrm{e}-1$ | 1.1603 |
| 10 | $2^{-7}$ | $6.1673 \mathrm{e}-2$ | $5.4026 \mathrm{e}-2$ | $3.9521 \mathrm{e}-2$ | $2.3001 \mathrm{e}-2$ | $7.0780 \mathrm{e}-2$ | 1.1477 |
|  |  |  |  |  |  |  |  |
| 20 | $2^{-4}$ | $2.9410 \mathrm{e}-1$ | $5.2280 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $2.0432 \mathrm{e}-2$ | $3.0865 \mathrm{e}-1$ | 1.0495 |
| 20 | $2^{-5}$ | $1.5116 \mathrm{e}-1$ | $3.7632 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $2.1589 \mathrm{e}-2$ | $1.6239 \mathrm{e}-1$ | 1.0743 |
| 20 | $2^{-6}$ | $7.7490 \mathrm{e}-2$ | $3.0209 \mathrm{e}-2$ | $7.8869 \mathrm{e}-2$ | $2.2472 \mathrm{e}-2$ | $8.7395 \mathrm{e}-2$ | 1.1278 |
| 20 | $2^{-7}$ | $4.6148 \mathrm{e}-2$ | $2.6833 \mathrm{e}-2$ | $3.9520 \mathrm{e}-2$ | $2.2995 \mathrm{e}-2$ | $5.3015 \mathrm{e}-2$ | 1.1488 |
|  |  |  |  |  |  |  |  |
| 40 | $2^{-4}$ | $2.9313 \mathrm{e}-1$ | $2.5771 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $2.0431 \mathrm{e}-2$ | $3.0528 \mathrm{e}-1$ | 1.0415 |
| 40 | $2^{-5}$ | $1.4913 \mathrm{e}-1$ | $1.8571 \mathrm{e}-2$ | $1.5649 \mathrm{e}-1$ | $2.1588 \mathrm{e}-2$ | $1.5906 \mathrm{e}-1$ | 1.0666 |
| 40 | $2^{-6}$ | $7.4518 \mathrm{e}-2$ | $1.4930 \mathrm{e}-2$ | $7.8869 \mathrm{e}-2$ | $2.2471 \mathrm{e}-2$ | $8.3356 \mathrm{e}-2$ | 1.1186 |
| 40 | $2^{-7}$ | $4.1056 \mathrm{e}-2$ | $1.3278 \mathrm{e}-2$ | $3.9520 \mathrm{e}-2$ | $2.2993 \mathrm{e}-2$ | $4.7611 \mathrm{e}-2$ | 1.1596 |
|  |  |  |  |  |  |  |  |
| 80 | $2^{-4}$ | $2.9480 \mathrm{e}-1$ | $1.2951 \mathrm{e}-2$ | $3.0351 \mathrm{e}-1$ | $2.0431 \mathrm{e}-2$ | $3.0447 \mathrm{e}-1$ | 1.0328 |
| 80 | $2^{-5}$ | $1.4910 \mathrm{e}-1$ | $9.3271 \mathrm{e}-3$ | $1.5649 \mathrm{e}-1$ | $2.1588 \mathrm{e}-2$ | $1.5825 \mathrm{e}-1$ | 1.0613 |
| 80 | $2^{-6}$ | $7.4526 \mathrm{e}-2$ | $7.4936 \mathrm{e}-3$ | $7.8869 \mathrm{e}-2$ | $2.2471 \mathrm{e}-2$ | $8.2349 \mathrm{e}-2$ | 1.1050 |
| 80 | $2^{-7}$ | $3.8214 \mathrm{e}-2$ | $6.6615 \mathrm{e}-3$ | $3.9520 \mathrm{e}-2$ | $2.2992 \mathrm{e}-2$ | $4.6204 \mathrm{e}-2$ | 1.2091 |

Table 4.5: Error, estimators and effectivity index for the first case (4.27) with $\varepsilon=0.1$.
1.58 and 1.13. The case $\varepsilon=0.2$ presents an intermediate stage with ratios $1.92,1.79$ and 1.56 . Similar reasoning can be made for any other cases, namely that the saturation of the error is well explained by the domination of one of the error estimators. To conclude on this example, we finally mention that the slight increase of e.i. when $\tau$ decreases in Table 4.5 is due to the fact that the error is computed with respect to a reference solution. Indeed, if we consider the deterministic case $\varepsilon=0$ with $N=80$ and $\tau=2^{-7}$, the error with respect to the reference solution is 0.0299529 yielding an effectivity index of about 1.36 , see also Table 4.1.

The results for the second case with $\varepsilon=0.5$ and $\varepsilon=0.25$ are provided in Tables 4.6 and 4.7, respectively.

| $N$ | $\tau$ | $e r r$ | $w_{\eta} \eta$ | $w_{\gamma} \gamma$ | $w_{\theta} \theta$ | $e s t$ | e.i. |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 10 | $2^{-4}$ | $1.0989 \mathrm{e}-2$ | $1.0500 \mathrm{e}-2$ | $2.4507 \mathrm{e}-3$ | $5.2493 \mathrm{e}-3$ | $1.1992 \mathrm{e}-2$ | 1.0913 |
| 10 | $2^{-5}$ | $1.1020 \mathrm{e}-2$ | $1.0491 \mathrm{e}-2$ | $1.2254 \mathrm{e}-3$ | $5.2393 \mathrm{e}-3$ | $1.1790 \mathrm{e}-2$ | 1.0699 |
| 10 | $2^{-6}$ | $1.1140 \mathrm{e}-2$ | $1.0488 \mathrm{e}-2$ | $6.1275 \mathrm{e}-4$ | $5.2356 \mathrm{e}-3$ | $1.1738 \mathrm{e}-2$ | 1.0537 |
|  |  |  |  |  |  |  |  |
| 20 | $2^{-4}$ | $7.2634 \mathrm{e}-3$ | $5.2838 \mathrm{e}-3$ | $2.4512 \mathrm{e}-3$ | $5.2568 \mathrm{e}-3$ | $7.8461 \mathrm{e}-3$ | 1.0802 |
| 20 | $2^{-5}$ | $7.1864 \mathrm{e}-3$ | $5.2790 \mathrm{e}-3$ | $1.2257 \mathrm{e}-3$ | $5.2469 \mathrm{e}-3$ | $7.5432 \mathrm{e}-3$ | 1.0496 |
| 20 | $2^{-6}$ | $6.8839 \mathrm{e}-3$ | $5.2777 \mathrm{e}-3$ | $6.1287 \mathrm{e}-4$ | $5.2431 \mathrm{e}-3$ | $7.4646 \mathrm{e}-3$ | 1.0844 |
|  |  |  |  |  |  |  |  |
| 40 | $2^{-4}$ | $5.7040 \mathrm{e}-3$ | $2.6335 \mathrm{e}-3$ | $2.4513 \mathrm{e}-3$ | $5.2591 \mathrm{e}-3$ | $6.3720 \mathrm{e}-3$ | 1.1171 |
| 40 | $2^{-5}$ | $5.3548 \mathrm{e}-3$ | $2.6311 \mathrm{e}-3$ | $1.2257 \mathrm{e}-3$ | $5.2491 \mathrm{e}-3$ | $5.9982 \mathrm{e}-3$ | 1.1202 |
| 40 | $2^{-6}$ | $5.5691 \mathrm{e}-3$ | $2.6304 \mathrm{e}-3$ | $6.1290 \mathrm{e}-4$ | $5.2454 \mathrm{e}-3$ | $5.8999 \mathrm{e}-3$ | 1.0594 |

Table 4.6: Error, estimators and effectivity index for the second case (4.28) with $\varepsilon=0.5$.

| $N$ | $\tau$ | $e r r$ | $w_{\eta} \eta$ | $w_{\gamma} \gamma$ | $w_{\theta} \theta$ | $e s t$ | e.i. |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 10 | $2^{-4}$ | $1.0142 \mathrm{e}-2$ | $1.0500 \mathrm{e}-2$ | $2.4507 \mathrm{e}-3$ | $2.6247 \mathrm{e}-3$ | $1.1097 \mathrm{e}-2$ | 1.0942 |
| 10 | $2^{-5}$ | $1.0155 \mathrm{e}-2$ | $1.0491 \mathrm{e}-2$ | $1.2254 \mathrm{e}-3$ | $2.6197 \mathrm{e}-3$ | $1.0882 \mathrm{e}-2$ | 1.0717 |
| 10 | $2^{-6}$ | $1.0167 \mathrm{e}-2$ | $1.0488 \mathrm{e}-2$ | $6.1275 \mathrm{e}-4$ | $2.6178 \mathrm{e}-3$ | $1.0827 \mathrm{e}-2$ | 1.0650 |
|  |  |  |  |  |  |  |  |
| 20 | $2^{-4}$ | $5.7001 \mathrm{e}-3$ | $5.2838 \mathrm{e}-3$ | $2.4512 \mathrm{e}-3$ | $2.6284 \mathrm{e}-3$ | $6.3903 \mathrm{e}-3$ | 1.1211 |
| 20 | $2^{-5}$ | $5.6392 \mathrm{e}-3$ | $5.2790 \mathrm{e}-3$ | $1.2257 \mathrm{e}-3$ | $2.6234 \mathrm{e}-3$ | $6.0210 \mathrm{e}-3$ | 1.0677 |
| 20 | $2^{-6}$ | $5.6824 \mathrm{e}-3$ | $5.2777 \mathrm{e}-3$ | $6.1287 \mathrm{e}-4$ | $2.6216 \mathrm{e}-3$ | $5.9247 \mathrm{e}-3$ | 1.0426 |
|  |  |  |  |  |  |  |  |
| 40 | $2^{-4}$ | $3.6562 \mathrm{e}-3$ | $2.6335 \mathrm{e}-3$ | $2.4513 \mathrm{e}-3$ | $2.6296 \mathrm{e}-3$ | $4.4563 \mathrm{e}-3$ | 1.2188 |
| 40 | $2^{-5}$ | $3.6174 \mathrm{e}-3$ | $2.6311 \mathrm{e}-3$ | $1.2257 \mathrm{e}-3$ | $2.6246 \mathrm{e}-3$ | $3.9132 \mathrm{e}-3$ | 1.0818 |
| 40 | $2^{-6}$ | $3.6337 \mathrm{e}-3$ | $2.6304 \mathrm{e}-3$ | $6.1290 \mathrm{e}-4$ | $2.6227 \mathrm{e}-3$ | $3.7647 \mathrm{e}-3$ | 1.0361 |

Table 4.7: Error, estimators and effectivity index for the second case (4.28) with $\varepsilon=0.25$.

Looking at the estimators for the case $\varepsilon=0.5$, we see that the FE error is dominant when $N=10$, the FE and stochastic errors balanced for $N=20$ and the stochastic error is dominant when $N=40$. We indeed observe this behaviour for the error. First, it remains more or less constant when changing the time step. Moreover, it decreases by a factor about 1.6 when
doubling $N$ from 10 to 20 , while the reduction of the error is only about 1.2 from $N=20$ to $N=40$. When diminishing the level of uncertainty, taking $\varepsilon=0.25$, the stochastic error is lower and the error decreases by a factor 1.8 when increasing $N$ from 10 to 20 and a factor 1.6 comparing the error for $N=20$ and $N=40$. Finally, the FE and stochastic error estimators are balanced when $N=40$.

## Conclusions

We have considered in this chapter the heat equation with random Robin boundary conditions. Under the assumption of small uncertainty, we have used a perturbation technique for the stochastic space approximation. In addition, the finite element method and the (implicit) backward Euler scheme have been used for the space and time discretizations, respectively. The a posteriori error estimator we have obtained for the approximation of the first term in the expansion consists in three distinct terms controlling each source of error. In the numerical experiments, we have introduced a weighted error estimator, with weights tuned numerically, and we have tested its efficiency on two different examples.

# 5 Error analysis for the stochastic collocation method 

## Introduction

In the previous chapters, we have used a perturbation approach for the stochastic space approximation. Such technique is no longer appropriate for problems with large variability. An alternative is to use the stochastic Galerkin or the stochastic collocation methods that present potentially much faster convergence rate than Monte-Carlo type methods and can handle large uncertainties. The advantage of the stochastic collocation method is that, as sampling methods, it requires only the solution of decoupled deterministic problems and thus allows the re-usability of deterministic solvers. However, it suffers from the so-called curse of dimensionality when tensor grids are used, namely the performance of the method deteriorates as the number of random variables increases. A remedy is then to exploit the possible anisotropy of the solution, in the sense that the different random variables might not have the same influence on the solution. Example of works in this direction are the anisotropic sparse grid method proposed in [96] or the quasi-optimal sparse grids method introduced in [20]. In the latter, the adaptive algorithm is based on a priori error estimates whose constants are numerically tuned during the process, yielding what the authors called an a priorila posteriori strategy for which the proof of convergence has been obtained in [94]. An a posteriori sparse grid algorithm has been proposed in [95], where the adaptive process is driven by profit indicators obtained by solving additional PDEs. The method is applicable to a wide range of problems, including for instance the case of unbounded random variables or non-nested grids and can be combined with a Monte Carlo sampling, using a control variate technique, to handle rough random field [98]. However, the error indicators proposed so far are heuristic and do not provide a certified control of the error. The goal here is to derive a guaranteed upper bound of the error and use the stochastic error estimator to steer an adaptive process yielding an approximate solution with prescribed accuracy.

In this chapter, we thus present a residual-based a posteriori error estimate accounting both the stochastic collocation and the Finite Element error. We consider again the model problem of Chapter 1, namely a diffusion equation with a random diffusion coefficient that depends
in an affine manner of a finite number of random variables. We start by briefly recalling the SC method before presenting the error estimate. We give then possible adaptive algorithms, focusing on the stochastic space adaptation since the physical space adaptation can be done following a standard procedure. Finally, we give some preliminary numerical results to test the efficiency of a simple version of our sparse grid adaptive strategy.

### 5.1 Problem statement

Let $D \subset \mathbb{R}^{d}$ be an open bounded domain with Lipschitz continuous boundary $\partial D$ and let $(\Omega, \mathscr{F}, P)$ be a complete probability space. We consider the diffusion problem:
find $u: D \times \Omega \rightarrow \mathbb{R}$ such that $P$-a.e. in $\Omega$, in other words a.s. in $\Omega$, the following equation holds

$$
\left\{\begin{array}{rll}
-\operatorname{div}(a(\mathbf{x}, \omega) \nabla u(\mathbf{x}, \omega)) & =f(\mathbf{x}) &  \tag{5.1}\\
\mathbf{x} \in D \\
u(\mathbf{x}, \omega) & =0 & \mathbf{x} \in \partial D
\end{array}\right.
$$

with deterministic forcing term $f \in L^{2}(D)$ and random field $a$ on $(\Omega, \mathscr{F}, P)$ over $L^{\infty}(D)$. We assume that the random diffusion coefficient $a$ is uniformly bounded from below and above and that it depends affinely on a finite number of random variables. More precisely, we assume that $a$ satisfies the two following properties:

$$
\begin{equation*}
\exists 0<a_{\min } \leq a_{\max }<\infty: \quad P\left(\omega \in \Omega: a_{\min } \leq a(\mathbf{x}, \omega) \leq a_{\max } \forall \mathbf{x} \in \bar{D}\right)=1 \tag{5.2}
\end{equation*}
$$

and

$$
\begin{equation*}
a(\mathbf{x}, \omega)=a_{0}(\mathbf{x})+\sum_{n=1}^{N} a_{n}(\mathbf{x}) Y_{n}(\omega) \tag{5.3}
\end{equation*}
$$

where $\left\{Y_{n}\right\}_{n=1}^{N}$ are independent random variables. Thanks to the Doob-Dynkin Lemma, the solution $u$ depends on the same random variables as the diffusion coefficient $a$, i.e. we have $u(\mathbf{x}, \omega)=u\left(\mathbf{x}, Y_{1}(\omega), \ldots, Y_{N}(\omega)\right)$. Let us introduce $\Gamma=\Gamma_{1} \times \ldots \times \Gamma_{N}$ with $\Gamma_{n}=Y_{n}(\Omega)$ for $n=1, \ldots, N$. Moreover, let $\rho: \Gamma \rightarrow \mathbb{R}_{+}$be the joint probability density function of the random vector $\mathbf{Y}=\left(Y_{1}, \ldots, Y_{N}\right)$, which factorizes as $\rho(\mathbf{y})=\Pi_{n=1}^{N} \rho_{n}\left(y_{n}\right)$ for all $\mathbf{y}=\left(y_{1}, \ldots, y_{N}\right) \in \Gamma$. We can then replace the probability space $(\Omega, \mathscr{F}, P)$ by $(\Gamma, B(\Gamma), \rho(\mathbf{y}) d \mathbf{y})$, where $B(\Gamma)$ denotes the Borel $\sigma$-algebra defined on $\Gamma$ and $\rho(\mathbf{y}) d \mathbf{y}$ the probability measure of $\mathbf{Y}$. Finally, we define the Bochner space

$$
\begin{equation*}
L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right):=\left\{v: \Gamma \rightarrow H_{0}^{1}(D) \mid v \text { is strongly measurable and }\|v\|_{L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)}<\infty\right\} \tag{5.4}
\end{equation*}
$$

with

$$
\|v\|_{L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)}:=\left(\int_{\Gamma}\|\nabla v(\mathbf{y})\|_{L^{2}(D)}^{2} \rho(\mathbf{y}) d \mathbf{y}\right)^{\frac{1}{2}} .
$$

The (parametric, pointwise) weak formulation of problem (5.1) reads:
find $u: \Gamma \rightarrow H_{0}^{1}(D)$ such that

$$
\begin{equation*}
\int_{D} a(\mathbf{x}, \mathbf{y}) \nabla u(\mathbf{x}, \mathbf{y}) \cdot \nabla v(\mathbf{x}) d \mathbf{x}=\int_{D} f(\mathbf{x}) v(\mathbf{x}) d \mathbf{x} \quad \forall v \in H_{0}^{1}(D), \rho \text {-a.e. in } \Gamma . \tag{5.5}
\end{equation*}
$$

By a straightforward application of Lax-Milgram's lemma, assumption (5.2) ensures the wellposedness of problem (5.5), namely that there exists a unique solution $u \in L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)$ which satisfies the a priori estimate

$$
\|u\|_{L_{\rho}^{2}\left(\Gamma ; H_{0}^{1}(D)\right)} \leq \frac{C_{P}}{a_{\min }}\|f\|_{L^{2}(D)}
$$

Moreover, it has been shown (see for instance [7]) that the parametric solution $u$ of problem (5.5) is analytic with respect to each parameter $y_{n} \in \Gamma_{n}, n=1, \ldots, N$.

### 5.2 Stochastic collocation finite element method

In this section, we briefly present the stochastic collocation finite element method (SC-FEM for short) for solving numerically PDEs with random input data, following closely [115] and focusing on the model problem (5.1). We also refer to [7, 124] for a complete discussion on this method. The idea is to proceed in two steps: first a semi-discretization of problem (5.5) using the FEM for the physical space approximation and then the application a collocation method for the stochastic space approximation using global polynomials in $\mathbf{y}$. We thus seek for an approximate solution in a space $\mathbb{P}(\Gamma) \otimes V_{h}$, with $\mathbb{P}(\Gamma) \subset L_{\rho}^{2}(\Gamma)$ a polynomial space on $\Gamma$ and $V_{h}$ a FE subspace of $V=H_{0}^{1}(D)$.

More precisely, for any $h>0$, let $\mathscr{T}_{h}$ be a regular triangulation of $D$ with elements $T$ of diameter $h_{T} \leq h$. We assume that the exists a constant $c>0$ satisfying

$$
\begin{equation*}
\frac{h_{T}}{\rho_{T}} \leq c \quad \forall T \in \mathscr{T}_{h}, \forall h>0 \tag{5.6}
\end{equation*}
$$

where $\rho_{T}=\sup \{\operatorname{diam}(B): B$ is a ball contained in $T\}$. We consider $V_{h} \subset V$ a finite element space of dimension $N_{h}$ constituted of continuous piecewise polynomials on $\mathscr{T}_{h}$. The semidiscretized problem is therefore given by:
find $u_{h}: \Gamma \rightarrow V_{h}$ such that

$$
\begin{equation*}
\int_{D} a(\mathbf{x}, \mathbf{y}) \nabla u_{h}(\mathbf{x}, \mathbf{y}) \cdot \nabla v_{h}(\mathbf{x}) d \mathbf{x}=\int_{D} f(\mathbf{x}) v_{h}(\mathbf{x}) d \mathbf{x} \quad \forall v_{h} \in V_{h}, \rho \text {-a.e. in } \Gamma . \tag{5.7}
\end{equation*}
$$

The problem (5.7) is then further discretized by considering a set $\left\{\mathbf{y}_{1}, \ldots, \mathbf{y}_{N_{c}}\right\}$ of $N_{c}$ collocation
points in $\Gamma$ and building the global polynomial approximation

$$
\begin{equation*}
u_{h, N_{c}}(\mathbf{y})=\sum_{k=1}^{N_{c}} u_{h}\left(\mathbf{y}_{k}\right) L_{k}(\mathbf{y}) \tag{5.8}
\end{equation*}
$$

for appropriate multivariate (for instance Lagrange) polynomials $L_{k}$, where $u_{h}\left(\mathbf{y}_{k}\right)$ is the solution of problem (5.7) with $\mathbf{y}=\mathbf{y}_{k}$. Notice that if $L_{k}$ satisfies $L_{k}\left(\mathbf{y}_{l}\right)=\delta_{k l}$, then the method presented above for the stochastic space approximation is a collocation method in the sense of [109], see [124].

A possible choice for the collocation points $\mathbf{y}_{k} \in \Gamma$ is to take the Cartesian product of the abscissas in each direction. However, using such tensor grid would rapidly become computationally unaffordable due to the curse of dimensionality: the number of nodes increases exponentially with $N$. To alleviate this drawback, the idea is to use a so-called sparse grid, first introduced by Smolyak in [113]. Let us define

$$
\begin{equation*}
\mathscr{U}_{n}^{m\left(i_{n}\right)}: C^{0}\left(\Gamma_{n}\right) \rightarrow \mathbb{P}_{m\left(i_{n}\right)-1}\left(\Gamma_{n}\right) \tag{5.9}
\end{equation*}
$$

a sequence of univariate polynomial interpolant operators along each direction $\Gamma_{n}$ for $n=$ $1, \ldots, N$. Here, $m\left(i_{n}\right)$ denotes the number of collocation points used to build the interpolant of level $i_{n}$ and $\mathbb{P}_{q}\left(\Gamma_{n}\right)$ is the space of polynomials in $y_{n}$ of degree at most $q$. The function $m$ should satisfy $m(0)=0, m(1)=1$ and $m(i)<m(i+1)$ for any $i \geq 1$. Moreover, let $I \subset \mathbb{N}_{+}^{N}$ be a multi-index set, where $\mathbb{N}_{+}=\{1,2, \ldots\}$ denotes the positive integers. Setting $\mathscr{U}_{n}^{0}=0$ for $n=1, \ldots, N$, we define then the sparse grid interpolant $S_{I}$ by

$$
\begin{equation*}
u_{h, I}(\mathbf{y})=S_{I}\left[u_{h}\right](\mathbf{y})=\sum_{\mathbf{i} \in I} \Delta^{\mathbf{m}(\mathbf{i})}\left(u_{h}\right)(\mathbf{y}) \tag{5.10}
\end{equation*}
$$

where

$$
\Delta^{\mathbf{m}(\mathbf{i})}=\bigotimes_{n=1}^{N} \Delta_{n}^{m\left(i_{n}\right)}=\bigotimes_{n=1}^{N}\left(\mathscr{U}_{n}^{m\left(i_{n}\right)}-\mathscr{U}_{n}^{m\left(i_{n}-1\right)}\right)
$$

and $\mathbf{m}(\mathbf{i})=\left(m\left(i_{1}\right), \ldots, m\left(i_{N}\right)\right)$. The operators $\Delta_{n}^{m\left(i_{n}\right)}$ and $\Delta^{\mathbf{m ( i )}}$ are often referred to as difference (or detail) and hierarchical surplus operators, respectively. In what follows, we assume that

$$
\begin{equation*}
u_{h}(\mathbf{y})=\sum_{\mathbf{i} \in \mathbb{N}_{+}^{N}} \Delta^{\mathbf{m}(\mathbf{i})}\left(u_{h}\right)(\mathbf{y}) \quad \rho \text {-a.e. in } \Gamma \tag{5.11}
\end{equation*}
$$

which holds if $u$ is sufficiently smooth in $\mathbf{y}$ and if the operators $\mathscr{U}_{n}^{m\left(i_{n}\right)}$ in (5.9) are such that $\otimes_{n=1}^{N} \mathscr{U}_{n}^{m\left(i_{n}\right)} u \rightarrow u$ in $V$ as $\mathbf{i} \rightarrow \infty$. Finally, we mention that the operator $S_{I}$ in (5.10) can be equivalently written as a linear combination of tensor grid interpolations, see for instance [122], as

$$
\begin{equation*}
S_{I}\left[u_{h}\right](\mathbf{y})=\sum_{\mathbf{i} \in I} c_{\mathbf{i}} \bigotimes_{n=1}^{N} \mathscr{u}_{n}^{m\left(i_{n}\right)}\left(u_{h}\right)(\mathbf{y}), \quad c_{\mathbf{i}}=\sum_{\substack{\mathbf{j} \in\{0,1\}^{N} \\(\mathbf{i}+\mathbf{j}) \in I}}(-1)^{|\mathbf{j}|} \tag{5.12}
\end{equation*}
$$

in which many of the coefficients $c_{\mathbf{i}}$ are actually zero, namely whenever $\mathbf{i}+\mathbf{1} \in I$. We then call
sparse grid the set of $N_{c}$ collocation points needed by (5.12) to compute $S_{I}\left[u_{h}\right]$. To summarize, the sparse grid interpolant $S_{I}$ is characterized by the multi-index set $I$, the function $m$ defining the number of collocation points on each level and the type of univariate nodes. One example, see for instance [18], is to consider

$$
I(l)=\left\{\mathbf{i} \in \mathbb{N}_{+}^{N}: \sum_{n=1}^{N}\left(i_{n}-1\right) \leq l\right\}
$$

with

$$
m(i)= \begin{cases}0 & \text { if } i=0  \tag{5.13}\\ 1 & \text { if } i=1 \\ 2^{i-1}+1 & \text { if } i>1\end{cases}
$$

and Clenshaw-Curtis nodes, yielding nested grids. Here $l$ denotes the level of the sparse grid. Remark that I must contain the multi-index 1, which allows to approximate constant functions.

In what follows, the only restriction on $I$ will be that it is a downward closed set (a.k.a. lower set), i.e. it satisfies

$$
\begin{equation*}
\forall \mathbf{i} \in I, \quad \mathbf{i}-\mathbf{e}_{j} \in I \quad \forall j=1, \ldots, N \text { such that } i_{j}>1 \tag{5.14}
\end{equation*}
$$

We give in Figure 5.1 an example of two multi-index sets satisfying or not this condition. The set on the left does not satisfy $(5.14)$ because $(3,2)$ is in the set while $(2,2)$ is not. This condition is necessary to get good approximation properties, see for instance [66]. Moreover, our error estimate will only be valid in the case $S_{I}$ is interpolatory, i.e. it satisfies $S_{I}[f]\left(\mathbf{y}_{k}\right)=f\left(\mathbf{y}_{k}\right)$ for $k=1, \ldots, N_{c}$ where $\left\{\mathbf{y}_{1}, \ldots, \mathbf{y}_{N_{c}}\right\}$ are the collocation points in the sparse grid underlying the multi-index set $I$ and function $m$. Notice that such property requires the use of nested nodes.

Finally, we introduce the notion of margin $M_{I}$, reduced margin $R_{I}$ and boundary $\partial I$ of a multi-index set $I$, see Figure 5.1-right for an illustration, defined respectively by

$$
\begin{aligned}
M_{I} & =\left\{\mathbf{i} \in \mathbb{N}_{+}^{N} \backslash I: \mathbf{i}-\mathbf{e}_{n} \in I \text { for some } n \in\{1, \ldots, N\}\right\} \\
R_{I} & =\left\{\mathbf{i} \in M_{I}: \mathbf{i}-\mathbf{e}_{n} \in I \text { for all } n=1, \ldots, N \text { with } i_{n}>1\right\} \\
\partial I & =\left\{\mathbf{i} \in I: \mathbf{i}+\mathbf{e}_{n} \notin I \text { for some } 1 \leq n \leq N\right\} .
\end{aligned}
$$

Notice that for a downward closed multi-index set $I$ and $\mathbf{j} \notin I$, then $I \cup\{\mathbf{j}\}$ is downward closed if and only if $\mathbf{j} \in R_{I}$.

From now on, unless otherwise clearly stated, we assume that $I$ is downward closed and that the operator $S_{I}$ is interpolatory.


Figure 5.1: Non-downward closed set (left), downward closed set (middle) and multi-index set with its margin and reduced margin (right).

### 5.3 Residual-based a posteriori error estimate

We will now derive an a posteriori error estimate for the error $u-S_{I}\left[u_{h}\right]$ which consists of two parts controlling the finite element and stochastic collocation errors, respectively. We first give two results that we will use in the derivation of the error estimate.

Proposition 5.3.1. Let $S_{I}$ be the operator defined in (5.10). Then for any $f, g \in C^{0}(\Gamma)$ we have

$$
S_{I}[f g]=S_{I}\left[f S_{I}[g]\right]
$$

Proof. Since $S_{I}$ is assumed to be interpolatory, we have $S_{I}[g]\left(\mathbf{y}_{k}\right)=g\left(\mathbf{y}_{k}\right)$ for all $k=1, \ldots, N_{c}$. By the definition of $S_{I}$, we get then for any $\mathbf{y} \in \Gamma$

$$
\begin{aligned}
S_{I}\left[f S_{I}[g]\right](\mathbf{y}) & =\sum_{k=1}^{N_{c}}\left(f S_{I}[g]\right)\left(\mathbf{y}_{k}\right) L_{k}(\mathbf{y})=\sum_{k=1}^{N_{c}} f\left(\mathbf{y}_{k}\right) S_{I}[g]\left(\mathbf{y}_{k}\right) L_{k}(\mathbf{y}) \\
& =\sum_{k=1}^{N_{c}} f\left(\mathbf{y}_{k}\right) g\left(\mathbf{y}_{k}\right) L_{k}(\mathbf{y})=S_{I}[f g](\mathbf{y}) .
\end{aligned}
$$

For any multi-index set $I$, let us define the polynomial space $\mathbb{P}_{I}$ by

$$
\mathbb{P}_{I}=\sum_{\mathbf{i} \in I} \mathbb{P}_{m\left(i_{1}\right)-1} \otimes \ldots \otimes \mathbb{P}_{m\left(i_{N}\right)-1}
$$

Notice that since we are using nested points, we have $N_{c}=\operatorname{dim}\left(\mathbb{P}_{I}\right)$ with $N_{c}$ the number of collocation points in the sparse grid. Moreover, we have the following crucial approximation properties.

Proposition 5.3.2. Let $S_{I}$ be the operator defined if (5.10). Then

1. $S_{I}[f] \in \mathbb{P}_{I} \quad \forall f \in C^{0}(\Gamma)$
2. $\quad S_{I}$ is exact on $\mathbb{P}_{I}$, i.e. $S_{I}[f]=f \quad \forall f \in \mathbb{P}_{I}$.

Proof. See [11].

Finally, we introduce the (generalized) jump of a function $\varphi$ across an edge $e \in \mathscr{T}_{h}$ in the direction of $\mathbf{n}_{e}$ as in Chapter 1 by

$$
[\varphi]_{\mathbf{n}_{e}}(\mathbf{x}):=\left\{\begin{aligned}
\lim _{t \rightarrow 0^{+}}\left(\varphi\left(\mathbf{x}+t \mathbf{n}_{e}\right)-\varphi\left(\mathbf{x}-t \mathbf{n}_{e}\right)\right) & \text { if } e \not \subset \partial D \\
0 & \text { if } e \subset \partial D
\end{aligned}\right.
$$

We can now state our residual-based a posteriori error estimate.
Proposition 5.3.3. Let $u$ and $u_{h}$ be the solutions of (5.5) and (5.7), respectively and let $S_{I}\left[u_{h}\right]$ be the sparse grid approximation of $u_{h}$ computed using the multi-index set $I$. There exists a constant $C>0$ depending only on the mesh aspect ratio such that for any $p \in[1, \infty]$ we have

$$
\begin{equation*}
\left\|u-S_{I}\left[u_{h}\right]\right\|_{L_{\rho}^{p}(\Gamma ; V)} \leq \frac{1}{a_{\min }}\left[C \eta_{I}+\zeta_{I}\right] \tag{5.15}
\end{equation*}
$$

where

$$
\begin{equation*}
\eta_{I}=\sum_{k=1}^{N_{c}} \eta_{I, k}\left\|L_{k}\right\|_{L_{\rho}^{p}(\Gamma)}, \quad \eta_{I, k}:=\left(\sum_{T \in \mathscr{T}_{h}} \eta_{I, k, T}^{2}\right)^{\frac{1}{2}} \tag{5.16}
\end{equation*}
$$

with

$$
\begin{equation*}
\eta_{I, k, T}:=h_{T}^{2}\left\|f+\nabla \cdot\left(a\left(\mathbf{y}_{k}\right) \nabla u_{h}\left(\mathbf{y}_{k}\right)\right)\right\|_{L^{2}(D)}^{2}+\sum_{e \subset \partial T} h_{e}\left\|\frac{1}{2}\left[a\left(\mathbf{y}_{k}\right) \nabla u_{h}\left(\mathbf{y}_{k}\right) \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2} \tag{5.17}
\end{equation*}
$$

and

$$
\begin{equation*}
\zeta_{I}=\sum_{\mathbf{i} \in M_{I}} \zeta_{I, \mathbf{i}}, \quad \zeta_{I, \mathbf{i}}:=\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla S_{I}\left[u_{h}\right]\right)\right\|_{L_{\rho}^{p}\left(\Gamma ; L^{2}(D)\right)} \tag{5.18}
\end{equation*}
$$

Proof. In what follows, all equations hold $\rho$-a.e. in $\Gamma$ without specifically mentioning it. Moreover, the dependence of each function on variables will not necessarily be indicated, unless ambiguity arises. For any $v \in V$ we have

$$
\begin{align*}
\int_{D} a \nabla\left(u-S_{I}\left[u_{h}\right]\right) \cdot \nabla v= & \int_{D} f v-\int_{D} a \nabla S_{I}\left[u_{h}\right] \cdot \nabla v \\
= & \underbrace{S_{I}\left[\int_{D} f v-\int_{D} a \nabla u_{h} \cdot \nabla v\right]}_{=: \mathrm{I}} \\
& +\underbrace{S_{I}\left[\int_{D} a \nabla u_{h} \cdot \nabla v\right]-\int_{D} a \nabla S_{I}\left[u_{h}\right] \cdot \nabla v}_{=: \mathrm{II}} \tag{5.19}
\end{align*}
$$

For the second equality, we have used that $f$ is assumed to be deterministic and thus $S_{I}[f]=f$ for any multi-index set $I$. We analyse the terms I and II separately. For the first term, thanks to the Galerkin orthogonality we have

$$
\begin{align*}
\mathrm{I} & =\sum_{k=1}^{N_{c}}\left[\int_{D} f v-\int_{D} a\left(\mathbf{y}_{k}\right) \nabla u_{h}\left(\mathbf{y}_{k}\right) \cdot \nabla v\right] L_{k}(\mathbf{y}) \\
& =\sum_{k=1}^{N_{c}}\left[\int_{D} f\left(v-v_{h}\right)-\int_{D} a\left(\mathbf{y}_{k}\right) \nabla u_{h}\left(\mathbf{y}_{k}\right) \cdot \nabla\left(v-v_{h}\right)\right] L_{k}(\mathbf{y}) \tag{5.20}
\end{align*}
$$

for any $v_{h} \in V_{h}$. We take $v_{h}=I_{h} v$ the Clément interpolant of $v$ for which we have the following interpolation error bounds, see also (1.26) and (1.28)

$$
\begin{equation*}
\left\|v-I_{h} v\right\|_{L^{2}(T)} \leq C h_{T}\|\nabla v\|_{L^{2}(N(T))} \quad \text { and } \quad\left\|v-I_{h} v\right\|_{L^{2}(e)} \leq C h_{e}^{\frac{1}{2}}\|\nabla v\|_{L^{2}\left(N\left(T_{e}\right)\right)} \tag{5.21}
\end{equation*}
$$

for any element $T$ and any edge $e$. Here, for an internal edge $e, T_{e}$ is the union of the two elements touching $e$ and $N(T)$ (resp. $N\left(T_{e}\right)$ ) denotes the patch of elements associated to $T$ (resp. $T_{e}$ ). After splitting the integral in (5.20) over each element $T$ and integrating by part, we obtain

$$
\begin{equation*}
\mathrm{I} \leq C \sum_{k=1}^{N_{c}}\left|L_{k}(\mathbf{y})\right| \eta_{I, k}\|\nabla v\|_{L^{2}(D)} \tag{5.22}
\end{equation*}
$$

with $\eta_{I, k}$ defined in (5.16). Notice that this term $\eta_{I, k}$ is deterministic, namely it does not depend on $\mathbf{y}$. It controls the FE error made when solving approximately the problem for the collocation point $\mathbf{y}_{k}$. We now bound the second term II. We first notice that, thanks to Proposition 5.3.1, we have $S_{I}\left[a \nabla u_{h}\right]=S_{I}\left[a \nabla S_{I}\left[u_{h}\right]\right]$ since $S_{I}$ is assumed to be interpolatory. Therefore, using relation (5.11) we get

$$
\begin{align*}
\text { II } & =\int_{D}\left(S_{I}\left[a \nabla S_{I}\left[u_{h}\right]\right]-a \nabla S_{I}\left[u_{h}\right]\right) \cdot \nabla v=-\int_{D} \sum_{\mathbf{i} \notin I} \Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla S_{I}\left[u_{h}\right]\right) \cdot \nabla v \\
& =-\int_{D_{\mathbf{i} \in M_{I}} \sum \Delta^{\mathbf{m ( i )}}\left(a \nabla S_{I}\left[u_{h}\right]\right) \cdot \nabla v} \\
& \leq\left\|\sum_{\mathbf{i} \in M_{I}} \Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla S_{I}\left[u_{h}\right]\right)\right\|_{L^{2}(D)}\|\nabla v\|_{L^{2}(D)} . \tag{5.23}
\end{align*}
$$

We have used the fact that $a$ depends in an affine way on the random variables, see (5.3), to restrict the summation over the multi-indices of the margin $M_{I}$ of $I$. Indeed, by Proposition 5.3.2 we have

$$
S_{I}\left[u_{h}\right] \in \mathbb{P}_{I}, \quad \text { where } \quad \mathbb{P}_{I}=\sum_{\mathbf{i} \in I} \mathbb{P}_{\mathbf{m}(\mathbf{i})-\mathbf{1}} \quad \text { with } \quad \mathbb{P}_{\mathbf{m}(\mathbf{i})-\mathbf{1}}=\mathbb{P}_{m\left(i_{1}\right)-1} \otimes \ldots \otimes \mathbb{P}_{m\left(i_{N}\right)-1}
$$

and by assumption

$$
a \in \mathbb{P}_{\mathbf{0}}+\sum_{n=1}^{N} \mathbb{P}_{\mathbf{e}_{n}}, \quad \text { with } \mathbb{P}_{\mathbf{e}_{n}}=\mathbb{P}_{0} \otimes \ldots \otimes \mathbb{P}_{0} \otimes \underbrace{\mathbb{P}_{1}}_{n^{t h} \text { index }} \otimes \mathbb{P}_{0} \ldots \otimes \mathbb{P}_{0} .
$$

Therefore, we have $a \nabla S_{I}\left[u_{h}\right] \in \sum_{n=1}^{N} \sum_{\mathbf{i} \in I} \mathbb{P}_{\mathbf{m}(\mathbf{i})-\mathbf{1}+\mathbf{e}_{n}} \subset \mathbb{P}_{I \cup M_{I}}$ and thus

$$
\begin{equation*}
\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla S_{I}\left[u_{h}\right]\right)=0 \quad \forall \mathbf{i} \notin I \cup M_{I} \tag{5.24}
\end{equation*}
$$

using again Proposition 5.3.2, namely that $S_{I \cup M_{I}}$ is exact on $\mathbb{P}_{I \cup M_{I}}$. Thanks to the uniform lower bound $a_{\min }$ on $a$, taking then $v=u(\mathbf{y})-S_{I}\left[u_{h}\right](\mathbf{y})$ in (5.19) and using the bounds (5.22) and (5.23) for the terms I and II, respectively, yields

$$
\begin{equation*}
\left\|\nabla\left(u(\mathbf{y})-S_{I}\left[u_{h}\right](\mathbf{y})\right)\right\|_{L^{2}(D)} \leq \frac{1}{a_{\min }}\left(C \sum_{k=1}^{N_{c}}\left|L_{k}(\mathbf{y})\right| \eta_{I, k}+\left\|\sum_{\mathbf{i} \in M_{I}} \Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla S_{I}\left[u_{h}\right]\right)(\mathbf{y})\right\|_{L^{2}(D)}\right) \tag{5.25}
\end{equation*}
$$

To conclude the proof, it only remains to take the $L_{\rho}^{p}(\Gamma)$ norm on both sides of the last inequality and to use the triangle inequality for the norm $L_{\rho}^{p}\left(\Gamma ; L^{2}(D)\right)$ to take out the sum over the multiindices $\mathbf{i} \in M_{I}$.

Notice that in this proof, we have strongly used the fact that $S_{I}$ is interpolatory and that $a$ depends in an affine way on the random variables. The latter allows us to restrict the summation over all the multi-indices outside $I$ in the bound of II to the multi-indices belonging to the margin $M_{I}$. Moreover, it is worth mentioning that equation (5.25) yields a pointwise (in y) error estimate.

Remark 5.3.4. The spatial error estimate $\eta_{I}$ in (5.16) depends on $\left\|L_{k}(\mathbf{y})\right\|_{L_{\rho}^{p}(\Gamma)}, k=1, \ldots, N_{c}$, i.e. on the stability constant of the operator $S_{I}$. These quantities can be bounded using the Lebesgue constant for $S_{I}$, whose growth depends on the choice of the function $m$ and the family of interpolation points used by $\mathscr{U}_{n}^{m(i)}, n=1, \ldots, N$. For instance, when using a doubling rule for $m$ as in (5.13) and Clenshaw-Curtis nodes, the Lebesgue constant associated with the operator $S_{I}$ can be bounded by $|I|^{2}$ [45]. As an alternative, we could bound the term I in (5.20) as follows

$$
\begin{aligned}
I= & \sum_{T \in \mathscr{T}_{h}}\left[\int_{T} \sum_{k=1}^{N_{c}} L_{k}(\mathbf{y})\left(f+\nabla \cdot\left(a\left(\mathbf{y}_{k}\right) \nabla u_{h}\left(\mathbf{y}_{k}\right)\right)\right)\left(v-v_{h}\right)+\right. \\
& \left.\frac{1}{2} \sum_{e \subset \partial T} \int_{e} \sum_{k=1}^{N_{c}} L_{k}(\mathbf{y})\left[a\left(\mathbf{y}_{k}\right) \nabla u_{h}\left(\mathbf{y}_{k}\right) \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\left(v-v_{h}\right)\right] \\
\leq & C\left(\sum_{T \in \mathscr{T}_{h}} \eta_{I, T}^{2}\right)^{\frac{1}{2}}\|\nabla v\|_{L^{2}(D)}
\end{aligned}
$$

with
$\eta_{I, T}(\mathbf{y})^{2}:=h_{T}^{2}\left\|\sum_{k=1}^{N_{c}} L_{k}(\mathbf{y})\left(f+\nabla \cdot\left(a\left(\mathbf{y}_{k}\right) \nabla u_{h}\left(\mathbf{y}_{k}\right)\right)\right)\right\|_{L^{2}(T)}^{2}+\sum_{e \subset \partial T} h_{e}\left\|\frac{1}{2} \sum_{k=1}^{N_{c}}\left[a\left(\mathbf{y}_{k}\right) \nabla u_{h}\left(\mathbf{y}_{k}\right) \cdot \mathbf{n}_{e}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2}$.

Since $\left(\sum_{T \in \mathscr{F}_{h}} \eta_{I, T}^{2}\right)^{\frac{1}{2}} \leq \sum_{T \in \mathscr{F}_{h}} \eta_{I, T}$, we can then replace (5.15) by

$$
\begin{equation*}
\left\|u-S_{I}\left[u_{h}\right]\right\|_{L_{\rho}^{p}(\Gamma ; V)} \leq \frac{1}{a_{\min }}\left[C \sum_{T \in \mathscr{T}_{h}}\left\|\eta_{I, T}\right\|_{L_{\rho}^{p}(\Gamma)}+\zeta_{I}\right] . \tag{5.27}
\end{equation*}
$$

Mesh refinement, using the error estimate of Proposition 5.3 .3 or the one proposed here would lead to different adaptive strategies. The estimator in (5.16) gives an estimation of the spatial error for each collocation point, that is further localized on each element $T \in \mathscr{T}_{h}$. Indeed, the estimator $\eta_{I, k, T}$ in (5.17) is an indicator of the FE error for element $T$ and collocation point $\mathbf{y}_{k}$. Therefore, different spatial meshes could be considered for each collocation point. On the contrary, the estimator in (5.26) gives an estimation of the spatial error for each element $T \in \mathscr{T}_{h}$ and contains the contribution of all the collocation point. In this case, the same spatial mesh would then be used for all the collocation points.

### 5.3.1 An abstract reformulation of the problem

We consider the (pointwise in $\mathbf{y}$ ) abstract problem:

$$
\begin{equation*}
\text { find: } u(\mathbf{y}) \in V \quad \text { such that } \quad \mathscr{A}(u, v ; \mathbf{y})=\mathscr{F}(v ; \mathbf{y}) \quad \forall v \in V, \rho \text {-a.e. in } \Gamma \text {. } \tag{5.28}
\end{equation*}
$$

Using the finite element method for the physical space approximation, we get the following semi-discretized problem:

$$
\begin{equation*}
\text { find } u_{h}(\mathbf{y}) \in V_{h} \quad \text { such that } \quad \mathscr{A}\left(u_{h}, v_{h} ; \mathbf{y}\right)=\mathscr{F}\left(v_{h} ; \mathbf{y}\right) \quad \forall v_{h} \in V_{h}, \rho \text {-a.e. in } \Gamma . \tag{5.29}
\end{equation*}
$$

Lax-Milgram's lemma ensures the well-posedness of problems (5.28) and (5.29) under the assumptions that the bilinear form $\mathscr{A}$ is (uniformly in $\mathbf{y}$ ) continuous and coercive and that the linear functional $\mathscr{F}$ is continuous. In particular, we assume that there exist two constants $\underline{\alpha}, \bar{\alpha}>0$ such that $\rho$-a.e. in $\Gamma$

$$
\underline{\alpha}\|v\|_{V}^{2} \leq \mathscr{A}(v, v ; \mathbf{y}) \quad \text { and } \quad|\mathscr{A}(u, v ; \mathbf{y})| \leq \bar{\alpha}\|u\|_{V}\|v\|_{V} \quad \forall u, v \in V .
$$

We can then derive the following a posteriori error estimate.

Proposition 5.3.5. Let $u$ and $u_{h}$ be the solutions of (5.28) and (5.29), respectively and let $u_{h, I}=S_{I}\left[u_{h}\right]$ be the sparse grid approximation of $u_{h}$ computed using the multi-index set I. If the series in (5.11) converge absolutely, then

$$
\left\|u-u_{h, I}\right\|_{L_{\rho}^{p}(\Gamma ; V)} \leq \frac{1}{\underline{\alpha}}\left[\left\|R\left(u_{h} ; \cdot\right)\right\|_{L_{\rho}^{p}\left(\Gamma ; V^{\prime}\right)}+\bar{\alpha} \sum_{\mathbf{i} \notin I}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left[u_{h}\right]\right\|_{L_{\rho}^{p}(\Gamma ; V)}\right]
$$

where the residual $R$ is defined for any $w, v \in V$ and any $\mathbf{y} \in \Gamma$ by

$$
<R(w ; \mathbf{y}), v>:=F(v ; \mathbf{y})-\mathscr{A}(w, v ; \mathbf{y})
$$

with $\langle\cdot, \cdot\rangle$ the duality pairing bracket between $V$ and $V^{\prime}$.

We highlight that in this proposition, $S_{I}$ is not assumed to be interpolatory and the dependence on $\mathbf{y}$ of the coefficients in $\mathscr{A}$ and $\mathscr{F}$ is not specified. In particular, we do not assume an affine dependency. However, the absolute convergence of the series in (5.11) is required and the estimator is not computable as is since it contains an infinite series. A computable estimator can however be obtained if we are able to provide estimation of the tail of the series.

Proof. For any $\nu \in V$ and $\rho$-a.e. in $\Gamma$ we have

$$
\begin{aligned}
\mathscr{A}\left(u(\mathbf{y})-u_{h, I}(\mathbf{y}), v ; \mathbf{y}\right) & =F(v ; \mathbf{y})-\mathscr{A}\left(u_{h, I}(\mathbf{y}), v ; \mathbf{y}\right) \\
& =\underbrace{F(v ; \mathbf{y})-A\left(u_{h}(\mathbf{y}), v ; \mathbf{y}\right)}_{=: \mathrm{I}}+\underbrace{A\left(u_{h}(\mathbf{y})-u_{h, I}(\mathbf{y}), v ; \mathbf{y}\right)}_{=: \mathrm{II}} .
\end{aligned}
$$

Bounding each term separately, we easily obtain

$$
\mathrm{I}=<R\left(u_{h}(\mathbf{y}) ; \mathbf{y}\right), v>\leq\left\|R\left(u_{h}(\mathbf{y}) ; \mathbf{y}\right)\right\|_{V^{\prime}}\|\nu\|_{V}
$$

and

$$
\mathrm{II} \leq \bar{\alpha}\left\|u_{h}(\mathbf{y})-u_{h, I}(\mathbf{y})\right\|_{V}\|v\|_{V}=\bar{\alpha}\left\|\left(i d-S_{I}\right) u_{h}(\mathbf{y})\right\|_{V}\|v\|_{V} \leq \bar{\alpha} \sum_{\mathbf{i} \notin I}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left[u_{h}\right](\mathbf{y})\right\|_{V}\|v\|_{V}
$$

where id denotes the identity operator. For the second term, we have used the relation (5.11), namely that the sparse grid approximation converges $\rho$-a.e. in $\Gamma$. Therefore, thanks to the coercivity of $\mathscr{A}$, taking $v=u(\mathbf{y})-u_{h, I}(\mathbf{y}) \rho$-a.e. in $\Gamma$ we get

$$
\left\|u(\mathbf{y})-u_{h, I}(\mathbf{y})\right\|_{V} \leq \frac{1}{\underline{\alpha}}\left[\left\|R\left(u_{h}(\mathbf{y}) ; \mathbf{y}\right)\right\|_{V^{\prime}}+\bar{\alpha} \sum_{\mathbf{i} \notin I}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left[u_{h}\right](\mathbf{y})\right\|_{V}\right] .
$$

The proof is complete by taking the $L_{\rho}^{p}(\Gamma)$ norm on both sides of this last inequality and using the triangle inequality.

Remark 5.3.6. In the special case where $\mathscr{A}(u, v ; \mathbf{y})=\int_{D} a(\mathbf{y}) \nabla u \cdot \nabla v$ and $F(v ; \mathbf{y})=\int_{D} f v$, which corresponds to problem (5.5), the dual norm of the residual $\left\|R\left(u_{h}(\mathbf{y}) ; \mathbf{y}\right)\right\|_{V^{\prime}}$ can be estimated by

$$
\left\|R\left(u_{h}(\mathbf{y}) ; \mathbf{y}\right)\right\|_{V^{\prime}} \leq C \eta(\mathbf{y}) \quad \text { with } \quad \eta(\mathbf{y})=\left(\sum_{T \in \mathscr{T}_{h}} \eta_{T}(\mathbf{y})^{2}\right)^{\frac{1}{2}}
$$

with

$$
\eta_{T}(\mathbf{y})^{2}:=h_{T}^{2}\left\|f+\nabla \cdot\left(a(\mathbf{y}) \nabla u_{h}(\mathbf{y})\right)\right\|_{L^{2}(T)}^{2}+\sum_{e \subset \partial T} h_{e}\left\|\frac{1}{2}\left[a(\mathbf{y}) \frac{\partial u_{h}(\mathbf{y})}{\partial \mathbf{n}_{e}}\right]_{\mathbf{n}_{e}}\right\|_{L^{2}(e)}^{2} .
$$

## Chapter 5. Error analysis for the stochastic collocation method

### 5.4 Adaptive algorithms

The error estimator deduced from Proposition 5.3.3 can be used to adaptively refine the mesh and increase the multi-index set. Such an adaptive strategy aims at reaching a given accuracy of the (FE and stochastic) error with computational cost as low as possible. The theory for mesh adaptation, often referred to as adaptive finite element method (AFEM), is well developed and studied. In particular, the convergence of some adaptive procedures has been provided in many different cases. The first result in this direction is the work by Dörfler [57], where the convergence of an adaptive algorithm for the Poisson equation is given. Over the past decades, much effort has been put in proving convergence of adaptive algorithms (with optimal rate) for various types of problems, see for instance [25, 42, 93, 114]. In the context of parametric/random PDEs, we mention the work by [46] where the convergence of an adaptive algorithm is given when the solution is approximated via a Taylor series. In [58, 59], where the random PDEs are solved with the Stochastic Galerkin FEM, the convergence is proved when the adaptation is performed in both physical and stochastic spaces. In this case, the extension of the results obtained for the AFEM in [42] is straightforward and strongly uses the so-called Galerkin orthogonality property. Finally, for the stochastic collocation method, we mention the paper [20] in which a (quasi-optimal) sparse grid method based on a a priori/a posteriori strategy is proposed and whose convergence is analysed in [94]. Moreover, an a posteriori sparse grid algorithm is given in [95]. So far, at least to our knowledge, there is no proof of convergence for adaptive stochastic collocation methods.

Here, we will use the a posteriori error estimate given in Proposition 5.3.3 to drive an adaptive procedure. We start by considering only stochastic space adaptation since mesh adaptation can be performed in a classical way. The error estimator $\zeta_{I}$ can be used to adaptively enrich the multi-index set $I$ in order to reach a prescribed accuracy while minimizing the computational cost. The proposed adaptive procedure is given in Algorithm 4.

```
Algorithm 4 Adaptive algorithm (stochastic space adaptation)
Require: \(\theta \in(0,1)\) and Tol \(>0\)
Ensure: multi-index set \(I\) such that \(\zeta_{I} \leq\) Tol
    \(I=\{\mathbf{1}\}, u_{I}=S_{I}\left[u_{h}\right], \zeta_{I}=\zeta_{I, \mathbf{1}}\)
    while \(\zeta_{I}>\) Tol do
        \(J=\) new_index \(\left(\theta, I, \zeta_{I}\right) \quad\) select a subset of \(M_{I}\) satisfying (5.30)
        \(I \leftarrow I \cup J \quad\) update the multi-index set
        \(u_{I}=S_{I}\left[u_{h}\right] \quad\) compute the new sparse grid approximation
        \(\zeta_{I}=\sum_{\mathbf{i} \in M_{I}} \zeta_{I, \mathbf{i}} \quad\) compute the error estimator (5.18)
    end while
```

It remains to define the routine new_index of Step 3, namely to define how we select the multi-index set $J \subset M_{I}$ to be added to the current set $I$. Following a so-called Dörfler marking, we choose to select $J$ according to

$$
\begin{equation*}
\text { find } J \subset M_{I}: \quad \sum_{\mathbf{i} \in J} \zeta_{I, \mathbf{i}} \geq \theta \sum_{\mathbf{i} \in M_{I}} \zeta_{I, \mathbf{i}} \quad \text { and } \quad I \cup J \text { downward closed. } \tag{5.30}
\end{equation*}
$$

We can think of several strategies to select $J$ satisfying (5.30), keeping in mind that the goal is to minimize the computational cost. Since the set should remain downward closed at each iteration of the adaptive algorithm, we associate to each multi-index $\mathbf{i}$ a set $A_{\mathbf{i}}$ which consists of all multi-indices that must also be included in $I$ if $\mathbf{i}$ is added to $I$ so that $I$ remains downward closed. Notice that $A_{\mathbf{i}}=\{\mathbf{i}\}$ if $\mathbf{i}$ belongs to the reduced margin. Moreover, we can define a notion of profit for each multi-index $\mathbf{i} \in M_{I}$ as follows

$$
\begin{equation*}
P_{\mathbf{i}}:=\frac{\sum_{\mathbf{j} \in A_{\mathbf{i}}} \zeta_{I, \mathbf{j}}}{\sum_{\mathbf{j} \in A_{\mathbf{i}}} W_{\mathbf{j}}} \tag{5.31}
\end{equation*}
$$

taking into account all elements of $A_{\mathbf{i}}$. Here, we have denoted by $W_{\mathbf{i}}$ the work contribution of the multi-index $\mathbf{i}$, which can be defined by [95]

$$
\begin{equation*}
W_{\mathbf{i}}=\Pi_{n=1}^{N}\left(m\left(i_{n}\right)-m\left(i_{n}-1\right)\right) \tag{5.32}
\end{equation*}
$$

In the case of nested sets of point, as considered here, it corresponds to the number of new points in $\Gamma$ introduced if $\mathbf{i}$ is added to $I$. We could also choose to set $W_{\mathbf{i}}=1$ if we want to drive the adaptation only based on the error indicators. With these definitions of $A_{\mathbf{i}}$ and $P_{\mathbf{i}}$, we can formulate a possible version of the routine new_index.

```
Algorithm 5 new_index
Require: \(\theta, I\) and \(\zeta_{I}\)
Ensure: multi-index set \(J \subset M_{I}\) satisfying (5.30)
    \(J=\varnothing, \varrho=0\)
    while \(\varrho<\theta \zeta_{I}\) do
        \(\mathbf{i}=\operatorname{argmax}_{\mathbf{i} \in M_{I} \backslash J} P_{\mathbf{i}}\)
        \(J \leftarrow J \cup A_{\mathbf{i}}\)
        \(\varrho=\sum_{\mathbf{j} \in J} \zeta_{I, \mathbf{j}}\)
    end while
```

Remark 5.4.1. Notice that the set J returned by Algorithm 5 might not be the optimal set satisfying (5.30). Indeed, a better set could be obtained by re-computing at each iteration the profit $P_{\mathbf{i}}$ in (5.31) of the multi-indices $\mathbf{i} \in M_{I} \backslash\left(R_{I} \cup J\right)$ for which $A_{\mathbf{i}}$ contains a multi-index added at the previous iteration. For such multi-index $\mathbf{i}$, the set $A_{\mathbf{i}}$ has changed and thus the profit.

To summarize, we have to choose the following parameters:

- the value of the Dörfler parameter $\theta \in(0,1)$,
- the value of $p \in[1, \infty]$ for the $L_{\rho}^{p}(\Gamma)$ norm,
- the definition of the work $W_{\mathbf{i}}$ by (5.32) or $W_{\mathbf{i}}=1$ in (5.31).


## Implementation

We give here some details about the computation of the error estimators $\zeta_{I, \mathbf{i}}$ defined in (5.18), with particular attention to the case $\mathbf{i} \in M_{I} \backslash R_{I}$.

We consider the case $p=\infty$. Since the images of the random variables $\Gamma_{n}, n=1, \ldots, N$, are bounded and $u$ is smooth with respect to $\mathbf{y}$, the essential supremum norm can be replaced by the maximum norm. Of course, not all the points of $\Gamma$ can be explored and we choose to approximate the maximum norm searching for the maximum over a given set $\Theta \subset \Gamma$ of finite cardinality. The error is therefore computed using

$$
\begin{aligned}
\left\|u-S_{I}\left[u_{h}\right]\right\|_{L_{\rho}^{\infty}(\Gamma ; V)} & =\max _{\mathbf{y} \in \Gamma}\left|\left\|\nabla\left(u-S_{I}\left[u_{h}\right]\right)(\mathbf{y})\right\|_{L^{2}(D)} \rho(\mathbf{y})\right| \\
& \approx \max _{\mathbf{y} \in \Theta}\left|\left\|\nabla\left(u-S_{I}\left[u_{h}\right]\right)(\mathbf{y})\right\|_{L^{2}(D)} \rho(\mathbf{y})\right|
\end{aligned}
$$

which requires the solution of $|\Theta|$ PDEs to get the value of $u(\mathbf{y})$ for each $\mathbf{y} \in \Theta$. Notice that since the FE error will not be accounted for in the numerical results, all the computation can be done on the same spatial mesh. The computation of the error estimators $\zeta_{I, \mathbf{i}}$ can be done as follows. Let $G$ be any downward closed multi-index set that does not contains $\mathbf{i}$ and such that $G \cup\{\mathbf{i}\}$ is also downward closed. The error estimator for $\mathbf{i}$ is then approximately

$$
\begin{align*}
\zeta_{I, \mathbf{i}} & =\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla S_{I}\left[u_{h}\right]\right)\right\|_{L_{\rho}^{\infty}\left(\Gamma ; L^{2}(D)\right)} \\
& =\left\|S_{G \cup\{\mathbf{i}\}}\left[a \nabla S_{I}\left[u_{h}\right]\right]-S_{G}\left[a \nabla S_{I}\left[u_{h}\right]\right]\right\|_{L_{\rho}^{\infty}\left(\Gamma ; L^{2}(D)\right)} \\
& \approx \max _{\mathbf{y} \in \Theta}\left|\left\|S_{G \cup\{\mathbf{i}\}}\left[a \nabla S_{I}\left[u_{h}\right]\right](\mathbf{y})-S_{G}\left[a \nabla S_{I}\left[u_{h}\right]\right](\mathbf{y})\right\|_{L^{2}(D)} \rho(\mathbf{y})\right| \tag{5.33}
\end{align*}
$$

The key-point here is that no PDE need to be solved to compute (5.33). This formula can be straightforwardly applied for all the multi-indices $\mathbf{i} \in R_{I}$ with $G=I$, since $G \cup\{\mathbf{i}\}$ is downward closed, but a special care is required for the elements in $M_{I} \backslash R_{I}$. The idea is to iteratively increase the multi-index set $I$ to cover the full margin in such a way that it remains downward closed throughout the process. We proceed layer by layer, starting by adding the elements of the reduced margin $R_{I}$, as described in the pseudo-code of Algorithm 6.

```
Algorithm 6 Computation of \(\zeta_{I, \mathbf{i}}\) for all \(\mathbf{i} \in M_{I}\)
Require: \(I, S_{I}\left[u_{h}\right], a\)
Ensure: \(\zeta_{I, \mathbf{i}} \forall \mathbf{i} \in M_{I}\)
    \(G=I\)
    while \(G \neq I \cup M_{I}\) do
        \(R=R_{G} \cap M_{I}\)
        for \(\mathbf{i} \in R\) do
            compute \(\zeta_{I, \mathbf{i}}\) using (5.33)
            \(G \leftarrow G \cup\{\mathbf{i}\}\)
        end for
    end while
```

Notice that $R$ at line 3 is a subset of the neighbours of the element of the previous previous layer. Moreover, the order of selection of the elements of $R$ in the for loop is irrelevant.

Remark 5.4.2. For the case $p \in[1, \infty)$, the $L_{\rho}^{p}(\Gamma)$ norm can be computed (either exactly or approximately) using a Gauss quadrature formula built upon $S_{I(l)}$ with level l high enough. Notice that the larger $p$ and the larger the polynomial degree of the integrand, the larger the level l should be.

## Simplified algorithm

Algorithm 4, based on a Dörfler marking, is designed in the spirit of AFEM. The idea for introducing such algorithm was to prove its convergence as it is done for example in [59] for the Stochastic Galerkin method. We have made several attempts in this direction, for instance to prove that the error estimator satisfies a certain contraction property or to use different markings as it is done in [93] to control the decrease of the data oscillation. Unfortunately, we have not been successful so far, mainly due to the lack of the so-called Galerkin orthogonality valid for both the physical and the stochastic spaces when using the SG-FEM. The proof of convergence of the proposed adaptive algorithm is thus still an open question.

In the numerical results of Section 5.5, we consider a simplified version of Algorithm 4. First of all, we allow the selection of elements of the reduced margin $R_{I}$ only and not of the full margin $M_{I}$. This simplifies the definition of the profits, since we do not need to introduce the sets $A_{\mathbf{i}}$. Indeed, we recall that if $I$ is downward closed, then so is $I \cup\{\mathbf{i}\}$ for any multi-index $\mathbf{i} \in R_{I}$. The second modification is that we add only one multi-index at a time. More precisely, the adaptive algorithm that is used for the numerical experiments of Section 5.5 reads as follows.

```
Algorithm 7 Simplified adaptive algorithm (stochastic space adaptation)
Require: Tol \(>0\)
Ensure: multi-index set \(I\) such that \(\zeta_{I} \leq\) Tol
    \(I=\{\mathbf{1}\}, u_{I}=S_{I}\left[u_{h}\right], \zeta_{I}=\zeta_{I, \mathbf{1}}\)
    while \(\zeta_{I}>\) Tol do
        \(\mathbf{i}=\operatorname{argmax}_{\mathbf{i} \in R_{I}} P_{\mathbf{i}} \quad\) select the multi-index with highest profit
        \(I \leftarrow I \cup\{\mathbf{i}\} \quad\) update the multi-index set
        \(u_{I}=S_{I}\left[u_{h}\right] \quad\) compute the new sparse grid approximation
        \(\zeta_{I}=\sum_{\mathbf{i} \in M_{I}} \zeta_{I, \mathbf{i}} \quad\) compute the error estimator (5.18)
    end while
```

Remark 5.4.3. The adaptive process of Algorithm 7 is driven only by the profit of the elements of the reduced margin $R_{I}$ of the current set $I$. To reduce the computational cost, we could therefore compute $\zeta_{I, \mathbf{i}}$ for $\mathbf{i} \in R_{I}$ only. However, the global error estimator $\zeta_{I}$ would no longer be available and we have to define another stopping criterion for the algorithm. For example, we can prescribe a tolerance Tol on the highest profit, i.e. stop the adaptive procedure when $\max _{\mathbf{i} \in R_{I}} P_{\mathbf{i}}<$ Tol.

### 5.5 Numerical results

We consider here numerical examples to test Algorithm 7. In all what follows, we choose $m$ as in (5.13) and we use Clenshaw-Curtis nodes. The FE error is not accounted here. Moreover, we consider the case $p=\infty$ and we thus consider the error and estimator defined by respectively

$$
\left\|u_{h}-S_{I}\left[u_{h}\right]\right\|_{L_{\rho}^{\infty}\left(\Gamma ; H_{0}^{1}(D)\right)} \quad \text { and } \quad \sum_{\mathbf{i} \in M_{I}}\left\|\Delta^{\mathbf{m ( i})}\left(a \nabla S_{I}\left[u_{h}\right]\right)\right\|_{L_{\rho}^{\infty}\left(\Gamma ; L^{2}(D)\right)}
$$

Be aware that the initialization step is not counted in the number of iterations given below. Therefore, the cardinality of the set $I$ at the $k^{\text {th }}$ iteration is equal to $k+1$.

## First example

For this first example, we consider an inclusion problem with $N=2$ random variables, similar to the one consider in [11] for $N=8$. The physical domain, depicted in Figure 5.1-left, is the unit square $D=(0,1)^{2}$. We identify three subdomains $F, C_{1}$ and $C_{2}$, with $F$ a square centred in the domain with side length equal to 0.2 and $C_{1}$ and $C_{2}$ two circular inclusions of radius 0.13 . We define the random diffusion coefficient by

$$
\begin{equation*}
a(\mathbf{x}, \mathbf{Y}(\omega))=a_{0}(\mathbf{x})+\sum_{n=1}^{2} \gamma_{n} \chi_{n}(\mathbf{x}) Y_{n}(\omega) \quad \text { with } a_{0}=1 \text { and } Y_{n} \sim \mathscr{U}[-0.99,0.99] \tag{5.34}
\end{equation*}
$$

and we set the forcing term to $f(\mathbf{x})=100 \chi_{F}(\mathbf{x})$, where $\chi_{F}$ and $\chi_{n}, n=1,2$, denote the indicator function of each subdomain. The parameters $\gamma_{1}$ and $\gamma_{2}$ are used to introduce anisotropy in the problem, assigning more importance to one or another direction $y_{1}$ or $y_{2}$.

For the numerical experiments of this first example, we have used the following setting. The FE mesh consists of 4961 vertices and 9696 triangles with minimal and maximal diameter $h_{T}$ of about $7.367 \mathrm{e}-3$ and $2.854 \mathrm{e}-2$, respectively. Since we would like to test the efficiency of our error estimator, namely to see if it is a good control of the (stochastic) error, we compute the estimator $\zeta_{I, \mathbf{i}}$ for each multi-index $\mathbf{i}$ of the margin $M_{I}$. We can therefore base the stopping criterion on the global estimator $\zeta_{I}$, see Remark 5.4.3. We set the tolerance to $\mathrm{Tol}=10^{-6}$. Finally, we compute the $L_{\rho}^{\infty}(\Gamma)$ norm approximately using for $\Theta$ a $20 \times 20$ Cartesian grid of equispaced points in each direction.

## Isotropic case

We start with the isotropic case $\gamma_{1}=\gamma_{2}=1$ in (5.34). The mean and the standard deviation of the solution is given in Figure 5.2, while the evolution of the set $I$ during the adaptive process is presented in Figure 5.3. The multi-index with the green cross indicates the selected element at the current iteration of Algorithm 7, i.e. the one with the highest profit that belongs to the reduced margin of the previous set.


Figure 5.2: Geometry of the problem (left), expected value (middle) and standard deviation (right) of the solution for the case $\gamma_{1}=\gamma_{2}=1$.


Figure 5.3: Evolution of $I$ during the adaptive process for the case $\gamma_{1}=\gamma_{2}=1$. From left to right and top to bottom: iterations 3,5,8,10,14 and order of selection of the multi-indices.

We can detect the isotropy of the problem by the symmetrical construction of the multi-index set. For instance, at iteration 11 the point $(2,4)$ is added while $(4,2)$ is selected at the next iteration. Moreover, we see that the estimator provides a good control of the error as shown in Figure 5.4, where the final multi-index set and the corresponding sparse grid are also given. It has been obtained after 17 iterations, yielding a sparse grids of 97 points and an error and an estimator of about $3.4649 \mathrm{e}-7$ and $8.1070 \mathrm{e}-7$, respectively. Finally, we mention that the highest profit of the elements of the reduced margin of this final stage is about $2.3702 \mathrm{e}-8$ and is achieved at $(2,5)$.


Figure 5.4: Final multi-index set $I$ (left), final sparse grid (middle) and error and estimator with respect to the number of points in semi-logarithmic scale (right) for the case $\gamma_{1}=\gamma_{2}=1$.

## Anisotropic case

We now set different values for $\gamma_{1}$ and $\gamma_{2}$ in (5.34) to see if the adaptive algorithm is able to capture the anisotropy of the problem. We start with the trivial case $\gamma_{1}=1$ and $\gamma_{2}=0$, for which no point should be added in the second direction $y_{2}$. This is indeed the result we get, as shown in Figure 5.5. At the end of the adaptive procedure, which requires 4 iterations, the sparse grid consists of 17 points and the error and estimator are about $1.4219 \mathrm{e}-10$ and $1.5276 \mathrm{e}-10$, respectively. The maximal profit among the elements of the reduced margin is $9.5472 \mathrm{e}-12$ and is attained at $(6,1)$.

Finally, we consider the case $\gamma_{1}=1$ and $\gamma_{2}=0.1$. We present in Figure 5.6 the set $I$ at various steps of the adaptive construction. As expected, we can clearly identify a preferred direction, namely the horizontal direction which corresponds to $y_{1}$.

The final situation, reached in 10 iterations, is given in Figure 5.7. In this case, there are 41 points in the sparse grid, the error and estimator are $6.8878 \mathrm{e}-8$ and $1.2500 \mathrm{e}-7$, respectively, and the maximal profit among the elements of the reduced margin is of $1.9995 \mathrm{e}-8$ at $(3,3)$.


Figure 5.5: Final multi-index set $I$ (left), final sparse grid (middle) and error and estimator with respect to the number of points in semi-logarithmic scale (right) for the case $\gamma_{1}=1$ and $\gamma_{2}=0$.


Figure 5.6: Evolution of the multi-index set $I$ during the adaptive process for the case $\gamma_{1}=1$ and $\gamma_{2}=0.1$. From left to right and top to bottom: iterations $4,6,8$ and order of selection of the multi-indices.


Figure 5.7: Final multi-index set $I$ (left) and error and estimator with respect to the number of points in semi-logarithmic scale (right) for the case $\gamma_{1}=1$ and $\gamma_{2}=0.1$.

## Anisotropic case $N=8$

To conclude on this inclusion problem, we consider the case $N=8$ as in [11] and we choose $a$ similarly to (5.34) with $Y_{n} \sim \mathscr{U}[-0.99,0.2]$ for $n=1, \ldots, 8$. The geometry is given in Figure 5.8 -left, where the value of the coefficients $\gamma_{n}, n=1, \ldots, 8$, is also given. The FE mesh we are using contains 3805 vertices and 7416 triangles with minimal and maximal diameter $h_{T}$ of about $1.0041 \mathrm{e}-2$ and $3.1153 \mathrm{e}-2$, respectively. Moreover, a set of 500 points randomly sampled from a multivariate uniform distribution is used for the approximation of the $L_{\rho}^{\infty}(\Gamma)$ norm. In Figure 5.8-right, we give the error and estimator for the 55 first iterations of Algorithm 7, after which the estimator is about $2.5102 \mathrm{e}-3$ and the sparse grid consists of 213 points in $\Gamma$. Moreover, the projection of the obtained multi-index set $I$ over two directions, namely $y_{1}$ and $y_{4}, y_{1}$ and $y_{5}$ and $y_{1}$ and $y_{7}$, is presented in Figure 5.9.


Figure 5.8: Geometry of the problem for $N=8$ with indication of the coefficients $\gamma_{n}, n=1, \ldots, 8$ (left) and error and estimator with respect to the number of points in logarihmic scale for the 55 first iterations (right).

Even though the estimator still provides a reasonable control of the error, it is less efficient than for the case $N=2$. We see several possible explanations for this behaviour and we give a


Figure 5.9: Projection of the multi-index set $I$ obtained after 55 iterarions on ( $y_{1}, y_{4}$ ) (left), $\left(y_{1}, y_{5}\right)$ (middle) and ( $y_{1}, y_{7}$ ) (right).
non-exhaustive list below. First of all, we have not been able to prove that the error estimator provides a lower bound for the error. The difficulties arise, among other, from the lack of Galerkin orthogonality but also from the use of the triangle inequality to localize the estimator on each multi-index of the margin. Moreover, we are not taking into account the error due to the approximation of the $L_{\rho}^{\infty}(\Gamma)$ norm and further investigation should be made in this direction, namely trying to quantify this additional error and perform additional tests with other training sets $\Theta$.

## Second example

As a second numerical experiment, we take again the 2D example investigated in Section 1.7.2 of Chapter 1, namely we choose $f(\mathbf{x})=32\left(x_{1}\left(1-x_{1}\right)+x_{2}\left(1-x_{2}\right)\right)$ and

$$
a(\mathbf{x}, \mathbf{Y}(\omega))=1+\sum_{n=1}^{N} \frac{\cos \left(2 \pi n x_{1}\right)+\cos \left(2 \pi n x_{2}\right)}{(\pi n)^{2}} Y_{n}(\omega) \quad \text { with } \quad Y_{n} \sim \mathscr{U}[-\sqrt{3}, \sqrt{3}]
$$

for $\mathbf{x}=\left(x_{1}, x_{2}\right) \in D$. We use a spatial mesh consisting of 2673 vertices and 5184 triangles with minimum and maximum diameter $h_{T}$ of about 0.01 and 0.04 , respectively. We set again the tolerance to $T o l=10^{-6}$ in Algorithm 7 and the set $\Theta$ for the approximation of the $L_{\rho}^{\infty}(\Gamma)$ norm consists of 500 points in $\Gamma$ randomly sampled from a multivariate uniform distribution. We consider the two cases $N=3$ and $N=5$.

The results for the case $N=3$ are given in Figure 5.10. We plot the error and the estimator with respect to the work, i.e. number of points in the sparse grid. We also give the projection of the final multi-index set $I$ over two directions, namely $y_{1}$ and $y_{3}$. For this final state, obtained in 27 iterations, the error and the estimator are about 4.1493e-7 and 9.1738e-7, respectively, and the grid contains 141 points. Finally, we mention that the multi-index that has been in the last iteration to the final set $I$ is $(4,3,1)$ and that the maximum profit among the elements of $R_{I}$ is about $3.0159 \mathrm{e}-8$ and is reached at $(3,2,3)$.

The Figure 5.11 contains the results for the case $N=5$. The final multi-index set $I$ is projected


Figure 5.10: Error and estimator with respect to the number of points in logarihmic scale (left) and projection of the final multi-index set on $\left(y_{1}, y_{3}\right)$ (right) for the case $N=3$.
on $y_{1}$ and $y_{5}$. The final grid has 469 points, for an error and estimator of about $2.2500 \mathrm{e}-6$ and $9.8095 \mathrm{e}-6$, respectively, and has been reached in 69 iterations. The last multi-index added to the set is $(4,4,1,1,1)$ and the maximum profit among the elements of the reduced margin of the final set is about $7.7365 \mathrm{e}-8$ at $(3,2,1,2,2)$.


Figure 5.11: Error and estimator with respect to the number of points in logarihmic scale (left) and projection of the final multi-index set on $\left(y_{1}, y_{5}\right)$ (right) for the case $N=5$.

In both cases $N=3$ and $N=5$, the error estimator provides a good control of the error, the overestimation being slightly bigger for $N=5$ than $N=3$. Moreover, due to the decay of the $a_{n}$ in $n^{-2}$, the random variables $Y_{n}$ should have less and less influence as $n$ increases. The adaptive algorithm is able to capture this feature, as seen for instance when projecting the obtained multi-index set over two different directions. From this experiment, together with the numerical results obtained for the inclusion problems, we see that the efficiency of the stochastic error estimator seems to be linked to the number of random variables. Further investigation should be made in this direction to determine whether this is indeed the case or if the reason is elsewhere, for instance the error due to the approximation of the $L_{\rho}^{\infty}(\Gamma)$ norm.

## Conclusions

In this last chapter, we went out of the framework of small uncertainties considered in the previous chapters and in which a perturbation technique has been used for the stochastic space approximation. Here, we have considered the stochastic collocation method which is also appropriate for problems with a large amount of randomness but its use becomes challenging for problem in high dimensions. We have proposed a residual-based a posteriori error estimate that controls both the physical and stochastic space discretization. This estimate is valid under quite strong assumptions but that are often meet in practise. First, we have assumed that the random diffusion coefficient depends in an affine way on a finite number of random variables, which is what we get for instance from a (truncated) Karhunen-Loève expansion of a random field. The second assumption is that the sparse grid operator is interpolatory, which requires the use of nested sequences of univariate nodes such as Clenshaw-Curtis or Leja nodes.

We have then proposed an adaptive sparse grid algorithm. The stochastic error estimator, which is localized on each element of the margin of the current multi-index set, is used to select the most profitable elements that should enter the set. The error estimator we have proposed presents the advantage to be computable without solving additional PDEs. However, it has the drawback that the profit need to be recomputed at each iteration of the adaptive process since the residual depends on $S_{I}\left[u_{h}\right]$. We have made some numerical experiments to test the efficiency of a simple version of the adaptive algorithm. These are just preliminary yet promising results. They open the door to many improvements and prospects, including but not limited to

- quantify the error of approximation of the $L_{\rho}^{\infty}(\Gamma)$ norm using a finite number of (deterministic or random) points in $\Gamma$
- test different choices of family of points, such as Leja-sequence of points
- make a comparison with other methods, adaptive or not
- analyse the complexity of the proposed adaptive strategy
- prove the convergence of Algorithm 4
- take the FE error into account and do mesh refinement when the FE error dominates the stochastic one; take either the same mesh for all the collocation points or allow different refinements for the various points, see Remark 5.3.4
- consider the case of infinite number of random variables


## 5.A Miscellaneous results

We give here some preliminary results which might be useful to prove the convergence of Algorithm 4. In what follows, we will write $I_{k}$ and $I_{k+1}$ two successive multi-index sets produced
by the adaptive algorithm, that is $I_{k+1}=I_{k} \cup J_{k}$ with $J_{k} \subset M_{I_{k}}$ obtained using new_index and thus satisfying the Dörfler condition (5.30). Moreover, since we perform only stochastic space adaptation, we assume that there is no error due to FE approximation and the subscript $h$ is no longer indicated in what follows. We write then $u_{k}=S_{I_{k}}[u]$ and $u_{k+1}=S_{I_{k+1}}[u]$ the sparse grid approximation corresponding to $I_{k}$ and $I_{k+1}$, respectively.

First of all, since $a$ depends affinely on the $y_{n}, n=1, \ldots, N$, we have that if $\mathbf{i} \in M_{I}$ then

$$
\begin{equation*}
\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla \Delta^{\mathbf{m}(\mathbf{j})}(u)\right)=0 \quad \forall \mathbf{j} \in I \backslash \partial I . \tag{5.35}
\end{equation*}
$$

Indeed, if $\mathbf{j} \in I \backslash \partial I$ then $\mathbf{j}+\mathbf{e}_{n} \in I$ for all $n=1, \ldots, N$.
For ease of notation, we will write $\|\cdot\|$ instead of $\|\cdot\|_{L_{\rho}^{p}\left(\Gamma ; L^{2}(D)\right)}$ in the sequel.

Proposition 5.A.1. (Estimator reduction I)
If $u_{k+1}=u_{k}$, then $\left\|\nabla\left(u-u_{k+1}\right)\right\|=\left\|\nabla\left(u-u_{k}\right)\right\|$ but

$$
\zeta_{I_{k+1}}<\zeta_{I_{k}}
$$

Proof. First of all, we split the margin of $I_{k+1}$ into two disjoint parts as

$$
M_{I_{k+1}}=\left(M_{I_{k}} \backslash J_{k}\right) \cup\left(M_{J_{k}} \backslash M_{I_{k}}\right)
$$

Using the assumption $u_{k+1}=u_{k}$ we get then

$$
\begin{aligned}
\zeta_{I_{k+1}} & =\sum_{\mathbf{i} \in M_{I_{k+1}}}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla u_{k+1}\right)\right\| \\
& =\sum_{\mathbf{i} \in M_{I_{k+1}}}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla u_{k}\right)\right\| \\
& =\sum_{\mathbf{i} \in M_{I_{k}} \backslash J_{k}}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla u_{k}\right)\right\|+\sum_{\mathbf{i} \in M_{J_{k}} \backslash M_{I_{k}}}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla u_{k}\right)\right\| \\
& =\sum_{\mathbf{i} \in M_{I_{k}} \backslash J_{k}}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla u_{k}\right)\right\| .
\end{aligned}
$$

For the last equality, we have used that $\Delta^{\mathbf{m ( i )}}\left(a \nabla u_{k}\right)=0$ for all $\mathbf{i} \in\left(M_{J_{k}} \backslash M_{I_{k}}\right)$ thanks to (5.24). Indeed, if $\mathbf{i} \in\left(M_{J_{k}} \backslash M_{I_{k}}\right)$ then $\mathbf{i} \notin\left(I_{k} \cup M_{I_{k}}\right)$. Finally, we use the property of $J_{k}$ in (5.30) to obtain

$$
\begin{aligned}
\zeta_{I_{k+1}} & =\sum_{\mathbf{i} \in M_{I_{k}}}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla u_{k}\right)\right\|-\sum_{\mathbf{i} \in J_{k}}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla u_{k}\right)\right\| \\
& \leq(1-\theta) \sum_{\mathbf{i} \in M_{I_{k}}}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla u_{k}\right)\right\| \\
& =\kappa \zeta_{I_{k}},
\end{aligned}
$$

with $\kappa=(1-\theta)<1$ for any choice of the Dörfler parameter $\theta \in(0,1)$.

One way to prove the convergence of Algorithm 4 is to prove a contraction property, for instance on the error, the estimator or some other quantity. The difficulty is therefore to first define the quantity on which we would like to prove a contraction property. We have tried to do it on the estimator, but, unfortunately, we have not been able yet to find a conclusion. So far, we have obtained the following relation

$$
\begin{aligned}
\zeta_{I_{k+1}} & \leq \sum_{\mathbf{i} \in M_{I_{k}} \backslash J_{k}}\left\|\Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla u_{k}\right)\right\|+\sum_{\mathbf{i} \in M_{I_{k+1}}} \| \Delta^{\mathbf{m}(\mathbf{i})}\left(a \nabla\left(u_{k+1}-u_{k}\right) \|\right. \\
& \leq(1-\theta) \zeta_{I_{k}}+\sum_{\mathbf{i} \in M_{I_{k+1}}} \| \Delta^{\mathbf{m} \mathbf{( i )}}\left(a \nabla\left(u_{k+1}-u_{k}\right) \| .\right.
\end{aligned}
$$

## Conclusions and perspectives

In this thesis, error analysis for PDEs with random input data has been performed on various problems with a focus on a posteriori error estimation.

The starting point was the well-studied elliptic diffusion model problem with random diffusion coefficient and affine dependence on the random variables. Assuming small amount of randomness in the model, characterized with the parameter $\varepsilon$, a perturbation technique was used expanding the exact random solution of this problem in powers of $\varepsilon$. Error estimation for the error between the exact solution and the finite element approximation of the truncated expansion has been established in great details, considering different measures of the error. Computing for instance only the first term in the expansion, which is deterministic, the a posteriori error estimate provides information about both sources of error, namely the physical space discretization and the uncertainty, and can be used to balance these two errors. Moreover, such error estimates are the basis for adaptive strategies designed to find an approximation of prescribed accuracy with computational cost as low as possible. Having $a$ posteriori error estimate for the approximation of any order allows us to adaptively choose between mesh refinement and increase of the order of the expansion. The theoretical results have been validated and illustrated through many numerical experiments in one and two physical space dimensions. We are looking forward to perform numerical experiments on adaptive schemes of higher-order in $\varepsilon$. A proof of the lower bound for the explicit stochastic error estimator of the first order approximation, required to prove its efficiency, is still missing at the moment.

Next, steady-state nonlinear problems in random domains have been investigated. For such problems, the so-called domain mapping method has been used to transform the PDEs in random domains into PDEs on a fixed reference domain with random coefficients. All the analysis can then be made on this fixed reference domain and, from a numerical point of view, this method prevents the need of remeshing. Application to the one-dimensional viscous Burger's equation and the incompressible Navier-Stokes equations has been proposed. The well-posedness has been shown, under suitable conditions on the mapping and the input data, using a fixed-point theorem for existence and a variational argument for uniqueness. A posteriori error estimation has been proposed for a specific but rather general form of the mapping, again under the assumption of small perturbation. For the Navier-Stokes problem, two different estimates have been developed, each of them presenting advantages and draw-
backs. Numerical results have been given for both problems. Possible extensions include the consideration of problems for which the mapping is not given analytically, numerical experiments on three-dimensional Navier-Stokes equations and analysis of the time-dependent Burgers and Navier-Stokes equations.

To extend the proposed methodology to other types of problems, a parabolic problem has been analysed next, namely the heat equation with random Robin boundary conditions. In addition to the perturbation technique and the finite element method for the stochastic and physical space approximations, respectively, an implicit time stepping scheme has been used for the time discretization. An a posteriori error estimate for the approximation of the first term in the expansion has been proposed and its efficiency has been investigated through two numerical examples. Application to problems of practical interest could be an interesting direction for a future work.

In the last part of this thesis, a residual-based a posteriori error estimate for the stochastic collocation finite element method has be proposed. The error estimator controlling the randomness in the problem has then be used to drive an adaptive sparse grid algorithm. Finally, promising preliminary numerical examples have been given that open the door to many thrilling perspectives, such as complexity analysis, comparison with other methods, combination with spatial mesh refinement or proof of convergence of adaptive scheme.

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# Curriculum Vitae 

## Personal Data

Name Diane Guignard<br>Date of birth January 14th, 1988<br>Nationality Swiss

## Education

$$
\begin{array}{ll}
2012-2016 & \text { PhD in Mathematics } \\
& \text { Ecole Polytechnique Fédérale de Lausanne, Switzerland. } \\
\text { Thesis advisers: Prof. F. Nobile and Prof. M. Picasso. } \\
2010-2012 & \begin{array}{l}
\text { Master of Science in Applied Mathematics } \\
\text { Ecole Polytechnique Fédérale de Lausanne, Switzerland. } \\
\\
\text { Master thesis at Caltech, Pasadena, USA under the supervision of Prof. T.Y. Hou. } \\
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\end{array}
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$$

## Publications

1. D. Guignard, F. Nobile and M. Picasso. A posteriori error estimation for elliptic partial differential equations with small uncertainties. Numer. Methods Partial Differential Equations, 32(1): 175-212, 2016.
2. D. Guignard, F. Nobile and M. Picasso. A posteriori error estimation for the steady Navier-Stokes equations in random domains. MATHICSE Technical Report 13.2016 (submitted for publication).

## Presentations

- 4th Workshop on Sparse Grids and Applications (Miami, Florida, USA, 4 October 2016)

Contributed talk: A posteriori error estimate and adaptive sparse grid algorithm for random PDEs.

- MATHICSE Retreat 2016 (Leysin, Switzerland, 28 June 2016)

Contributed talk: A posteriori error estimation for PDEs with random input data.

- ECCOMAS 2016 (Crète, Greece, 10 June 2016)

Talk in minisymposium: A posteriori error estimate for the Navier-Stokes equations in random domains solved with a perturbation technique.

- Colloque Numérique Suisse 2016 (Fribourg, Switzerland, 22 April 2016)

Talk: A posteriori error estimation for the steady Navier-Stokes equations in random domains.

- SIAM UQ 2016 (Lausanne, Switzerland, 8 April 2016)

Contributed talk: A Posteriori Error Estimates for Navier-Stokes Equations with Small Uncertainties.

- Reliable Methods of Mathematical Modeling (RMMM) 2015 (Zürich, Switzerland, 30 June 2015)

Talk in minisymposium: A posteriori error estimation for elliptic partial differential equations with small uncertainties.

- International Conference on Adaptive Modeling and Simulation (ADMOS) 2015 (Nantes, France, 9 June 2015)

Contributed talk: A Posteriori Error Estimation for PDEs with Small Uncertainties.

- Swiss Numerical Day 2015 (Geneva, Switzerland, 17 April 2015)

Poster: A posteriori error estimation for partial differential equations with small uncertainties.

- MATHICSE Retreat 2014 (Leysin, Switzerland, 12 June 2014)

Contributed talk: A posteriori error estimation for PDEs with small uncertainties using perturbation methods.


[^0]:    ${ }^{1}$ The function $u_{0}$ in (1.93) is the solution of the problem $-\tau u_{0}^{\prime \prime}+u_{0}^{\prime}=1$ in $(0,1)$ with homogeneous Dirichlet boundary conditions.

[^1]:    ${ }^{2}$ From a numerical point of view, any element which does not contribute to the sum for $\eta_{1}$ will not be refined, i.e. any element which is numerically zero due to machine precision.

[^2]:    ${ }^{3}$ If the factor $\frac{1}{2}$ is replaced by $\frac{1}{4}$ for the jump contribution, see Remark 1.7.1, then we should take $C_{H_{0}^{1}}:=1 / 5$. See Appendix 1.C for more details.

[^3]:    ${ }^{1}$ If we have homogeneous Dirichlet conditions in $x=0$ and $x=1$, we have an a priori estimate. Indeed, it is easy to show that $\left\|u^{\prime}\right\|_{L^{2}(0,1)} \leq \frac{1}{a}\|f\|_{L^{2}(0,1)}$ taking $v=u$ in (3.3) and using the fact that $\int_{0}^{1} b u u^{\prime} u=0$. The existence of a solution can then be proved using for instance Schauder's fixed point theorem while for the uniqueness, it holds under the constraint $C_{F}\|f\|_{L^{2}(0,1)}<\frac{a^{2}}{b}$.

[^4]:    ${ }^{2}$ Notice that we get comparable results if we bound these two terms separately, in which case the estimator due to the uncertainty reads $\eta_{\varepsilon}^{2}=2 \varepsilon^{2} C_{F}^{2}\left(\left\|2 s_{0} f_{0}+s_{0}^{2} f_{1}\right\|_{L^{2}(D)}^{2}+b^{2}\left\|u_{0, h} u_{0, h}^{\prime}\right\|_{L^{2}(D)}^{2}\right)$.

[^5]:    ${ }^{1}$ This first example is similar to the case (3a) considered in [103]. The difference is that here we impose Robin (random) boundary conditions on a part of the boundary.

