ABSTRACT

Title of dissertation: FAST AND OPTIMAL SOLUTION ALGORITHMS FOR PARAMETERIZED PARTIAL DIFFERENTIAL EQUATIONS

Kookjin Lee, Doctor of Philosophy, 2017

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This dissertation presents efficient and optimal numerical algorithms for the solution of parameterized partial differential equations (PDEs) in the context of stochastic Galerkin discretization. The stochastic Galerkin method often leads to a large coupled system of algebraic equations, whose solution is computationally expensive to compute using traditional solvers. For efficient computation of such solutions, we present low-rank iterative solvers, which compute low-rank approximations to the solutions of those systems while not losing much accuracy. We first introduce a low-rank iterative solver for linear systems obtained from the stochastic Galerkin discretization of linear elliptic parameterized PDEs. Then we present a low-rank nonlinear iterative solver for efficiently computing approximate solutions of nonlinear parameterized PDEs, the incompressible Navier–Stokes equations.

Along with the computational issue, the stochastic Galerkin method suffers
from an optimality issue. The method, in general, does not minimize the solution error in any measure. To address this issue, we present an optimal projection method, a least-squares Petrov-Galerkin (LSPG) method. The proposed method is optimal in the sense that it produces the solution that minimizes a weighted $\ell^2$-norm of the solution error over all solutions in a given finite-dimensional subspace. The method can be adapted to minimize the solution error in different weighted $\ell^2$-norms by simply choosing a specific weighting function within the least-squares formulation.
FAST AND OPTIMAL SOLUTION ALGORITHMS FOR PARAMETERIZED PARTIAL DIFFERENTIAL EQUATIONS

by

Kookjin Lee

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<tbody>
<tr>
<td>AE</td>
<td>Absolute difference exponential</td>
</tr>
<tr>
<td>AMG</td>
<td>Algebraic multigrid</td>
</tr>
<tr>
<td>CoV</td>
<td>Coefficient of variation</td>
</tr>
<tr>
<td>CG</td>
<td>Conjugate gradient</td>
</tr>
<tr>
<td>FEM</td>
<td>Finite element method</td>
</tr>
<tr>
<td>GCR</td>
<td>Generalized conjugate residual</td>
</tr>
<tr>
<td>GMRES</td>
<td>General minimum residual</td>
</tr>
<tr>
<td>gPC</td>
<td>Generalized polynomial chaos</td>
</tr>
<tr>
<td>IFISS</td>
<td>Incompressible Flow and Iterative Solver Software</td>
</tr>
<tr>
<td>KL</td>
<td>Karhunen-Loève</td>
</tr>
<tr>
<td>LRP</td>
<td>Low-rank projection</td>
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<tr>
<td>LSC</td>
<td>Least-squares commutator</td>
</tr>
<tr>
<td>LSPG</td>
<td>Least-squares Petrov–Galerkin</td>
</tr>
<tr>
<td>MINRES</td>
<td>Minimum residual</td>
</tr>
<tr>
<td>OoM</td>
<td>Out of memory</td>
</tr>
<tr>
<td>PDE</td>
<td>Partial differential equation</td>
</tr>
<tr>
<td>PGD</td>
<td>Proper generalized decomposition</td>
</tr>
<tr>
<td>PS</td>
<td>Pseudo-spectral</td>
</tr>
<tr>
<td>QoI</td>
<td>Quantity of interest</td>
</tr>
<tr>
<td>SE</td>
<td>Squared difference exponential</td>
</tr>
<tr>
<td>SG</td>
<td>Stochastic Galerkin</td>
</tr>
<tr>
<td>s.p.d.</td>
<td>Symmetric positive definite</td>
</tr>
<tr>
<td>SVD</td>
<td>Singular value decomposition</td>
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<tr>
<td>TD</td>
<td>Total degree</td>
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<tr>
<td>TP</td>
<td>Tensor product</td>
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Chapter 1: Introduction

Forward uncertainty propagation for parameterized algebraic systems is important in a range of applications to characterize the effects of input uncertainties on the output of computational models. Such parameterized algebraic systems arise in many important problems in science and engineering, for which models using stochastic partial differential equations (PDEs) are formulated and where uncertain input parameters are treated as a set of random variables. Examples of such problems include diffusion/ground water flow simulations with uncertain diffusivity/permeability [53,115], solid mechanics with uncertain material properties [43,44], incompressible fluid flow problems with uncertain viscosity [63,113], thermofluid flow problems [58,62], and reacting flow problems with chemical kinetics [29,88] with uncertain inputs. Parameterized algebraic systems also arise in other computational models such as models for reconstructing a high-resolution image from a set of low resolution images [22], and the PageRank algorithm [15,26].

There is a number of sources that cause input uncertainties, for example, an inherent stochastic nature of physical phenomena, and errors in measuring physical properties of objects of interest [46]. If the source of uncertainty comes from a lack of knowledge about physical properties, one feasible approach to handle this
is to collect a finite number of observations, characterize the statistical quantities of the properties, and model the properties as random fields governed by a set of random variables. That is, the physical property can be modeled as a random function such that the value of the random field varies over the spatial domain and the “stochastic domain” (i.e., the image space of the random variables). Suppose, for example, that we are interested in diffusion of chemicals in a medium with an unknown diffusivity. Then the diffusivity can be modeled as a random field based on the statistical quantities (e.g., sample mean and covariance) obtained from a finite number of observations.

There are several ways to model a random field. If the mean and the covariance function of a random field over the spatial domain are known, the random field can be represented as a Karhunen-Loève expansion [67], which is a linear expansion of orthogonal functions that depend on the spatial parameters and for which the coefficients of those functions are pairwise uncorrelated random variables. The orthogonal functions can be obtained by solving an eigenvalue problem associated with the covariance function. In a discrete sense, the KL-expansion is equivalent to principal component analysis [79]. To simulate a random field in terms of a finite number of random variables, the random field can be approximated by truncating the KL-expansion with a finite number of terms so that only the terms with larger variances are retained. There are also alternatives to using the KL-expansion; a random field can be modeled as a linear expansion of certain orthogonal polynomials, which will be introduced in Section 2.1, and as a linear expansion of trigonometric polynomials for weakly stationary random fields [49].
When the uncertain input is modeled as a random field, the model output (i.e., the solution of the algebraic system) also can be modeled as a random field, that is, a random function depending on the spatial location and the same random variables associated with the input random field. This can be seen from the fact that a specific realization of the input parameters gives rise to a deterministic problem and, consequently, leads to a specific realization of the solution function, the function evaluated at those specific values of the input parameters. Thus, the effects of the uncertain input on the model output (i.e., the solution) can be characterized by the statistical properties of the solution such as the mean, the variance, and higher moments of the solution.

The most straightforward approach to obtain statistical moments of the solution is to use the Monte Carlo method [72], which estimates the statistical moments of the solution from a finite number of sample solutions. That is, the Monte Carlo method requires a set of realizations of an input random field and collects solutions of deterministic problems associated with given input random field realizations. Then the statistical moments of the solution can be approximated by the sample moments. The Monte Carlo method is very simple and powerful; the method exhibits \( \frac{1}{\sqrt{N}} \) convergence, where \( N \) is the number of samples, regardless of the dimension of the sample space. At the same time, however, if high accuracy is required in the approximation, \( N \) may need to be very large. Moreover, the Monte Carlo method can be very expensive if solving each deterministic problem associated with a sample is expensive. For faster convergence, there have been many improvements made such as Quasi-Monte Carlo methods [74] using pseudo-random sequences and sampling
methods based on Markov-Chain Monte Carlo methods [71]. Developing an optimal sampling strategy in the uncertainty quantification framework is an active area of research.

Recently, many advanced algorithms have been developed to achieve such statistical characterization of the solution with efficiency. One of the widely used approaches is spectral methods [46, 110], where the solution is expressed as a linear expansion of a finite number of certain orthogonal basis polynomials depending on the input random variables. Once the solution expansion is computed using numerical algorithms, the statistical quantities of the solution can be computed directly and inexpensively by sampling the solution expansion. This approach was inspired by the work [108], which studied the decomposition of a Gaussian random process (or, Gaussian random field), where the Gaussian random process is represented as a linear expansion of Hermite polynomials, which are orthogonal with respect to an inner product induced by the Gaussian probability density function. The series expansion displays a mean-square convergence; the expected value of a squared error goes to zero as the number of terms in the expansion goes to infinity.

In the early work on spectral methods [42, 44, 102] in uncertainty quantification, the Hermite polynomials are used to represent the solution function and numerical algorithms were developed to compute the coefficients of the solution expansion. This approach has been shown to be very successful when the random variables parameterizing the problems follow the Gaussian distribution; the Hermite polynomial expansion exhibits exponential convergence rate for the Gaussian random field [68]. For non-Gaussian random fields, however, the use of the Hermite
polynomials may result in significantly slower convergence [112]. With the recent development of the generalized polynomial chaos (gPC) expansion [112], the type of the orthogonal basis polynomials can be chosen based on the underlying measure of the input random variables, which results in better convergence. The effectiveness of the gPC expansion in characterizing the solution statistics of the stochastic PDEs has been demonstrated in [112–114].

After choosing the type of the orthogonal polynomials, spectral methods require a numerical algorithm to compute coefficients of the solution expansion. The first class of numerical algorithms developed for spectral methods is known as stochastic Galerkin methods [1, 3, 28, 46, 69], which extends a classical Galerkin approach for deterministic equations and, thus, is based on a Galerkin projection technique. As in the finite element method (FEM) for solving PDEs, the stochastic Galerkin method enforces a Galerkin orthogonality condition on the residual of stochastic PDEs with respect to the span of the gPC polynomial basis using an inner product associated with an underlying probability measure of the random variables. This procedure results in a system of (non-)linear equations where the number of equations is the same as the number of unknown coefficients and, thus, the coefficients can be obtained by solving the system of equations. The stochastic Galerkin method is popular for its simplicity (i.e., the trial and test bases are the same) and its optimality in terms of minimizing an energy norm of solution errors when the underlying PDE operator is elliptic and self-adjoint.

Another class of numerical algorithms can be thought of as specialized sampling-based methods. These methods generate a set of independent realizations of random
inputs based on their probability distribution and solve the corresponding deterministic realizations of the problems, and then use the resulting solutions to construct spectral approximations that can be used for simulation. There are several types of numerical algorithms of this type. One of the popular methods is the *stochastic collocation method* \([2,111]\), which solves deterministic problems on a set of predetermined nodes in the space defined by the random variables. The stochastic collocation method computes the spectral approximation by constructing a Lagrange interpolating polynomial. A variant of the stochastic collocation is to compute the coefficients of the solution expansion using quadrature rules. In this approach, the solution is directly projected onto each polynomial basis function exploiting the orthogonality of the polynomial basis functions. This approach is known as a *pseudo-spectral* approach \([109,110]\). Another sampling-based approach for the spectral methods is a polynomial-regression-type approach with gPC expansion (e.g., least-squares regression \([54,96,97]\), least angle regression \([12]\), compressive sampling \([32,50]\)).

Compared with sampling-based methods, the stochastic Galerkin method can lead to smaller errors for a fixed basis dimension \([37,110,111]\). In general, however, the stochastic Galerkin method suffers from two main problems. First, the method typically leads to a large set of coupled deterministic equations, for which computations will be expensive for large-scale applications. The solution function lies on a tensor product space of a spatial domain (a physical space) and a “stochastic domain” (a parameter space), and, after discretization, the number of coupled deterministic equations to be solved is the product of the numbers of basis polynomials in the spatial domain and the stochastic domain (i.e., degrees of freedom in each
domain). When the solution is sought in a high-dimensional space (i.e., dimension of the discrete physical space or the number of basis polynomials on the parameter space is large), the computation of the solution can be very expensive. Secondly, the method may not produce numerical solutions that minimize any measure of the solution error if the underlying PDE operator is not symmetric positive-definite. For many practical problems such as flow problems, PDE operators are not self-adjoint.

In an effort to alleviate the first difficulty, sparse structures of system matrices that can be obtained from the stochastic Galerkin method have been studied. To compute solutions of those systems, efficient iterative algorithms such as Krylov subspace methods \cite{33, 34, 40, 56, 80, 81} and multigrid methods \cite{27, 35, 64, 94} are applied. Matrix-vector products are essential matrix operations in those iterative solvers and those products can be performed very efficiently by exploiting the sparsity structures of system matrices. In combination with specially designed preconditioners \cite{81, 84, 101, 105}, those iterative solvers have been adjusted and successfully applied to many stochastic PDEs. As the size of the problems become larger, however, the computational costs of the iterative solvers increase rapidly, which makes use of those iterative solvers for high-dimensional problem less attractive.

The second issue of the stochastic Galerkin method has not been explored much. In many applications, however, quantities such as solution error or solution residual can be considered as more important metrics for measuring performance of solution methods. Numerical experiments in \cite{73} demonstrated that, for certain classes of stochastic PDEs, the stochastic Galerkin method fails to generate a solution that minimizes a certain norm of solution error. Weighted projection methods
(i.e., Petrov–Galerkin projection techniques) have been proposed to resolve the issue. The weighted projection method successfully minimizes the solution error although the proposed methods require problem-specific projection bases.

In this thesis, we have developed efficient and optimal numerical methods to address and overcome the issues raised above. To address the first problem (high costs), we have developed efficient iterative solvers that decouple matrix operations associated with the spatial domain and the stochastic domain, which makes the computational complexity depend on the sum of the numbers of degrees of freedom in the spatial domain and the stochastic domain rather than their product. In particular, we consider a tensor variant of the Krylov subspace methods that operates in such a decoupled manner so that the computational costs and memory requirements can be significantly reduced. In addition, the variant of the Krylov subspace method will be used to compute a low-rank approximate solution, which further reduces the computational costs. The second problem is addressed using an optimal projection method, the stochastic least-squares Petrov–Galerkin method, which produces solution coefficients that minimize a certain measure of the solution error. We study the behavior of the stochastic Galerkin solution in several error measures and propose an optimization framework that provides an optimal projection basis to minimize a certain measure of the solution error.
1.1 Outline of Thesis

An outline of the thesis is as follows. We begin in Chapter 2 by introducing the stochastic Galerkin method and deriving the stochastic Galerkin system that arises from stochastic diffusion equations. Then we briefly review existing iterative solution methods for a large coupled deterministic system arising from the stochastic Galerkin method.

In Chapter 3, we discuss the use of a low-rank tensor variant of the Krylov subspace method in the stochastic Galerkin setting. For the efficient computation, we propose a two-level rank reduction scheme, which identifies an important subspace in the stochastic domain and compresses tensors of high rank on-the-fly during the iterations. The proposed reduction scheme is a coarse-grid method in that the important subspace can be identified inexpensively in a coarse spatial grid setting. The efficiency of the proposed method is illustrated by numerical experiments on benchmark elliptic linear stochastic PDE problems.

In Chapter 4, we develop a low-rank tensor variant of Newton–Krylov subspace methods for stochastic Navier–Stokes problems in the stochastic Galerkin setting. We base our development on a deterministic variant of a “linearization” scheme and solve a linear system at each nonlinear iteration step using the low-rank Krylov subspace method. We test our method under various settings of the Navier–Stokes equations and compare results with the conventional full-rank method.

In Chapter 5, we propose a new projection framework, stochastic Least-Squares Petrov–Galerkin (LSPG) method, which provides an optimal projection method.
The proposed method is optimal in the sense that it produces the solution that minimizes a weighted $\ell^2$-norm of the residual over all solutions in a given finite-dimensional subspace. With extensive numerical experiments, we show that the weighted LSPG methods outperforms other spectral methods in minimizing corresponding target weighted norms.

In Chapter 6, we draw some conclusions.
Chapter 2: Background: The stochastic Galerkin method

In this chapter, we begin with a brief introduction of the stochastic Galerkin method with stochastic diffusion equations as a model problem. The stochastic Galerkin discretization procedure is discussed only with the stochastic diffusion problem, an extension of the stochastic Galerkin formulation to other linear elliptic PDEs with uncertain input is straightforward.

2.1 Overview of the stochastic Galerkin method

Consider the steady-state stochastic diffusion equation with homogeneous Dirichlet boundary conditions,

\[ \begin{align*}
-\nabla \cdot (a(x, \omega) \nabla u(x, \omega)) &= f(x, \omega) \quad \text{in } D \times \Omega, \\
u(x, \omega) &= 0 \quad \text{on } \partial D \times \Omega,
\end{align*} \tag{2.1} \]

where the diffusion coefficient \( a(x, \omega) \) is a random field and \( \omega \) is an elementary event in a probability space \( (\Omega, \mathcal{F}, P) \). Here, \( \Omega \) is a sample space, \( \mathcal{F} \) and \( P \) are a \( \sigma \)-algebra on \( \Omega \) and a probability measure on \( \Omega \), respectively. The gradient operator \( \nabla \) only acts on the physical domain \( D \). We begin by introducing a weak formulation of a deterministic problem of \( (2.1) \), which arises from sampling an elementary event \( \omega^{(k)} \).
from the probability space $\Omega$: Find $u(x, \omega^{(k)}) \in H^1_0(D)$ such that

$$\int_D a(x, \omega^{(k)})\nabla u(x, \omega^{(k)}) \cdot \nabla v(x) dx = \int_D f v(x), \quad \forall v(x) \in H^1_0(D). \quad (2.2)$$

The stochastic Galerkin method seeks a solution satisfying an “extended” weak formulation of (2.1): Find $u(x, \xi)$ in $V = H^1_0(D) \otimes L^2(\Omega)$ such that

$$\left< \int_D a(x, \omega)\nabla u(x, \omega) \cdot \nabla v(x, \omega) dx \right> = \left< \int_D f v(x, \omega) \right>, \quad \forall v(x, \omega) \in V \quad (2.3)$$

where $\langle \cdot \rangle$ refers to expected value with respect to the probability measure on $L^2(\Omega)$ and $V$ is equipped with the gradient norm

$$\|v\|^2_V = \int_\Omega \int_D a(x, \omega)|\nabla v(x, \omega)|^2 dx dP(\omega). \quad (2.4)$$

If $a(x, \omega)$ is bounded and uniformly positive,

$$0 < a_{\min} \leq a(x, \omega) \leq a_{\max} < +\infty, \quad \text{a.e. in } D \times \Omega, \quad (2.5)$$

then the Lax-Milgram lemma can be applied to establish existence and uniqueness of a solution $u(x, \omega) \in V$ of the variational problem (2.3). The gradient norm is also called an energy norm [3]. It has been shown that the solution error in the energy norm is minimized by the stochastic Galerkin solution [3,73] as in this example, the underlying PDE operator is self-adjoint and coercive.

For the uncertain diffusivity $a(x, \omega)$, we consider a spectral representation of
the random field using gPC expansion,

\[ a(x, \xi(\omega)) = \sum_{i=0}^{\infty} a_i(x) \psi_i(\xi(\omega)), \quad (2.6) \]

where \( \xi(\omega) = \{\xi_1(\omega), \ldots, \xi_M(\omega)\} \) is an \( M \)-dimensional random variable with joint probability density function \( \rho(\xi) \). We assume that the random variables are independent and identically distributed and the stochastic domain is denoted by \( \Gamma = \prod_{i=1}^{M} \Gamma_i \) (i.e., the joint image of \( \xi \) where \( \xi_i : \Omega \to \Gamma_i \). Here, \( \{\psi_i(\xi)\} \) is an orthogonal gPC basis, for which the details will be introduced in Section 2.2. In the sequel, we denote the random diffusivity by \( a(x, \xi) \) as the random diffusivity is parameterized with a set of random variables \( \xi \).

For simplifying a derivation of the stochastic Galerkin system, we consider a special case of the random field expansion (2.6) where the expansion consists of polynomials with degree \( \leq 1 \). Such random field can be simulated by using Karhunen-Loève expansion [67] or considering a piecewise constant random field. Note that the derivation of the stochastic Galerkin system with a general random field expansion (2.6) is a straightforward extension of the derivation described in the following. For the discussion in this chapter, we consider a truncated Karhunen-Loève expansion [67],

\[ a(x, \omega) = a_0 + \sigma \sum_{i=1}^{M} \sqrt{\lambda_i} a_i(x) \xi_i(\omega), \quad (2.7) \]

where \( a_0 \) and \( \sigma^2 \) are the mean and variance of the random field, respectively, and
$(\lambda_i, a_i)$ is an eigenpair of the covariance kernel of the random field, $C(x, y)$. That is, eigenpairs consist of the solutions of the eigenproblem of the integral operator $C$: $L^2(D) \rightarrow L^2(D)$,

$$(Cu)(x) = \int_D c(x, y)u(y)dy, \quad (Ca_m)(x) = \lambda_m a_m(x),$$

If the random field $a(x, \xi)$ is parameterized by a finite number of random variables, then the solution $u(x, \omega)$ can be described by this same set of random variables by Doob-Dynkin’s Lemma [86] (i.e., $u(x, \omega) \approx u(x, \xi_1, \ldots, \xi_M)$).

### 2.2 Discretization

The discrete stochastic Galerkin method employs a standard approximation in the spatial domain and a polynomial approximation in the probability domain [1, 3, 46]. The stochastic Galerkin method seeks a finite-dimensional solution $u_{hp}(x, \xi) \in W^h = X_h \otimes S_M$ such that

$$\left\langle \int_D a(x, \xi)\nabla u_{hp}(x, \xi) \cdot \nabla v(x, \xi) dx \right\rangle_{\rho} = \left\langle \int_D f v(x, \xi) \right\rangle_{\rho} \quad v(x, \xi) \in W^h \quad (2.8)$$

where $X_h$ and $S_M$ are finite-dimensional subspaces of $H^1_0(D)$ and $L^2_\rho(\Gamma)$,

$$X_h = \text{span}\{\phi_r(x)\}_{r=1}^{n_x} \subset H^1_0(D), \quad (2.9)$$

$$S_M = \text{span}\{\psi_s(\xi)\}_{s=1}^{n_\xi} \subset L^2_\rho(\Gamma), \quad (2.10)$$
and
\[ u_{hp}(x, \xi) = \sum_{s=1}^{n_\xi} \sum_{r=1}^{n_x} u_{r,s}(x) \phi_r(x) \psi_s(\xi). \] (2.11)

Here, \( \{\phi_r\} \) is a set of standard finite element basis functions and \( \{\psi_s\} \) is a set of basis functions for the generalized polynomial chaos (gPC) expansion [112] consisting of products of orthonormal univariate polynomials: \( \psi_s(\xi) = \psi_{\alpha(s)}(\xi) = \prod_{i=1}^{M} \pi_{\alpha_i(s)}(\xi_i) \) where \( \{\pi_{\alpha_i(s)}(\xi_i)\}_{i=1}^{M} \) is a set of univariate polynomials and \( \alpha(s) = (\alpha_1(s), \ldots, \alpha_M(s)) \in \mathbb{N}_0^M \) is a multi-index, where \( \alpha_i \) represents the degree of a polynomial in \( \xi_i \). The univariate polynomials \( \{\pi_{\alpha_i(s)}(\xi_i)\}_{i=1}^{M} \) are orthonormal with respect to underlying probability density functions \( \{\rho_i(\xi_i)\}_{i=1}^{M} \),
\[
\int_{\Gamma_i} \pi_k(\xi_i) \pi_l(\xi_i) \rho(\xi_i) d\xi_i = \kappa_i \delta_{kl}, \quad k, l \in \mathbb{N}_0, \ i = 1, \ldots, M
\]
where \( \delta_{kl} = 1 \) if \( k = l \) and 0 otherwise. Due to the orthonormality of the univariate polynomials \( \{\pi_{\alpha_i(s)}(\xi_i)\}_{i=1}^{M} \) and the independence among the random variables, the stochastic basis functions \( \{\psi_s\} \) are orthonormal with respect to the joint probability density function \( \rho(\xi) = \prod_{i=1}^{M} \rho_i(\xi_i) \),
\[
\int_{\Gamma} \psi_k(\xi) \psi_l(\xi) \rho(\xi) d\xi = \prod_{i=1}^{M} \int_{\Gamma_i} \pi_{\alpha_i(k)}(\xi_i) \pi_{\alpha_i(l)}(\xi_i) \rho(\xi_i) d\xi_i = \prod_{i=1}^{M} \delta_{\alpha_i(k)\alpha_i(l)}.
\]
If \( \rho \) is the density function corresponding to \( M \)-variate uniform distribution, \( \psi_s \) is a product of \( M \) univariate Legendre polynomials. Table 2.1 lists different types of probability measures (and probability density functions) and the types of gPC basis polynomials associated with those probability measures.
Table 2.1: Probability distribution and the type of gPC basis.

<table>
<thead>
<tr>
<th>Probability distribution</th>
<th>pdf</th>
<th>gPC basis polynomial</th>
<th>Support</th>
</tr>
</thead>
<tbody>
<tr>
<td>Normal</td>
<td>$\frac{1}{\sqrt{2\pi}} \exp\left(-\frac{\xi^2}{2}\right)$</td>
<td>Hermite</td>
<td>$[-\infty, \infty]$</td>
</tr>
<tr>
<td>Uniform</td>
<td>$\frac{1}{2}$</td>
<td>Legendre</td>
<td>$[0, 1]$</td>
</tr>
<tr>
<td>Exponential</td>
<td>$\exp(-\xi)$</td>
<td>Laguerre</td>
<td>$[0, \infty]$</td>
</tr>
<tr>
<td>Gamma($\alpha$)</td>
<td>$\frac{\Gamma(\alpha+1)}{(\alpha+\xi^2)^{\frac{\alpha}{2}}}$</td>
<td>Generalized Laguerre</td>
<td>$[0, \infty]$</td>
</tr>
<tr>
<td>Beta($\alpha, \beta$)</td>
<td>$\frac{\Gamma(\alpha+\xi^2)^{\frac{\beta}{2}}}{2^{\frac{\alpha+\beta}{2}} \bar{\Gamma}(\alpha+\beta+1)}$</td>
<td>Jacobi</td>
<td>$[0, \infty]$</td>
</tr>
</tbody>
</table>

Once the type of the gPC basis polynomials is chosen, the finite-dimensional polynomial space, $S_M = \text{span}\{\psi_s(\xi)\}_{s=1}^{n_\xi}$, can be constructed. The most naive approach in constructing $S_M$ is called “Tensor Product (TP) space,” for which the multi-index set can be defined as

$$\Lambda_{M,p}^{\text{TP}} = \{\alpha(s) \in \mathbb{N}_0^M : \max\{\alpha_1(s), \ldots, \alpha_M(s)\} \leq p\}. \quad (2.12)$$

Although the TP space is easy to construct, the cardinality of the set $\Lambda_{M,p}^{\text{TP}}$ is $M^p$, which increases exponentially as the maximum polynomial degree $p$ increases. Instead, in this study, we set $\Lambda_M$ to be the Total Degree (TD) space $\Lambda_{M,p}^{\text{TD}}$, given by

$$\Lambda_{M,p}^{\text{TD}} = \{\alpha(s) \in \mathbb{N}_0^M : \|\alpha(s)\|_1 \leq p\} \quad (2.13)$$

where $\mathbb{N}_0^M$ is the set of non-negative integers, $\|\alpha(s)\|_1 = \sum_{k=1}^{M} \alpha_k(s)$, and $p$ defines the maximal degree of $\{\psi_i\}_{i=1}^{n_\xi}$. Then, the number of gPC basis functions is $n_\xi = \dim(\Lambda_{M,p}) = \frac{(M+p)!}{M!p!}$. The TD space has been known to be very effective in approximating the solutions of many stochastic PDEs. In particular, if the stochastic diffusion equations with the random field of the form (2.7) is considered, the TD
space is known to provide the best N-term approximation for the solutions \([23,24]\).

In \([7,8]\), the decay of Legendre coefficients for the solutions of the elliptic stochastic PDEs is studied and the TD space has been shown to be quasi-optimal. Thus, in this thesis, we use the TD space as the finite-dimensional approximation space in the stochastic domain. An example of the TD space with \(M = 2\) and \(p = 3\) is

\[
\Lambda_{2,3}^{TD} = \{(0, 0), (0, 1), (0, 2), (0, 3), (1, 0), (1, 1), (2, 0), (2, 1), (3, 0)\},
\]

which is lexicographically ordered, and the cardinality of the space is \(n_\xi = \frac{(2+3)!}{2!3!} = 10\).

If the coefficients of \((2.11)\) are ordered by grouping spatial indices together as

\[
u_{11}, u_{21}, \ldots, u_{n_x1}, u_{12}, u_{22}, \ldots, u_{n_x2}, u_{13}, \ldots, u_{n_xn_\xi}, \tag{2.14}\]

then, it follows from \((2.8)\) and \((2.11)\) that the Galerkin system

\[
Au = f \tag{2.15}
\]

can be represented using Kronecker-product notation \([81]\),

\[
\left( G_0 \otimes K_0 + \sum_{l=1}^{M} G_l \otimes K_l \right) u = g_0 \otimes f_0, \tag{2.16}
\]

where the Kronecker product between two matrices \(G \in \mathbb{R}^{n_\xi \times n_\xi}\) and \(K \in \mathbb{R}^{n_x \times n_x}\) is
defined as follows:

\[
G \otimes K = \begin{bmatrix}
[G]_{11}K & \ldots & [G]_{1n_\xi}K \\
\vdots & \ddots & \vdots \\
[G]_{n_\xi1}K & \ldots & [G]_{n_\xi n_\xi}K
\end{bmatrix},
\]

\(K_i\) refers to the \(i\)th weighted stiffness matrix defined via

\[
[K_0]_{ij} = \int_D a_0 \nabla \phi_i(x) \nabla \phi_j(x) dx,
\]

\[
[K_l]_{ij} = \sigma \sqrt{\lambda_l} \int_D a_l(x) \nabla \phi_i(x) \nabla \phi_j(x) dx, \quad l = 1, \ldots, M,
\]

(2.17)

\(G_i\) refers to the \(i\)th “stochastic” matrices defined via

\[
[G_0]_{ij} = \langle \psi_i(\xi) \psi_j(\xi) \rangle_\rho,
\]

\[
[G_l]_{ij} = \langle \xi_l \psi_i(\xi) \psi_j(\xi) \rangle_\rho \quad l = 1, \ldots, M,
\]

(2.18)

and the vectors \(f_0\) and \(g_0\) are defined via

\[
[f_0]_i = \int_D f \phi_i(x) dx,
\]

\[
[g_0]_i = \langle \psi_i(\xi) \rangle_\rho.
\]

(2.19)

Note that \(\{G_l\}_{l=1}^M\) of (2.18) are highly sparse because of the orthogonality properties of the stochastic basis functions \([41]\).

The global Galerkin system shown in (2.16) is of order \(n_x n_\xi\), which becomes very large if the solution is sought on a fine spatial grid (i.e., large \(n_x\)) and a high-dimensional stochastic space (i.e., large \(M\) or \(p\) and, consequently, large \(n_\xi\)). The
Kronecker-product structure, however, leads to a block-sparse matrix, where the block nonzero structure of the matrix follows the nonzero structure of the stochastic matrices \( \{G_i\}_{i=0}^M \). Figure 2.1 depicts the block nonzero structure of the Galerkin matrix of order \( n_x n_\xi \) where \( n_\xi = 56 \) by setting \( M = 5 \) and \( p = 3 \) (i.e., \( 56 = (5+3)! / 5!3! \)), and each square in the figure represent a weighted stiffness matrix of order \( n_x \).

With this block sparse structure, it is natural to consider development of sparse linear solvers for use with the stochastic Galerkin methods, which will be discussed in the next section.

![Block nonzero structure of the Galerkin matrix.](image)

Figure 2.1: Block nonzero structure of the Galerkin matrix.

### 2.3 Iterative solvers for stochastic Galerkin systems

As for deterministic PDE problems, use of Krylov subspace methods has been very successful for stochastic PDE problems. Here, we review briefly review Krylov subspace methods and some notable results concerning Krylov subspace methods.
in the context of the stochastic Galerkin method. A Krylov subspace method seeks an approximate solution of a linear system $Ax = b$ on an affine subspace $x_0 + \mathcal{K}_m$, where $\mathcal{K}_m$ is the $m$-dimensional Krylov subspace

$$
\mathcal{K}_m(A, r_0) = \text{span}\{r_0, Ar_0, \ldots, A^{m-1}r_0\}.
$$

Here, $x_0$ denotes a possibly arbitrary initial iterate for an approximate solution and $r_0 = b - Ax_0$ denotes the initial residual. An approximate solution $x_m$ can be found by employing an orthogonal projection of the residual $r_m = b - Ax_m$ onto another $m$-dimensional subspace $\mathcal{L}_m (r_m \perp \mathcal{L}_m)$. There are two well-known projection techniques: the Galerkin projection technique with $\mathcal{L}_m = \mathcal{K}_m$ and the Petrov–Galerkin projection technique $\mathcal{L}_m = A\mathcal{K}_m$. Such projection techniques give rise to effective Krylov subspace methods. The Galerkin projection technique characterizes the Conjugate Gradient (CG) method [51] for a symmetric positive definite $A$. The Petrov–Galerkin projection technique characterizes the minimum residual method (MINRES) [78] for a nonsingular symmetric indefinite $A$ and the generalized minimum residual method (GMRES) [92] for a general nonsingular $A$.

An initial attempt to solve the global Galerkin system (2.16) using Krylov subspace solution methods can be found in [80], followed by more advanced studies on iterative methods for the Galerkin system [33,34,56]. In those studies, an efficient structure-aware matrix-vector product exploiting the block sparse structure has been studied. For faster convergence of the iterative methods, a preconditioned system
is considered,
\[ M^{-1}A = M^{-1}f \]  
(2.20)

where \( M \) is the preconditioner. Or, alternatively,
\[ AM^{-1}\tilde{u} = f, \quad u = M^{-1}\tilde{u}. \]  
(2.21)

A Krylov subspace method constructs an \( m \)-dimensional preconditioned Krylov subspace \( \{ r_0, M^{-1}Ar_0, \ldots, M^{-1}A^{m-1}r_0 \} \) for a left preconditioned system (2.20) and \( \{ r_0, AM^{-1}r_0, \ldots, A^{m-1}M^{-1}r_0 \} \) for a right preconditioned system (2.21).

The use of this preconditioner \( M \) in Krylov subspace methods requires an application of the action of its inverse \( M^{-1} \), or approximating it.

In those early studies, a preconditioned conjugate gradient method [51] with a simple block-diagonal preconditioning strategy, which incorporates the mean component of the random field is widely used, i.e.,
\[ M = G_0 \otimes K_0. \]  
(2.22)

For the efficient application of the preconditioner, an incomplete Cholesky factor as a preconditioner has been considered [45, 80].

In more recent work [81], a preconditioned CG method with a black-box Algebraic Multigrid (AMG) [90] preconditioner was considered, where the action of \( K_0^{-1} \) is replaced by the V-cycle of the black-box AMG. In [81], it has been shown that an eigenvalue bound of the preconditioned Galerkin system associated with the normal
and the uniform random variables is independent of the spatial discretization parameter and $p$ for bounded random variables such as ones with uniform distribution, but it depends on the variance of the random field $\sigma^2$ in (2.7). Also, it has been shown that the black-box AMG preconditioner is robust with respect to the spatial discretization parameter and requires less memory than the factorization method for fine spatial meshes.

So far, the preconditioning strategy only takes into account the matrix associated with the mean coefficient of the random field (i.e., $G_0 \otimes K_0$). The mean-based preconditioner may not be effective if the variance of the random field becomes large compared to the mean. To resolve this issue, a preconditioner proposed in [105] incorporates the entire information in the global Galerkin matrix (2.16). Inspired by the work of [66], the new preconditioner is constructed by solving a minimization problem

$$\min \| A - G \otimes K_0 \|_F$$  \hspace{1cm} (2.23)

where $\| \cdot \|_F$ is a Frobenius norm, $A$ is the global Galerkin matrix, $K_0$ is the mean stiffness matrix, and $G \in \mathbb{R}^{n_\xi \times n_\xi}$ to be solved. The minimization can be solved as

$$G = \sum_{i=0}^{M} \frac{\text{trace}(K_i^T K_0)}{\text{trace}(K_0^T K_0)} G_i,$$  \hspace{1cm} (2.24)

where

$$\text{trace}(A) = \sum_{i=1}^{n_\xi n_\xi} [A]_{ii}$$

is the trace operator, which sums of diagonal entries of $A \in \mathbb{R}^{n_\xi \times n_\xi \times n_\xi}$. Combined
with this preconditioner, the conjugate gradient method performs efficiently even for the large variance stochastic diffusion coefficients.

There are other effective preconditioning strategies [36, 89, 100, 101] including a preconditioner based on matrix splitting technique (e.g., Jacobi, Gauss-Seidel) applied to the stochastic matrices [89], and a hierarchical Gauss-Seidel preconditioner and a Schur complement preconditioner and its efficient variant proposed by [100, 101], which exploit recursive hierarchical structure of the global Galerkin matrix. We also note that there are several attempts to solve the stochastic Galerkin system using the multigrid solver. Initiated by [64], the practical application and the theoretical aspects of the multigrid method for the stochastic diffusion equations have been studied in [35, 94], and further extended in [27, 89].

Another successful approach shown in [40] considers a stochastic variant of the mixed variational formulation [16, 87] to discretize the stochastic diffusion equations, which results in a symmetric and indefinite system matrix. To solve the saddle point problem, a preconditioned MINRES method [78] is considered. For a preconditioner, a mean-based Schur complement of the indefinite system is computed, where action of inverse required to apply the preconditioner is replaced with the application of AMG V-cycle. Further studies on a preconditioner for the saddle point system in the stochastic Galerkin mixed variational formulation have been conducted in [84].

When the lognormal diffusion coefficient $a(x, \xi) = \exp(g(x, \xi))$, where $g(x, \xi)$ is a Gaussian random field, is considered, the coefficient $a(x, \xi)$ is typically approximated as a finite-term gPC expansion. After discretization, the resulting systems are block dense [69, 106], which makes matrix operations required by Krylov subspace
methods expensive. To resolve this issue, in [106], a “log-transformed” reformulation [115] of the problem as a convection-diffusion problem is considered. By multiplying \( \exp(-g(x, \xi)) \) on both sides of \(- \nabla \cdot (a(x, \xi)\nabla u(x, \xi)) = f\) and algebraically manipulating the equation, one can obtain a convection-diffusion equation

\[
\Delta u(x, \xi) - \nabla g(x, \xi) \cdot \nabla u(x, \xi) = f(x) \exp(-g(x, \xi)).
\]  

\[
(2.25)
\]

The stochastic Galerkin discretization can be applied to the convection-diffusion equations, which results in a nonsymmetric system of equations. To compute solutions of the nonsymmetric system, the generalized minimum residual method [92] is used. For preconditioning, two types of a mean-based preconditioner have been used: one constructed from a diffusion term only, and one constructed from the diffusion term and a convection term associated with the mean coefficients of the random field. A mixed variational formulation of the log-transformed equations and an associated iterative solution method are studied in [107].

Those solution methods have explored various formulations of the stochastic diffusion equations and applied iterative solvers to the resulting Galerkin system, which is preconditioned by various preconditioning strategies. With the numerical experiments with benchmark problems, those methods have been shown to be efficient and effective for moderate dimensional problems. However, the computational complexity \( O(n_x n_\xi) \) grows rapidly as the problem posed on a finer spatial grid or the number of random variables parameterizing the problem increases.
Chapter 3: Low-rank approximation method for linear PDEs

3.1 Introduction

In this chapter, we present a low-rank approximation method for the stochastic linear elliptic boundary value problem: Find a random function, \( u(x, \xi) : \bar{D} \times \Gamma \rightarrow \mathbb{R} \) that satisfies

\[
\mathcal{L}(a(x, \xi))(u(x, \xi)) = f(x) \quad \text{in } D \times \Gamma,
\]  

(3.1)

where \( \mathcal{L} \) is a linear elliptic operator and \( a(x, \xi) \) is a positive random field parameterized by a set of random variables \( \xi = \{\xi_1, \ldots, \xi_M\} \). The problem is posed on a bounded domain \( D \subset \mathbb{R}^2 \) with appropriate boundary conditions.

After the stochastic Galerkin method [1,3,46], which, after discretization described in Section 2.2, leads to a large coupled deterministic system (2.16) for which computations will be expensive for large-scale applications. When the coefficient \( a(x, \xi) \) has an affine structure depending on a finite number of random variables, the system matrix \( A \) can be represented by a sum of Kronecker products of smaller matrices. Matrix operations such as matrix-vector products that take advantage of the tensor format can be performed efficiently, which makes the use of iterative solvers attractive. In this study, we develop a new efficient iterative solver for
systems represented in the Kronecker-product structure.

In recent years, many authors started to explore the Kronecker-product structure of such problems and developed iterative algorithms that exploit the structure to reduce computational efforts [6,59–61,70,83,93]. In particular, thorough use of tensor Krylov subspace methods, which operate in tensor format, have been studied. Variants of this approach have been developed for the Richardson iteration [61,70], the conjugate gradient method [61], the BiCGstab method [61], the minimum residual method [103], and the general minimum residual (GMRES) method [6]. In addition, it has been shown that the solution of (3.1) in the stochastic Galerkin setting can be approximated by a tensor of low rank, which further reduces computational effort [4,5]. If Krylov subspace methods are used to compute such a solution, however, it may happen that approximate solutions or other auxiliary terms obtained during the course of an iteration do not have low rank, and rank-reduction schemes are required to keep costs under control.

In this study, we will explore a variant of the generalized minimum residual (GMRES) method combined with a rank-reduction strategy that exploits specific features of the stochastic Galerkin formulation. The strategy we propose is a two-level scheme that first identifies a low-dimensional subspace, obtained from a coarse-grid spatial discretization, on which a low-rank coarse-grid tensor solution is computed. This solution can be used to estimate the rank of the tensor solution for the desired fine-grid solution. This information is used to define a strategy for rank reduction to be used with iteration on the fine grid space. We show that this strategy enhances the efficiency of preconditioned GMRES for computing the
solution.

The proposed method can be viewed as a dimension-reduction method as it identifies a dominant subspace and computes an approximate solution in that subspace. Other approaches developed for dimension reduction for the solutions of stochastic PDEs include reduced basis methods \([52,85]\), which construct dominant subspace associated with parameterized models using greedy search methods, and active subspace methods \([25]\), which detect a subspace of strong variability for a scalar-valued multivariate functions using gradient computations. Another model reduction approach developed in \([31]\) identifies a dominant subspace based on the covariance structure of the solution on the coarse grid and uses the subspace for the fine-grid computation. The approach developed here uses inexpensive low-rank approximation technique to construct the desired subspace on coarse-grid computations. Then the identified subspace is used to truncate tensors of high ranks in the iteration process to construct a solution on a finer spatial discretization.

An outline of the chapter is as follows. In section 3.2, we review the stochastic Galerkin method and present the Kronecker-product structure of Galerkin systems. In section 3.3, we present a preconditioned projection method for computing approximate solutions in low-rank tensor format. In section 3.4, we review the conventional approaches and propose a coarse-grid rank-reduction scheme, which is the main contribution of this work. In section 3.5, we illustrate the effectiveness of the low-rank projection method combined with the proposed truncation scheme by numerical experiments on benchmark problems. In section 3.6, we discuss the impact of truncation on solution statistics. Finally, in section 3.7, we draw some
3.2 Stochastic Galerkin formulation in tensor notation

Recall the stochastic Galerkin discretization discussed in Chapter 2, where we consider the steady-state stochastic diffusion equation with homogeneous Dirichlet boundary conditions shown in (2.1) with the diffusion coefficient \( a(x, \xi) \) parameterized by using a truncated Karhunen-Loève expansion (2.7). Here, \( \xi \) is an \( M \)-dimensional random variable with joint probability density function \( \rho(\xi) \). We let \( \Gamma = \prod_{i=1}^{M} \Gamma_i \) denote the joint image of \( \xi \), which we refer to as the stochastic domain. The expected value of a random variable \( v(\xi) \) on \( \Gamma \) is then \( \langle v(\xi) \rangle_\rho = \int_\Gamma v(\xi) \rho(\xi) d\xi \).

The stochastic Galerkin method [1, 3, 46] seeks a finite-dimensional solution \( u_{hp}(x, \xi) = \sum_{s=1}^{n_\xi} \sum_{r=1}^{n_x} u_{rs} \phi_r(x) \psi_s(\xi) \) as shown in (2.11). We consider set the Total Degree (TD) space \( \Lambda_{TD}^{M,p} : \Lambda_{TD}^{M,p} = \{ \alpha(s) \in \mathbb{N}_0^M : \|\alpha(s)\|_1 \leq p \} \) (2.13). Consequently, the number of gPC basis functions is \( n_\xi = \dim(\Lambda_{M,p}) = \frac{(M+p)!}{M!p!} \). Ordering the coefficients of (2.11) based on lexicographical order as shown in (2.14) gives the linear system \( Au = f \) of (2.16) represented in tensor product notation [81],

\[
\left( G_0 \otimes K_0 + \sum_{l=1}^{M} G_l \otimes K_l \right) u = g_0 \otimes f_0
\]

(3.2)

where \( \{K_l\}_{l=0}^{M}, \{G_l\}_{l=0}^{M}, f_0, \) and \( g_0 \) are defined in (2.17)–(2.19).

We will make use of an isomorphism between \( \mathbb{R}^{n_x n_\xi} \) and \( \mathbb{R}^{n_x \times n_\xi} \) determined...
by the operators vec(·) and mat(·): $u = \text{vec}(U)$, $U = \text{mat}(u)$ where

\begin{align}
    u &= [u_1^T, \ldots, u_{n_x}^T]^T \in \mathbb{R}^{n_x n_x} \\
    U &= [u_1, \ldots, u_{n_x}] \in \mathbb{R}^{n_x \times n_x}
\end{align}

(3.3) (3.4)

with each $u_i$ of length $n_x$. In particular, (3.2) is equivalent to its “matricized” form

\begin{align}
    \sum_{i=0}^{M} K_i U G_i^T &= f_0 g_0^T,
\end{align}

(3.5)

and $u$ and $U$ can be used interchangeably to represent a solution of the Galerkin system. A solution $u$ can be represented by a sum of vectors of Kronecker structure, or equivalently $U = \text{mat}(u)$ can be represented by a sum of rank-one matrices,

\begin{align}
    u &= \sum_{k=1}^{\kappa_u} z_k \otimes y_k \\
    \Leftrightarrow U &= \sum_{k=1}^{\kappa_u} y_k z_k^T = Y_{\kappa_u} Z_{\kappa_u}^T
\end{align}

(3.6) (3.7)

where $y_i \in \mathbb{R}^{n_x}$, $z_i \in \mathbb{R}^{n_x}$, and $Y_{\kappa_u} = [y_1, \ldots, y_{\kappa_u}] \in \mathbb{R}^{n_x \times \kappa_u}$ and $Z_{\kappa_u} = [z_1, \ldots, z_{\kappa_u}] \in \mathbb{R}^{n_x \times \kappa_u}$. A tensor of the form (3.6) is often referred to as having a canonical decomposition [21] (e.g., $x = \sum_{i=1}^{\kappa_x} \otimes_{j=1}^{d} x_i^j$ where $x \in \mathbb{R}^{n_1 \ldots n_d}$, $x_i^j \in \mathbb{R}^{n_j}$ for $i = 1, \ldots, \kappa_x$, $j = 1, \ldots, d$, and $d$ refers to the dimension of the tensor). The tensor rank $\kappa_u$ is defined as the smallest number of terms needed to represent $u$. In this study, the dimension of the tensor $u$ is two and the tensor rank $\kappa_u$ of the tensor $u$ coincides with the rank of the matrix $U$. Thus, in the sequel, we also use $\kappa_u$ to refer to the
rank of $u$. With this notation, the stochastic Galerkin solution $u_{hp}(x, \xi)$ can be represented as

$$u_{hp}(x, \xi) = \Phi(x)^T Y_{\kappa_u} Z_{\kappa_u}^T \Psi(\xi) = (Y_{\kappa_u}^T \Phi(x))^T (Z_{\kappa_u}^T \Psi(\xi))$$  \hspace{1cm} (3.8)$$

where $\Phi : D \to \mathbb{R}^{n_x}$ is given by $\Phi(x) = [\phi_1(x), \ldots, \phi_{n_x}(x)]^T$ and $\Psi : \Gamma \to \mathbb{R}^{n_\xi}$ is given by $\Psi(\xi) = [\psi_1(\xi), \ldots, \psi_{n_\xi}(\xi)]^T$. As shown in [104], (3.8) corresponds to a separated representation [11],

$$u_{hp}(x, \xi) = \sum_{i=1}^{\kappa_u} \hat{y}_i(x) \hat{z}_i(\xi),$$  \hspace{1cm} (3.9)$$

where $\hat{y}_i(x) = (\Phi(x))^T y_i$ and $\hat{z}_i(\xi) = (\Psi(\xi))^T z_i$. We will use this representation to construct a new rank-reduction operator. In the discrete model (3.8), the rank of the solution is typically $\kappa_u = \min(n_x, n_\xi)$.

In [10, 48], it was shown that the solution to (3.2) can be approximated well by a quantity $\tilde{u}$ of rank $\kappa_{\tilde{u}} \ll \min(n_x, n_\xi)$ if the system matrix and the right-hand side has Kronecker-product structure. Thus, we seek a low-rank approximation to the solution $\tilde{u}$ to (3.2) for which

$$A\tilde{u} = \left( \sum_{l=0}^{M} G_l \otimes K_l \right) \left( \sum_{k=1}^{\kappa_{\tilde{u}}} \tilde{z}_i \otimes \tilde{y}_k \right) \approx g_0 \otimes f_0.$$  \hspace{1cm} (3.10)$$
3.2.1 Basic operations in tensor notation

We point out here a feature of the basic operations required by Krylov subspace methods in the setting we are considering, where the operators and data of interest have tensor format. The $m$th step of such methods results in the Krylov subspace, 
\[ \mathcal{K}_m(A, v_1) = \text{span}\{v_1, Av_1, \ldots, A^{m-1}v_1\}, \]
which is generated using matrix-vector products and addition/subtraction of vectors.

The matrix-vector product in (3.10) can be represented as a sum of rank-one tensors by exploiting the properties of the Kronecker product,
\begin{equation}
Au = \sum_{l=0}^{M} \sum_{k=1}^{\kappa_u} G_l z_k \otimes K_l y_k = \sum_{i=1}^{(M+1)\kappa_u} \hat{z}_i \otimes \hat{y}_i.
\end{equation}

The latter expression in (3.11) suggests that in tensor notation, the matrix-vector product typically results in a vector with a higher rank. Similarly, the addition of two vectors $u$ and $v$ of rank $\kappa_u$ and $\kappa_v$ in tensor notation gives
\begin{equation}
u + v = \sum_{i=1}^{\kappa_u} z_i \otimes y_i + \sum_{j=1}^{\kappa_v} \hat{z}_j \otimes \hat{y}_j = \sum_{i=1}^{\kappa_u+\kappa_v} z_i \otimes y_i,
\end{equation}
where $y_{i+\kappa_u} = \hat{y}_i$ and $z_{i+\kappa_u} = \hat{z}_i$, $i = 1, \ldots, \kappa_v$, so that the resulting sum may have rank as large as $\kappa_u + \kappa_v$. Thus, although the goal is to find an approximate solution to (3.2) of low rank, two of the fundamental operations used in Krylov subspace methods tend to increase the rank of the quantities produced. Following [6], we will address this point in the next section.
3.3 A preconditioned projection method in tensor format

As is well known, the generalized minimum residual method (GMRES) \cite{92} constructs an approximate solution \( u_m \in u_0 + \mathcal{K}_m(A, v_1) \) where \( u_0 \) is an initial vector with residual \( r_0 = f - Au_0 \), \( v_1 = r_0/\|r_0\|_2 \), and \( \mathcal{K}_m \) is the Krylov space. This is done by generating \( V_m = [v_1, \ldots, v_m] \), where \( \{v_j\}_{j=1}^m \) is an orthogonal basis for \( \mathcal{K}_m \), and then computing \( u_m \) whose residual \( r_m \) is orthogonal to \( W_m = AV_m \). The method is shown in Algorithm 1. In this section, we discuss a variant of this method based on low-rank projection, where advantage is taken of the Kronecker format of the matrix \( A \) and the fact that we seek an approximation of \( u \) with low-rank structure.

\begin{algorithm}
1: set the initial solution \( u_0 \)
2: \( r_0 := f - Au_0 \)
3: \( \tilde{v}_1 := r_0 \)
4: \( v_1 := \tilde{v}_1/\|\tilde{v}_1\| \)
5: for \( j = 1, \ldots, m \) do
6: \( w_j := Av_j \)
7: solve \( (V_j^T V_j) \alpha = V_j^T w_j \)
8: \( \tilde{v}_{j+1} := w_j - \sum_{i=1}^j \alpha_i v_i \)
9: \( v_{j+1} := \tilde{v}_{j+1}/\|\tilde{v}_{j+1}\| \)
10: end for
11: solve \( (W_m^T A V_m) y = W_m^T r_0 \)
12: \( u_m := u_0 + V_m y \)
\end{algorithm}

3.3.1 Low-rank projection method with restarting

As we observed in Section 3.2, matrix-vector products and vector sums in tensor structure tend to increase the rank of the resulting objects. Thus, although we seek a solution of low rank, straightforward use of the GMRES method may lead to
approximate solutions of higher rank than the desired solutions. This complication can be addressed using truncation operators [6,60,61,70,93], whereby vectors of high rank are replaced by ones of low rank. The truncation is inserted into the GMRES algorithm and is interleaved with the basic operations such as matrix-vector product and addition so that the ranks of the vectors used in the algorithm are kept low.

**Algorithm 2** Restarted low-rank projection method in tensor format

1: set the initial solution \( \tilde{u}_0 \)
2: for \( k = 0, 1, \ldots \) do
3: \( r_k := f - A\tilde{u}_k \)
4: if \( \|r_k\|/\|f\| < \epsilon \) then
5: return \( \tilde{u}_k \)
6: end if
7: \( \tilde{v}_1 := \mathcal{T}_\kappa(r_k) \)
8: \( v_1 := \tilde{v}_1/\|	ilde{v}_1\| \)
9: for \( j = 1, \ldots, m \) do
10: \( w_j := Av_j \)
11: solve \((V_j^T V_j)\alpha = V_j^T w_j\)
12: \( \tilde{v}_{j+1} := \mathcal{T}_\kappa\left(w_j - \sum_{i=1}^j \alpha_i v_i\right) \)
13: \( v_{j+1} := \tilde{v}_{j+1}/\|	ilde{v}_{j+1}\| \)
14: end for
15: solve \((W_m^T A V_m)\beta = W_m^T r_k\)
16: \( \tilde{u}_{k+1} := \mathcal{T}_\kappa(\tilde{u}_k + V_m \beta) \)
17: end for

Algorithm 2 summarizes the restarted low-rank projection method in tensor format [6]. As in the standard Arnoldi iteration used by GMRES, a new vector is constructed by applying the linear operator \( A \) to the previous basis vector \( v_j \) and orthogonalizing the new basis vector \( w_j \) with respect to the previous basis vectors \( \{v_i\}_{i=1}^j \). The resulting vector is truncated to a vector \( \tilde{v}_{j+1} \) of low rank and normalized to \( v_{j+1} \), which is then added to the set of basis vectors. The truncation operator \( \mathcal{T}_\kappa \) truncates a tensor of higher rank to one of rank \( \kappa \). Thus, all the
basis vectors \( \{v_i\}_{i=1}^m \) are of the same rank, \( \kappa \). The basis vectors determine the subspace \( \mathcal{K}_m = \text{span}\{v_1, \ldots, v_m\} \), but because of truncation the basis vectors are not orthogonal and \( \mathcal{K}_m \) is not a Krylov subspace. However, it is still possible to project the residual onto the subspace \( \mathcal{W}_m = \text{span}\{w_1, \ldots, w_m\} \) to find out whether the residual can be decreased by forming a new iterate \( \tilde{u}_k + V_m\beta \). Note that all the vectors used in the entire iteration process are stored as the product of two matrices in the form like that shown in (3.7). The ranks of these vectors will be discussed below.

3.3.2 Preconditioned low-rank projection method

To speed the convergence of the projection method, we consider a right-preconditioned system:

\[
AM^{-1}\hat{u} = f, \quad \hat{u} = Mu. \tag{3.13}
\]

For the stochastic diffusion problem, we consider \( M = G_0 \otimes \tilde{K}_0 \approx G_0 \otimes K_0 \) as the preconditioner, a mean-based preconditioner [81]. For the practical application of the preconditioner, we employ algebraic multigrid methods [90], where the action of \( K_0^{-1} \) is replaced by \( \tilde{K}_0^{-1} \), an application of a single V-cycle of an algebraic multigrid method. The multigrid algorithm used point damped Jacobi smoothing with damping parameter \( .5 \) and two presmoothing and two postsmoothing steps, together with bilinear interpolation for grid transfer (as implemented in [98]). The preconditioned
matrix-vector product is then

\[
AM^{-1} \hat{u} = \sum_{l=0}^{M} \sum_{k=1}^{\kappa_l} G_l \hat{z}_k \otimes K_l \hat{y}_k, \quad \hat{u} = Mu = \sum_{i=1}^{\kappa_u} \hat{z}_i \otimes \hat{y}_i.
\]

Note that \(G_0^{-1}\) is the identity matrix because of the orthonormality of the stochastic basis functions. With right preconditioning and this preconditioner, the strategy for handling tensor rank is largely unaffected by preconditioning.

### 3.4 Truncation methods

As discussed in Section 3.3.1, in the low-rank projection method, truncation of tensors is essential for the efficient computation of approximate solutions. In this section, we discuss the conventional approach for truncation and we introduce a new coarse-grid truncation method based on a coarse-grid solution.

#### 3.4.1 Truncation based on singular values

Given a matricized vector \(U = Y_{\kappa'} Z_{\kappa'}^T\) of rank \(\kappa'\), a standard approach for truncation \([6, 70]\) is to compute the singular value decomposition (SVD) of \(U\) and compress \(U\) into an approximation of desired rank \(\kappa \ll \kappa'\). This can be done efficiently by computing QR factorizations of \(Y_{\kappa'}\) and \(Z_{\kappa'}\):

\[
Y_{\kappa'} = Q_Y R_Y \in \mathbb{R}^{n_x \times \kappa'}, \quad Z_{\kappa'} = Q_Z R_Z \in \mathbb{R}^{n_z \times \kappa'}.
\]
Then, one can compute the SVD of $R_Y R_Z^T$:

$$R_Y R_Z^T = \hat{U}_{\kappa'} \hat{\Sigma}_{\kappa'} \hat{V}_{\kappa'}^T = \sum_{k=1}^{\kappa'} \hat{\sigma}_k \hat{u}_k \hat{v}_k^T$$

and truncate the sum with $\kappa$ terms to produce

$$\hat{Y}_\kappa = Q_Y \hat{U}_\kappa \hat{\Sigma}_\kappa \in \mathbb{R}^{n_x \times \kappa}, \quad \hat{Z}_\kappa = Q_Z \hat{V}_\kappa \in \mathbb{R}^{n_\xi \times \kappa}.$$

The truncated approximation of $U$ is then $\tilde{U} = \hat{Y}_\kappa \hat{Z}_\kappa^T$. The computational complexity of the truncation is $O((n_x + n_\xi + \kappa)(\kappa')^2)$ [47], which grows quadratically with respect to $\kappa'$. In the next section, we introduce a new truncation method that avoids this computation.

### 3.4.2 Truncation based on coarse-grid rank-reduction

We now propose a coarse-grid rank-reduction strategy. We obtain insight into the rank structure of the solution using a coarse spatial grid computation. Then, we define a truncation operator based on the information obtained from this coarse-grid computation.

Let $u^c(x, \xi)$ represent a solution obtained on a coarse spatial grid (i.e., $n_x$ is small). As in (3.8), $u^c(x, \xi)$ can be represented as

$$u^c(x, \xi) = (\Phi^c(x))^T U^c \Psi(\xi) = \left((Y^c)^T \Phi^c(x)\right)^T \left((Z^c)^T \Psi(\xi)\right). \quad (3.14)$$
Here, we propose to use $Z^c$ to define a truncation operator for use in the projection method to compute a solution for the problem on a finer grid. That is, the truncation operator is defined such that, given a matricized vector $U = Y_{\kappa'}Z_{\kappa'}^T$ of rank $\kappa'$,

$$T_\kappa(U) \equiv (Y_{\kappa'}Z_{\kappa'}^T Z_{\kappa}^c)(Z_{\kappa}^c)^T = \tilde{U}$$  \hspace{1cm} (3.15)

where the resulting quantity $\tilde{U} = \tilde{Y}_\kappa \tilde{Z}_{\kappa}^T$ is of rank $\kappa$,

$$\tilde{Y}_\kappa = Y_{\kappa'}Z_{\kappa'}^T Z_{\kappa}^c \in \mathbb{R}^{n_x \times \kappa}, \hspace{0.5cm} \tilde{Z}_{\kappa} = Z_{\kappa}^c \in \mathbb{R}^{n_\xi \times \kappa}.$$  

The desired rank $\kappa$ is determined such that the relative residual $\|f^c - A^c u_{\kappa',\kappa}\|_2/\|f^c\|_2$ is smaller than a certain tolerance $\epsilon^c$ where $u_{\kappa',\kappa}$ is a $\kappa$-term approximation of $u^c$. This truncation operation requires two matrix-matrix products, and the computational complexity of truncating a vector from $\kappa'$ to $\kappa$ is $O(\kappa' \kappa (n_x + n_\xi))$. Note that with the proposed truncation strategy, the fine-grid computation is equivalent to applying GMRES to $\sum_{i=0}^{M} K_i U G_i Z_{\kappa}^c(Z_{\kappa}^c)^T = f_0 g_0^T$.

For efficient coarse-grid computation, we use the Proper Generalized Decomposition (PGD) method developed in [76, 104], which computes a separated representation of a coarse-grid solution:

$$u_{\kappa',\kappa}(x, \xi) = \sum_{i=1}^{\kappa} \tilde{y}_i(x) \tilde{z}_i(\xi).$$  \hspace{1cm} (3.16)
With the stochastic Galerkin discretization, each function can be represented as

\[
\tilde{y}_i(x) = \sum_{k=1}^{n_x} \tilde{y}_k^{(i)} \phi_k(x), \quad \tilde{z}_i(\xi) = \sum_{l=1}^{n_\xi} \tilde{z}_l^{(i)} \psi_l(\xi).
\]

As a result, as in (3.8),

\[
u^{c,\kappa}(x, \xi) = \left( (\tilde{Y}_\kappa^c)^T \Phi^c(x) \right)^T \left( (\tilde{Z}_\kappa^c)^T \Psi(\xi) \right)
\]

where \(\tilde{Y}_\kappa^c = [\tilde{y}^{(1)}, \ldots, \tilde{y}^{(\kappa)}] \in \mathbb{R}^{n_x \times \kappa}\) and \(\tilde{Z}_\kappa^c = [\tilde{z}^{(1)}, \ldots, \tilde{z}^{(\kappa)}] \in \mathbb{R}^{n_\xi \times \kappa}\) are coefficient matrices such that the \(i\)th elements of \(\tilde{y}^{(j)}\) and \(\tilde{z}^{(j)}\) are \(\tilde{y}_i^{(j)}\) and \(\tilde{z}_i^{(j)}\), respectively.

Now, the discrete solution \(U^c\) in (3.14) is approximated by \(U^{c,\kappa} = \tilde{Y}_\kappa^c (\tilde{Z}_\kappa^c)^T\), and we can obtain \(Z_\kappa^c\) by computing the SVD of \(U^{c,\kappa} = \hat{U}\hat{\Sigma}\hat{V}^T\), and, as a result, \(Z_\kappa^c = \hat{V}\).

We briefly explain how the PGD method computes a \(\kappa\)-term approximation in the next section.

### 3.4.3 Proper Generalized Decomposition method

The PGD method is a successive rank-1 approximation method. That is, the method incrementally identifies the function pairs \((\tilde{y}_i(x), \tilde{z}_i(\xi))\) of (3.16) one at a time. Once \(i\) such pairs have been computed, the next pair \((\tilde{y}_{i+1}, \tilde{z}_{i+1})\) is sought in \(X_h \times S_M\) by imposing Galerkin orthogonality with respect to the tangent manifold of the set of rank-one elements at \(\tilde{y}_{i+1}\tilde{z}_{i+1}\), which is \(\{\tilde{y}_{i+1}\zeta + \nu\tilde{z}_{i+1}; \nu \in X_h, \zeta \in S_M\}\):
find $\tilde{y}_{i+1} \tilde{z}_{i+1}$ such that $\forall (v, \zeta) \in X_h \times S_M$

$$\left\langle \int_D a(x, \xi) \nabla (u^{c,i} + \tilde{y}_{i+1} \tilde{z}_{i+1}) \cdot \nabla (\tilde{y}_{i+1} \zeta + v \tilde{z}_{i+1}) \right\rangle = \left\langle \int_D f(\tilde{y}_{i+1} \zeta + v \tilde{z}_{i+1}) \right\rangle.$$  \hfill (3.17)

It follows from (3.17) that each component of a pair $(\tilde{y}_{i+1}, \tilde{z}_{i+1})$ can be computed by solving two coupled problems: a deterministic problem (3.18) and a stochastic problem (3.19). The deterministic problem is as follows: given $\tilde{z}_{i+1}$, find $\tilde{y}_{i+1} \in X_h$ such that

$$\left\langle \int_D a(x, \xi) \nabla (u^{c,i} + \tilde{y}_{i+1} \tilde{z}_{i+1}) \cdot \nabla (\phi_j \tilde{z}_{i+1}) \right\rangle = \left\langle \int_D f \tilde{y}_{i+1} \psi_j \right\rangle, \quad j = 1, \ldots, n_x.$$  \hfill (3.18)

The first basis function $\tilde{z}_1$ can be chosen arbitrarily at the beginning of the PGD method. The finite element discretization of $u_{i+1}$ yields a linear system of order $n_x^c$. Analogously, the stochastic problem starts with $\tilde{y}_{i+1}$ and finds $\tilde{z}_{i+1} \in S_M$ such that

$$\left\langle \int_D a(x, \xi) \nabla (u^{c,i} + \tilde{y}_{i+1} \tilde{z}_{i+1}) \cdot \nabla (\tilde{y}_{i+1} \psi_j) \right\rangle = \left\langle \int_D f \tilde{y}_{i+1} \psi_j \right\rangle, \quad j = 1, \ldots, n_{\xi}.$$  \hfill (3.19)

Since $\tilde{z}_{i+1}$ is approximated by the gPC, $n_{\xi}$ unknowns have to be determined by solving a linear system of order $n_{\xi}$.

Solutions of these sets of $\kappa$ systems of order $n_x^c$ and $\kappa$ systems of order $n_{\xi}$ produce the $\kappa$-term approximation to the solution. The PGD method seeks solution
pairs until the relative residual of the computed solution satisfies a given tolerance,

\[ \| f^c - A^c u^{c,\kappa} \|_2 / \| f^c \|_2 < \epsilon^c. \] (3.20)

The accuracy of the \( \kappa \)-term approximation can also be improved by solving a set of \( \kappa \) coupled equations: given \( \{ \tilde{y}_i \}_{i=1}^{\kappa} \), find \( \{ \tilde{z}_i \}_{i=1}^{\kappa} \) such that

\[ \left\langle \int_D a(x, \xi) \nabla (u^{(\kappa)}) : \nabla (\tilde{y}_i \psi_j) \right\rangle = \left\langle \int_D f \tilde{y}_i \psi_j \right\rangle, \quad i = 1, \ldots, \kappa, \ j = 1, \ldots, n_{\xi}. \] (3.21)

This update requires the solution of a linear system of order \( \kappa n_{\xi} \). For the stochastic diffusion problems, the update problem is solved once at the end of the PGD method.

Note that the update problem could also be formulated for finding the deterministic parts \( \{ u_i \}_{i=1}^{\kappa} \) if \( n_x \ll n_{\xi} \), which requires a solution of a linear system of order \( \kappa n_x \).

With the proposed truncation strategy, Algorithm 3 summarizes the entire procedure to compute a solution on a finer grid.

**Algorithm 3** Preconditioned low-rank projection method with the coarse-grid rank-reduction

1: Compute \( u^{c,\kappa} \) that satisfies \( \| f^c - A^c u^{c,\kappa} \|_2 / \| f^c \|_2 < \epsilon^c \) using the PGD method
2: Compute \( Z^c_{\kappa} \) such that \( U^{c,\kappa} = Y^c_{\kappa} (Z^c_{\kappa})^T \) and define \( T_{\kappa}(U) \equiv (U Z^c_{\kappa}) (Z^c_{\kappa})^T \)
3: Run Algorithm 2 with \( L = AM^{-1} \), \( f_{\kappa} \) and \( T_{\kappa} \)

3.5 Numerical experiments

In this section, we present the results of numerical experiments in which the proposed iterative solver is applied to some benchmark problems. The implementa-
tion of the spatial discretization is based on the Incompressible Flow and Iterative Solver Software (IFISS) package [98]. Example problems are posed on a square domain and $\ell$ is the spatial discretization parameter (i.e., $n_x = (2^\ell + 1)^2$).

For $a(x, \xi)$ in (2.7), we consider independent random variables $\{\xi_i\}_{i=1}^M$ that are uniformly distributed over $[-\sqrt{3}, \sqrt{3}]$, $a_0 = 1$ and unless otherwise specified, $\sigma = 0.05$. As the covariance kernel, we use

$$C(x, y) = \sigma^2 \exp\left(-\frac{|x_1 - y_1|}{\gamma} - \frac{|x_2 - y_2|}{\gamma}\right)$$

(3.22)

where $\gamma$ is the correlation length. The number of terms $M$ in the truncated expansion (2.7) is determined such that 95% of the total variance is captured by $M$ terms (i.e., $(\sum_{i=1}^{M} \lambda_i) / (\sum_{i=1}^{n_x} \lambda_i) > 0.95$). We use bilinear $Q_1$ elements to generate the finite element basis and Legendre polynomials as the stochastic basis functions because the underlying random variables have a uniform distribution. The default setting of the maximal polynomial degree $p$ is 3.

3.5.1 Stochastic diffusion problem

We consider the steady-state stochastic diffusion equation in (2.8) on a domain $D = [0, 1] \times [0, 1]$ with forcing term $f(x) = 1$ and homogeneous Dirichlet boundary conditions, $u(x, \omega) = 0$ on $\partial D \times \Gamma$.

Coarse spatial grid computation. We compute $\kappa$-term approximations using the PGD method on a coarser spatial grid. Here $\ell^c$ is the refinement level for the coarse grid and $n_x^c$ is the number of degrees of freedom in the corresponding
spatial domain excluding boundary nodes. We discuss choices of coarse spatial grid in Section 3.5.3. Table 3.1 shows the rank $\kappa$ of solutions that satisfy the tolerance $\epsilon^c$ for varying correlation lengths $\gamma$ and $M$ and the computation time $t_c$. In PGD, the linear systems arising from (3.18), (3.19), and (3.21) are solved using MATLAB’s backslash operator.

Table 3.1: Rank ($\kappa$) of coarse-grid solutions satisfying $\epsilon^c$ of (3.20), and CPU time ($t_c$) for coarse-grid computation using the PGD method, for varying $\gamma$ and $M$.

<table>
<thead>
<tr>
<th>$\gamma$</th>
<th>$\epsilon^c = 10^{-5}$</th>
<th>$\epsilon^c = 10^{-6}$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$M, n_\xi$</td>
<td>4, 3, 2.5, 2</td>
<td>4, 3, 2.5, 2</td>
</tr>
<tr>
<td>$n_\xi^c(P^c)$</td>
<td>225(4), 225(4), 961(5), 961(5)</td>
<td>225(4), 225(4), 961(5), 961(5)</td>
</tr>
<tr>
<td>Rank($\kappa$)</td>
<td>25, 40, 65, 115</td>
<td>35, 65, 100, 210</td>
</tr>
<tr>
<td>CPU time($t_c$)</td>
<td>2.49, 3.47, 8.35, 45.08</td>
<td>2.93, 5.04, 14.83, 162.71</td>
</tr>
</tbody>
</table>

**Fine spatial grid computation.** With the truncation operator $T_\kappa$ (3.15) obtained from the coarse-grid solution (i.e., $Z_c^\kappa$), we solve the same stochastic diffusion problems on finer spatial grids $\ell = \{7, 8, 9\}$. For the fine-grid low-rank solutions, we use the rank $\kappa$ obtained from the coarse-grid solutions. For example, the third column of Table 3.2 shows the time required to find solutions of rank 25 satisfying the relative residual tolerance $10^{-5}$ when the number of terms in (2.7) is $M = 5$. In Algorithm 2, we set $m = 8$ (like restarted GMRES(8)). In examining performance, we identify the number of cycles, $k$, performed for the outer for-loop in Algorithm 2; this means that the number of matrix-vector products (i.e., the number of times line 10 is executed) is $mk$. Tables 3.2 and 3.3 show the number of cycles, $k$, and the computation time in seconds needed to compute approximate solutions with
\( \epsilon = 10^{-5} \) and \( 10^{-6} \), respectively, (see line 4 of Algorithm 2). Here, \( t \) is the total time and \( t_f \) excludes the time to compute the coarse-grid solution, \( t_c \). The fine-grid computation time, \( t_f \), consists of algorithm execution time and preconditioner set-up time, \( t_{\text{setup}} \).

Table 3.2: CPU times to compute low-rank solutions of the diffusion equation for \( \epsilon^c = \epsilon = 10^{-5} \) using the preconditioned low-rank projection method. Numbers of GMRES cycles are shown in parentheses.

<table>
<thead>
<tr>
<th>( n_x(\ell) )</th>
<th>( M=5 )</th>
<th>( M=7 )</th>
<th>( M=10 )</th>
<th>( M=15 )</th>
<th>( t_{\text{setup}} )</th>
</tr>
</thead>
<tbody>
<tr>
<td>129 (7)</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td>1.76</td>
</tr>
<tr>
<td>( t_f )</td>
<td>4.12 (1)</td>
<td>7.22 (1)</td>
<td>18.79 (1)</td>
<td>86.29 (1)</td>
<td></td>
</tr>
<tr>
<td>( t )</td>
<td>8.35</td>
<td>12.43</td>
<td>28.88</td>
<td>132.15</td>
<td></td>
</tr>
<tr>
<td>257 (8)</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td>10.16</td>
</tr>
<tr>
<td>( t_f )</td>
<td>12.55 (1)</td>
<td>24.70 (1)</td>
<td>74.71 (1)</td>
<td>330.45 (1)</td>
<td></td>
</tr>
<tr>
<td>( t )</td>
<td>25.17</td>
<td>38.37</td>
<td>93.20</td>
<td>385.59</td>
<td></td>
</tr>
<tr>
<td>513 (9)</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>( t_f )</td>
<td>92.83 (1)</td>
<td>102.42 (1)</td>
<td>353.07 (1)</td>
<td>2717.03 (1)</td>
<td>92.41</td>
</tr>
<tr>
<td>( t )</td>
<td>147.17</td>
<td>197.87</td>
<td>453.71</td>
<td>2854.62</td>
<td></td>
</tr>
</tbody>
</table>

Table 3.3: CPU times to compute low-rank solutions of the diffusion equation for \( \epsilon^c = \epsilon = 10^{-6} \) using the preconditioned low-rank projection method. Numbers of GMRES cycles are shown in parentheses.

<table>
<thead>
<tr>
<th>( n_x(\ell) )</th>
<th>( M=5 )</th>
<th>( M=7 )</th>
<th>( M=10 )</th>
<th>( M=15 )</th>
<th>( t_{\text{setup}} )</th>
</tr>
</thead>
<tbody>
<tr>
<td>129 (7)</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td>1.79</td>
</tr>
<tr>
<td>( t_f )</td>
<td>5.40 (1)</td>
<td>12.50 (1)</td>
<td>35.09 (1)</td>
<td>233.54 (1)</td>
<td></td>
</tr>
<tr>
<td>( t )</td>
<td>10.14</td>
<td>19.32</td>
<td>51.69</td>
<td>398.06</td>
<td></td>
</tr>
<tr>
<td>257 (8)</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td>10.53</td>
</tr>
<tr>
<td>( t_f )</td>
<td>17.23 (1)</td>
<td>46.07 (1)</td>
<td>137.19 (1)</td>
<td>1004.40 (1)</td>
<td></td>
</tr>
<tr>
<td>( t )</td>
<td>30.55</td>
<td>61.41</td>
<td>162.90</td>
<td>1177.68</td>
<td></td>
</tr>
<tr>
<td>513 (9)</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>( t_f )</td>
<td>70.37 (1)</td>
<td>217.12 (1)</td>
<td>1225.77 (1)</td>
<td>OoM</td>
<td>92.81</td>
</tr>
<tr>
<td>( t )</td>
<td>166.24</td>
<td>315.18</td>
<td>1333.63</td>
<td>OoM</td>
<td></td>
</tr>
</tbody>
</table>

The execution times show “textbook” behavior, i.e., they grow linearly with the size of the spatial grid.\(^1\) Note that the computational cost for the coarse-grid computation becomes negligible as the size of the problem becomes higher. If the

---

\(^{1}\)An exception to this statement is when both \( M \) and \( n_x \) are large. For these cases, the problem does not fit into physical memory and memory swap-in/out time dominates the execution time.
required memory for running Algorithm 2 exceeds the resources of our computing environment, solutions could not be computed and we denote these cases by OoM for “Out-of-Memory”. Table 3.4 shows the number of degrees of freedom of the fine spatial-grid problems for varying stochastic dimensions, \( M \).

In these experiments and in all those described below, we used \( \epsilon^c = \epsilon \) (the stopping tolerance specified in line 4 of Algorithm 2), and for this choice, the solver always satisfied the stopping criterion. We also tested both larger \( \epsilon^c \) and smaller \( \epsilon^c \). For \( \epsilon^c > \epsilon \), the solver sometimes failed to satisfy the stopping criterion. For \( \epsilon^c < \epsilon \), the solver was robust but consistently more expensive.

Table 3.4: Number of degrees of freedom of the fine-grid discretizations with \( p = 3 \), for varying spatial-grid refinement level, \( \ell \), and number of random variables, \( M \).

<table>
<thead>
<tr>
<th>( \ell )</th>
<th>( M=5 )</th>
<th>( M=7 )</th>
<th>( M=10 )</th>
<th>( M=15 )</th>
</tr>
</thead>
<tbody>
<tr>
<td>7</td>
<td>931,896</td>
<td>1,996,920</td>
<td>4,759,326</td>
<td>13,579,056</td>
</tr>
<tr>
<td>8</td>
<td>3,698,744</td>
<td>7,925,880</td>
<td>18,890,014</td>
<td>53,895,984</td>
</tr>
<tr>
<td>9</td>
<td>14,737,464</td>
<td>31,580,280</td>
<td>75,266,334</td>
<td>214,745,904</td>
</tr>
</tbody>
</table>

**Example problems with varying \( \sigma \) and \( p \).** We examine the rank structure of the numerical solutions of the stochastic diffusion problems and assess the performance of the proposed solution algorithm for different values of maximal degree of stochastic polynomial, \( p \) in (2.13), and variance \( \sigma^2 \) of the random field \( a(x, \xi) \).

As in the previous numerical experiments, we first identify the rank structure and define the truncation operator from coarse-grid computation. Then, we solve the same problems on a finer grid by using the proposed low-rank projection method with the coarse-grid rank-reduction scheme.
Table 3.5 shows the computation time needed to compute approximate solutions of the stochastic diffusion problems with $M = 7$ for varying maximal polynomial degree $p$. The required ranks of the approximate solutions are not affected by the number of terms in the polynomial expansion. However, the computation time is increased for the polynomial expansion with higher maximal polynomial degree because the size of $\{G_i\}_{i=0}^M$ and the size of the stochastic part of the solution gets larger as the number of terms in the gPC is increased.

Table 3.5: CPU times $t$ to compute low-rank solutions of the diffusion equation for $\epsilon^c = \epsilon = 10^{-5}$ and $10^{-6}$ using the preconditioned low-rank projection method for varying maximal polynomial degree $p$ (stochastic dofs, $n_\xi$, in the parenthesis).

<table>
<thead>
<tr>
<th>$n_x(\ell)$</th>
<th>$\epsilon^c = \epsilon = 10^{-5}$ ($\kappa = 40$)</th>
<th>$\epsilon^c = \epsilon = 10^{-6}$ ($\kappa = 65$)</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>$p = 3$ (120)</td>
<td>$p = 4$ (330)</td>
</tr>
<tr>
<td>129$^2$(7)</td>
<td>12.43</td>
<td>15.55</td>
</tr>
<tr>
<td>257$^2$(8)</td>
<td>38.37</td>
<td>44.27</td>
</tr>
<tr>
<td>513$^2$(9)</td>
<td>197.87</td>
<td>217.38</td>
</tr>
</tbody>
</table>

Table 3.6 shows the computation time $t$ needed to compute approximate solutions of the stochastic diffusion problems that satisfy the tolerance $10^{-5}$ and $10^{-6}$ for varying variance, $\sigma^2$. In general, the example problem with a larger variance requires a higher rank to satisfy the stopping tolerance, which, therefore, requires more computational effort.

Comparison to a truncation operator based on singular values. We compare the performance of the proposed solver to the preconditioned low-rank projection method combined with the conventional truncation operator from [61]. Table 3.7 shows the computation time required to compute approximate solutions
Table 3.6: CPU times \( t \) and rank \( \kappa \) to compute low-rank solutions of the diffusion equation for \( \epsilon^c = \epsilon = 10^{-5} \) and \( 10^{-6} \) using the preconditioned low-rank projection method for varying \( \sigma \).

<table>
<thead>
<tr>
<th>( \sigma )</th>
<th>( n_x )</th>
<th>( \epsilon = 10^{-5} )</th>
<th>( \epsilon = 10^{-6} )</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td></td>
<td>( M=5 )</td>
<td>( M=7 )</td>
</tr>
<tr>
<td>0.01</td>
<td></td>
<td>( \kappa = 15 )</td>
<td>( \kappa = 20 )</td>
</tr>
<tr>
<td></td>
<td>129 (^2)</td>
<td>7.28</td>
<td>8.65</td>
</tr>
<tr>
<td></td>
<td>257 (^2)</td>
<td>21.47</td>
<td>26.08</td>
</tr>
<tr>
<td></td>
<td>513 (^2)</td>
<td>130.93</td>
<td>150.85</td>
</tr>
<tr>
<td>0.05</td>
<td></td>
<td>( \kappa = 25 )</td>
<td>( \kappa = 40 )</td>
</tr>
<tr>
<td></td>
<td>129 (^2)</td>
<td>8.35</td>
<td>12.43</td>
</tr>
<tr>
<td></td>
<td>257 (^2)</td>
<td>25.17</td>
<td>38.37</td>
</tr>
<tr>
<td></td>
<td>513 (^2)</td>
<td>147.17</td>
<td>197.87</td>
</tr>
<tr>
<td>0.1</td>
<td></td>
<td>( \kappa = 35 )</td>
<td>( \kappa = 60 )</td>
</tr>
<tr>
<td></td>
<td>129 (^2)</td>
<td>9.78</td>
<td>17.24</td>
</tr>
<tr>
<td></td>
<td>257 (^2)</td>
<td>29.98</td>
<td>54.94</td>
</tr>
<tr>
<td></td>
<td>513 (^2)</td>
<td>164.48</td>
<td>273.33</td>
</tr>
</tbody>
</table>

using the conventional and new truncation strategies. The total computation time, \( t \), of the low-rank projection method with the coarse-grid rank reduction includes both coarse-grid, \( t_c \), and fine-grid computations, \( t_f \). The low-rank projection method with the SVD-based truncation operator, which is implemented based on [6], does not require a coarse-grid computation and can start with any arbitrary initial guess for rank, \( \kappa \). For these computations, we used the values of rank identified in the coarse-grid computations, which are illustrated in Table 3.1, for the initial rank.
Table 3.7: CPU times to compute low-rank solutions of the diffusion equation for $\epsilon^c = \epsilon = 10^{-5}$ and $10^{-6}$ using the preconditioned low-rank projection (LRP) methods with the coarse-grid rank-reduction and the singular value based truncation on the level 8 spatial grid (i.e., $n_x = 257^2$).

<table>
<thead>
<tr>
<th>$\epsilon = 10^{-5}$</th>
<th>Solver</th>
<th>$M=5$</th>
<th>$M=7$</th>
<th>$M=10$</th>
<th>$M=15$</th>
<th>$M=20$</th>
</tr>
</thead>
<tbody>
<tr>
<td>LRP-SVD</td>
<td>$t_{SVD}$</td>
<td>55.04</td>
<td>108.11</td>
<td>284.27</td>
<td>1280.65</td>
<td>5691.19</td>
</tr>
<tr>
<td>LRP-Coarse</td>
<td>$t$</td>
<td>25.17</td>
<td>38.37</td>
<td>93.20</td>
<td>385.59</td>
<td>1943.49</td>
</tr>
<tr>
<td>$\epsilon = 10^{-6}$</td>
<td>LRP-SVD</td>
<td>$t_{SVD}$</td>
<td>76.03</td>
<td>198.20</td>
<td>564.12</td>
<td>5131.32</td>
</tr>
<tr>
<td>LRP-Coarse</td>
<td>$t$</td>
<td>30.55</td>
<td>61.41</td>
<td>162.90</td>
<td>1177.68</td>
<td>OoM</td>
</tr>
</tbody>
</table>

Table 3.8: CPU times to compute low-rank solutions of the diffusion equation for $\epsilon^c = \epsilon = 10^{-5}$ and $10^{-6}$ using the PGD method and the preconditioned low-rank projection methods on the level 8 spatial grid (i.e., $n_x = 257^2$).

<table>
<thead>
<tr>
<th>$\epsilon = 10^{-5}$</th>
<th>Solver</th>
<th>$M = 5$</th>
<th>$M = 7$</th>
<th>$M = 10$</th>
<th>$M = 15$</th>
<th>$M = 20$</th>
</tr>
</thead>
<tbody>
<tr>
<td>PGD</td>
<td>$\kappa$</td>
<td>25</td>
<td>45</td>
<td>65</td>
<td>125</td>
<td>195</td>
</tr>
<tr>
<td></td>
<td>$t$</td>
<td>43.78</td>
<td>109.72</td>
<td>228.73</td>
<td>940.69</td>
<td>3066.87</td>
</tr>
<tr>
<td>LRP-Coarse</td>
<td>$\kappa$</td>
<td>25</td>
<td>40</td>
<td>65</td>
<td>115</td>
<td>180</td>
</tr>
<tr>
<td></td>
<td>$t$</td>
<td>25.17</td>
<td>38.37</td>
<td>93.20</td>
<td>385.59</td>
<td>1943.49</td>
</tr>
<tr>
<td>$\epsilon = 10^{-6}$</td>
<td>PGD</td>
<td>$\kappa$</td>
<td>40</td>
<td>70</td>
<td>110</td>
<td>225</td>
</tr>
<tr>
<td></td>
<td>$t$</td>
<td>74.43</td>
<td>214.82</td>
<td>533.10</td>
<td>2713.70</td>
<td>OoM</td>
</tr>
<tr>
<td>LRP-Coarse</td>
<td>$\kappa$</td>
<td>35</td>
<td>65</td>
<td>100</td>
<td>210</td>
<td>OoM</td>
</tr>
<tr>
<td></td>
<td>$t$</td>
<td>30.55</td>
<td>61.41</td>
<td>162.90</td>
<td>1177.68</td>
<td>OoM</td>
</tr>
</tbody>
</table>

**PGD as a solver on a finer spatial grid.** The PGD method could be applied directly to the fine-grid problems. We assess the performance of the PGD method for computing fine-grid solutions in Table 3.8, which shows the rank and computation time for computing approximate solutions that satisfy the tolerance $10^{-5}$ and $10^{-6}$ using PGD on a finer spatial grid. For the low-rank projection method, we record total computation time, $t$, which includes coarse-grid computation, $t_c$, AMG preconditioner set-up, $t_{setup}$, and fine-grid computation time, $t_f$. 

47
We compare the rank and the computation time for computing solutions using the PGD method and the proposed projection method. The proposed low-rank projection method runs faster and requires somewhat smaller ranks than the PGD method.

**Remark.** We also tested the techniques compared in Tables 3.7 and 3.8 for different values of $\sigma$, $\sigma = 0.01$ and 0.1, with similar results. Indeed, the performance of LRP-Coarse is more favorable for the larger value $\sigma = 0.1$.

### 3.5.2 Stochastic convection-diffusion problem

For a second benchmark problem, we consider the steady-state convection-diffusion equation defined on $D = [-1, 1] \times [-1, 1]$ with non-homogeneous Dirichlet boundary conditions, constant vertical wind $\bar{w} = (0, 1)$, and $f = 0$,

$$
\begin{align*}
\nu \nabla \cdot (a(x, \xi) \nabla u(x, \xi)) + \bar{w} \cdot \nabla u(x, \xi) &= f(x, \xi) \quad \text{in } D \times \Gamma, \\
u \nabla \cdot (a(x, \xi) \nabla u(x, \xi)) + \bar{w} \cdot \nabla u(x, \xi) &= f(x, \xi) \quad \text{in } D \times \Gamma, \\
\end{align*}
$$

where $g_D(x)$ is determined by

$$
g_D(x) = \begin{cases} 
g_D(x, -1) = x, & g_D(x, 1) = 0, \\
g_D(-1, y) = -1, & g_D(1, y) = 1, 
\end{cases}
$$

where the latter two approximations hold except near $y = 1$, and $\nu$ is the viscosity parameter. We consider the convection-dominated case (i.e., $\nu < 1$) and employ the streamline-diffusion method for stabilization [17]. Here, we define the element
Peclet number

\[ P_k = \frac{||\vec{w}_k||_2h_k}{2\nu} \]  \hspace{1cm} (3.25)

where \( ||\vec{w}_k||_2 \) is the \( \ell_2 \) norm of the wind at the element centroid and \( h_k \) is a measure of the element length in the direction of the wind. Note that the solution has an exponential boundary layer near \( y = 1 \) where the value of the solution dramatically changes essentially from \(-1\) to \(0\) on the left and \(+1\) to \(0\) on the right [39]. Figure 3.1 illustrates the mean of solutions \( \langle u(x, \xi) \rangle_\rho \) computed on the level 6 spatial grid and corresponding contour plots for varying viscosity parameter, \( \nu \).

![Figure 3.1: Mean solutions and contour plots on the level 6 spatial grid for varying \( \nu \).](image)

Given \( a(x, \xi) \) in (2.7), we again discretize (3.23) using the finite element method and the gPC expansion. The result is a linear system in tensor product
\[
\left( G_0 \otimes \nu K_0 + \sum_{l=1}^{M} G_l \otimes \nu K_l + G_0 \otimes N + G_0 \otimes S \right) u = g_0 \otimes f_0 \quad (3.26)
\]

where the convection term \( N \) and the streamline-diffusion term \( S \) are given by

\[
\begin{align*}
[N]_{ij} &= \int_D \vec{w} \cdot \nabla \phi_i(x) \phi_j(x) dx, \\
[S]_{ij} &= \sum_{k=1}^{n_e} \delta_k \int_D (\vec{w} \cdot \nabla \phi_i)(\vec{w} \cdot \nabla \phi_j) d, x
\end{align*}
\]

\( n_e \) is the number of elements in the finite element discretization, and

\[
\delta_k = \begin{cases} 
\frac{h_k}{2\|\vec{w}\|_2} \left( 1 - \frac{1}{P_k} \right) & \text{if } P_k > 1, \\
0 & \text{if } P_k \leq 1 
\end{cases} \quad (3.27)
\]

As the preconditioner, we choose \( M \approx G_0 \otimes (K_0 + N + S) \) where the action of \((K_0 + N + S)^{-1}\) is replaced by application of a single V-cycle of an AMG method. In the PGD method, the non-homogeneous Dirichlet boundary condition is handled by introducing an extended affine space \([76]\): \( u^e \approx u_{bc} + u^{e,k} \) where \( u_{bc} \) is the boundary nodal functions such as \( u_{bc} = \sum_{k \in \partial D} u_{k}^{(bc)} \phi_k(x) \). For the stochastic convection-diffusion problems, the update problems \((3.21)\) need to be solved more often to compute an approximate solution of a desired accuracy with fewer terms.

**Numerical results.** To cope with the existence of the exponential boundary layer in the solution, we use vertically stretched spatial grids. We examine the performance of the low-rank projection method for varying viscosity parameter \( \nu \),
and we set \( m = 10 \) for Algorithm 2. Table 3.9 and 3.10 show \( \kappa \), \( \kappa \) computed by the PGD method, coarse-grid computation time \( t_c \), and fine grid computation time \( t_f \) to compute approximate solutions on fine spatial grids \( \ell = \{7, 8, 9\} \) satisfying \( 10^{-5} \) and \( 10^{-6} \), respectively. Underlined numbers in the spatial grid level indicates cases where streamline diffusion is not needed.

Table 3.9: CPU times to compute low-rank solutions of the convection-diffusion equation for \( \epsilon^c = \epsilon = 10^{-5} \) using the preconditioned low-rank projection methods for varying \( \nu \). Numbers of GMRES cycles are shown in parentheses.

<table>
<thead>
<tr>
<th>( \nu )</th>
<th>( \ell )</th>
<th>( M = 5 )</th>
<th>( M = 7 )</th>
<th>( M = 10 )</th>
<th>( M = 15 )</th>
<th>( t_{\text{setup}} )</th>
</tr>
</thead>
<tbody>
<tr>
<td>( \frac{1}{20} )</td>
<td>4</td>
<td>( \kappa )</td>
<td>25</td>
<td>35</td>
<td>55</td>
<td>65*</td>
</tr>
<tr>
<td></td>
<td></td>
<td>( t_c )</td>
<td>2.56</td>
<td>4.83</td>
<td>26.34</td>
<td>58.92*</td>
</tr>
<tr>
<td></td>
<td>7</td>
<td>( t_f )</td>
<td>5.73 (1)</td>
<td>9.47 (1)</td>
<td>24.86 (1)</td>
<td>72.29 (1)</td>
</tr>
<tr>
<td></td>
<td>8</td>
<td>( t_f )</td>
<td>20.52 (1)</td>
<td>36.66 (1)</td>
<td>98.72 (1)</td>
<td>248.31 (1)</td>
</tr>
<tr>
<td></td>
<td>9</td>
<td>( t_f )</td>
<td>84.55 (1)</td>
<td>152.69 (1)</td>
<td>592.63 (1)</td>
<td>1953.52 (2)</td>
</tr>
<tr>
<td>( \frac{1}{100} )</td>
<td>4</td>
<td>( \kappa )</td>
<td>20</td>
<td>25</td>
<td>45</td>
<td>55*</td>
</tr>
<tr>
<td></td>
<td></td>
<td>( t_c )</td>
<td>2.94</td>
<td>3.12</td>
<td>16.28</td>
<td>47.24*</td>
</tr>
<tr>
<td></td>
<td>7</td>
<td>( t_f )</td>
<td>5.06 (1)</td>
<td>7.28 (1)</td>
<td>18.90 (1)</td>
<td>60.66 (1)</td>
</tr>
<tr>
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<td>( t_f )</td>
<td>121.98 (2)</td>
<td>201.62 (2)</td>
<td>745.92 (2)</td>
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</tr>
<tr>
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<td>5</td>
<td>( \kappa )</td>
<td>20</td>
<td>25</td>
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</tr>
<tr>
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<td>4.79</td>
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<td>5.16 (1)</td>
<td>7.21 (1)</td>
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<td>17.57 (1)</td>
<td>25.05 (1)</td>
<td>63.56 (1)</td>
<td>175.30 (1)</td>
</tr>
<tr>
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<td>9</td>
<td>( t_f )</td>
<td>123.73 (2)</td>
<td>200.10 (2)</td>
<td>605.50 (2)</td>
<td>2568.41 (2)</td>
</tr>
<tr>
<td>( \frac{1}{400} )</td>
<td>5</td>
<td>( \kappa )</td>
<td>20</td>
<td>20</td>
<td>35</td>
<td>45*</td>
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<tr>
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<td></td>
<td>( t_c )</td>
<td>2.94</td>
<td>3.79</td>
<td>12.49</td>
<td>82.06*</td>
</tr>
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<td>7</td>
<td>( t_f )</td>
<td>8.61 (2)</td>
<td>9.84 (2)</td>
<td>26.97 (2)</td>
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<td>( \frac{1}{600} )</td>
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<td>( \kappa )</td>
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<td>20</td>
<td>35</td>
<td>45</td>
</tr>
<tr>
<td></td>
<td></td>
<td>( t_c )</td>
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<td>13.20</td>
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</tr>
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<td>( t_f )</td>
<td>31.94 (2)</td>
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</tr>
<tr>
<td></td>
<td>9</td>
<td>( t_f )</td>
<td>133.80 (2)</td>
<td>163.90 (2)</td>
<td>506.42 (2)</td>
<td>1977.83 (2)</td>
</tr>
</tbody>
</table>

When the viscosity parameter is small (i.e., \( \nu = 1/600 \)), the coarse-grid com-
putation requires the $\kappa$-term approximation on a relatively fine spatial grid (i.e., $\ell = 6$). The exponential boundary layer gets narrower as the viscosity parameter gets smaller, which requires the use of a finer spatial grid for the coarse-grid computation. If the coarse-grid computation is performed on coarser spatial grids, it fails to identify the rank structure of solutions and to yield a proper truncation operator. Analogously, when the number of terms, $M$, in the KL expansion (2.7) is large, the coarse-grid computation has to be done on a relatively fine spatial grid because the KL expansion contains more spatially oscillatory terms. In the last columns of Table 3.9 and 3.10, * and † indicate that the coarse-grid solutions are computed on the level 5 and the level 6 spatial grid, respectively.
Table 3.10: CPU times to compute approximate solutions of the convection-diffusion equation for $\epsilon = \epsilon = 10^{-6}$ using the preconditioned low-rank projection methods for varying $\nu$. Numbers of GMRES cycles are shown in parentheses.

<table>
<thead>
<tr>
<th>$\nu$</th>
<th>$\ell$</th>
<th>$M = 5$</th>
<th>$M = 7$</th>
<th>$M = 10$</th>
<th>$M = 15$</th>
<th>$t_{\text{setup}}$</th>
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<tr>
<td>1/20</td>
<td>4</td>
<td>$\kappa$</td>
<td>35</td>
<td>50</td>
<td>75</td>
<td>105*</td>
</tr>
<tr>
<td></td>
<td></td>
<td>$t_c$</td>
<td>3.31</td>
<td>9.17</td>
<td>60.51</td>
<td>194.33*</td>
</tr>
<tr>
<td></td>
<td></td>
<td>$t_f$</td>
<td>13.92 (2)</td>
<td>27.47 (2)</td>
<td>80.78 (2)</td>
<td>275.96 (2)</td>
</tr>
<tr>
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<td></td>
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<tr>
<td></td>
<td></td>
<td>$t_f$</td>
<td>52.45 (2)</td>
<td>106.11 (2)</td>
<td>311.59 (2)</td>
<td>1042.40 (2)</td>
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<td>9</td>
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<td></td>
<td></td>
<td>$t_f$</td>
<td>220.67 (2)</td>
<td>534.61 (2)</td>
<td>2694.26 (2)</td>
<td>8101.20 (2)</td>
</tr>
<tr>
<td>4/100</td>
<td>4</td>
<td>$\kappa$</td>
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<td>40</td>
<td>65</td>
<td>95*</td>
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<td>$t_c$</td>
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<td>6.25</td>
<td>38.39</td>
<td>155.83*</td>
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<td>12.34 (2)</td>
<td>21.28 (2)</td>
<td>65.02 (2)</td>
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<td>$t_f$</td>
<td>46.67 (2)</td>
<td>85.66 (2)</td>
<td>255.79 (2)</td>
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<td>273.45 (3)</td>
<td>549.82 (3)</td>
<td>3069.96 (3)</td>
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<td>1/200</td>
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<td>$\kappa$</td>
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<td>547.62 (3)</td>
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</tr>
<tr>
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<td>$\kappa$</td>
<td>25</td>
<td>35</td>
<td>55</td>
<td>75†</td>
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<td>3.49</td>
<td>6.63</td>
<td>30.50</td>
<td>151.46†</td>
</tr>
<tr>
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<td>$t_f$</td>
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<td>17.96 (2)</td>
<td>50.96 (2)</td>
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<td></td>
<td>$t_f$</td>
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</tr>
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<td>65</td>
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<td>281.27 (3)</td>
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<td>$t_f$</td>
<td>48.39 (2)</td>
<td>74.40 (2)</td>
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<td>$t_f$</td>
<td>281.27 (3)</td>
<td>462.52 (3)</td>
<td>1184.74 (3)</td>
<td>6261.34 (3)</td>
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</table>

Comparison to a truncation operator based on singular values. We again compare the performance of the proposed solver to the preconditioned low-rank projection method combined with the conventional truncation operator, the SVD-based truncation operator. Table 3.11 shows the computation time required to compute approximate solutions using the conventional and the new truncation strategy. When the low-rank projection method with SVD-based truncation oper-
ator is used, initial values for rank $\kappa$ in Algorithm 2 are obtained from coarse-grid computations of the proposed rank reduction strategy.

Table 3.11: CPU times to compute low-rank solutions of the convection-diffusion equation for $\epsilon^c = \epsilon = 10^{-5}$ and $10^{-6}$ using the preconditioned low-rank projection (LRP) methods with the coarse-grid rank-reduction and the singular value based truncation on the level 8 spatial grid (i.e., $n_x = 257^2$).

<table>
<thead>
<tr>
<th>Viscosity ($\nu$)</th>
<th>Solver</th>
<th>$M = 5$</th>
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<th>$M = 10$</th>
<th>$M = 15$</th>
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<td>$\epsilon = 10^{-5}$</td>
<td>LRP-SVD</td>
<td>$t_{\text{SVD}}$</td>
<td>68.45</td>
<td>100.83</td>
<td>201.34</td>
</tr>
<tr>
<td></td>
<td>LRP-Coarse</td>
<td>$t$</td>
<td>54.06</td>
<td>72.08</td>
<td>154.79</td>
</tr>
<tr>
<td></td>
<td>LRP-SVD</td>
<td>$t_{\text{SVD}}$</td>
<td>93.91</td>
<td>121.89</td>
<td>295.27</td>
</tr>
<tr>
<td></td>
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<td>$t$</td>
<td>55.28</td>
<td>64.36</td>
<td>125.88</td>
</tr>
<tr>
<td>$\epsilon = 10^{-6}$</td>
<td>LRP-SVD</td>
<td>$t_{\text{SVD}}$</td>
<td>90.70</td>
<td>122.56</td>
<td>251.60</td>
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<td>55.42</td>
<td>66.08</td>
<td>115.68</td>
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<tr>
<td></td>
<td>LRP-SVD</td>
<td>$t_{\text{SVD}}$</td>
<td>91.11</td>
<td>107.47</td>
<td>221.32</td>
</tr>
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<td></td>
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<td>$t$</td>
<td>69.01</td>
<td>76.63</td>
<td>158.07</td>
</tr>
<tr>
<td></td>
<td>LRP-SVD</td>
<td>$t_{\text{SVD}}$</td>
<td>90.33</td>
<td>103.44</td>
<td>218.35</td>
</tr>
<tr>
<td></td>
<td>LRP-Coarse</td>
<td>$t$</td>
<td>75.26</td>
<td>86.48</td>
<td>176.93</td>
</tr>
<tr>
<td>$\epsilon = 10^{-6}$</td>
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<td>$t_{\text{SVD}}$</td>
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<td>234.15</td>
<td>570.56</td>
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<td>145.86</td>
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<td>$t_{\text{SVD}}$</td>
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<td>471.11</td>
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<td>126.77</td>
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<td>LRP-SVD</td>
<td>$t_{\text{SVD}}$</td>
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<td>188.76</td>
<td>416.52</td>
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<td>LRP-Coarse</td>
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<td>128.96</td>
<td>293.30</td>
</tr>
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<td>97.00</td>
<td>129.87</td>
<td>234.00</td>
</tr>
</tbody>
</table>

3.5.3 Choices of coarse spatial grid

Finally, we discuss criteria for choosing the coarse grid used to generate truncation operators. The basic idea is that the coarse grid needs to be fine enough
so that important features of the problem are represented. This quality is problem dependent, and we outline what is needed for the two types of problems we examined.

First consider the diffusion equation of Section 5.1. The issue is the oscillatory nature of components of the random field \( a(x, \xi) \). In the KL expansion (2.7), the eigenpairs, \( \{(\lambda_i, a_i(x))\}_{i=1}^M \), can be obtained by solving the following integral equation,

\[
\int_D C(x, y)a_i(y)dy = \lambda_i a_i(x), \quad i = 1, \ldots, M
\]

where \( C(x, y) \) is the covariance kernel (3.22). Since the kernel is separable, the eigenfunctions of the integral problem (3.28) can be represented as

\[
a_i(x) = a_k^1(x_1)a_j^2(x_2),
\]

where \( \{a_k^1\}_{k=1}^\infty \) and \( \{a_j^2\}_{j=1}^\infty \) are the eigenfunctions of the one-dimensional integral problem (i.e., \( \int_D \exp(-|x_l - y_l|/\gamma) a_k^1(y_l)dy_l = \lambda_k^1 a_k^1(x_l), \ l = 1, 2 \)). The eigenvalues, \( \{\lambda_i\}_{i=1}^M \), are in decreasing order and \( \lambda_i \) is the \( i \)th largest value of products \( \lambda_k^1 \lambda_j^2 \) for \( k, j = 1, 2, \ldots \). Analytic expressions for the 1D eigenfunctions are given in [46] as,

\[
a_k^l(x) = \cos(\theta_k x)/\sqrt{1/2 + \sin(\theta_k)/2\theta_k}, \quad \text{for even } k,
\]

\[
a_k^{l*}(x) = \sin(\theta_k^* x)/\sqrt{1/2 - \sin(\theta_k)/2\theta_k}, \quad \text{for odd } k,
\]

where \( \theta_k \) and \( \theta_k^* \) are the solutions of

\[
\frac{1}{c} - \theta \tan\left(\frac{\theta}{2}\right) = 0 \quad \text{and} \quad \theta^* + \frac{1}{c} \tan\left(\frac{\theta^*}{2}\right) = 0,
\]

55
respectively, when the 1D integral problem is posed on \([-\frac{1}{2}, \frac{1}{2}]\). As \(i\) in the KL expansion \((2.7)\) increases, the eigenfunctions \(a_i(x)\) become more oscillatory over the spatial domain (i.e., \(\theta_k\) or \(\theta^*_k\) become larger), so that finer coarse spatial grids are required to capture the oscillatory features of the KL expansion. Table 3.12 shows the largest value of \(\{\theta_k, \theta^*_k\}\) of the eigenfunctions in the KL expansion, the half-wavelength of the functions from \((3.29)\) and our choice of coarse spatial grid refinement levels, \(\ell^c\), for different values of \(M\). With these coarse grids, there are approximately eight grid points per half wave, enough to capture the qualitative character of the wave.

We turn now to the convection-diffusion equation of Section 5.2. This problem has the same diffusion coefficient \((2.7)\) as the diffusion problem, but in addition its solution has an exponential boundary layer. In particular, for small \(\nu\), the width of the layer is smaller than the finest interval needed to represent the eigenfunctions in \((2.7)\), and in this case the coarse grid must be finer than that needed for the diffusion problem (whose solution is smooth). In Figure 3.2, the top plot illustrates the mean solutions \(\langle u(x, \xi) \rangle_{\rho}\) of the weak formulation of \((3.23)\) at \(x = 1\), which are computed on two coarse spatial grids \(\ell = \{4, 5\}\) using PGD and a fine spatial grid \(\ell = 8\).
Figure 3.2: Mean solutions $\langle u(x, \xi) \rangle_\rho$ at $x = 1$ and $y = [0.9, 1]$ illustrating the exponential boundary layer for varying spatial grid refinement level, $\ell = \{4, 5, 8\}$; (top) and lengths in $y$-direction of first few elements from $y = 1$ (bottom).

using the proposed method, with the viscosity parameter, $\nu = \frac{1}{200}$, and $M = 10$ random variables. The bottom plot shows the lengths of the first few elements in the $y$-direction near $y = 1$ for these refinement levels. If the level-4 spatial grid is used for the coarse grid computation (i.e., $\ell = 4$, red line in Figure 3.2), the width of exponential boundary layer is much narrower than the length of the smallest element and the coarse-grid solution gives a poor representation of the boundary layer. When this coarse grid is used to construct the truncation operator, the proposed scheme fails to compute an accurate approximate solution on a fine spatial grid (i.e., $\ell = 8$, black line in Figure 3.2). On the other hand, the level-5 spatial grid (i.e., $\ell = 5$, blue line in Figure 3.2) is fine enough for the coarse-grid solution to represent the character of the exponential layer, and with this coarse-grid, the resulting proposed scheme efficiently computes an accurate fine-grid solution.
Although this discussion shows that some a priori knowledge of the problem is needed to identify the coarse grid operator, in general this information is not difficult to come by. In particular, we are assuming that the expansion (2.5) is known, and it is straightforward to identify the resolution needed to represent its components, for example by examining one-dimensional cross-sections of them. If as for the second problem some knowledge of the solution is needed, this can be obtained cheaply from the solution of a deterministic problem derived from the mean of the diffusion coefficient; indeed, for the convection-diffusion problem, the boundary layer for the deterministic solution has essentially the same character as that of the stochastic solution whose mean is shown in Figure 3.2.

3.6 Statistical Computations

In this section, we explore the impact of truncation on statistical quantities associated with the solutions. In particular, we examine the mean and the variance of the solution $u_{hp}(x, \xi)$, which are defined as

$$
\mu = E[u_{hp}], \quad \sigma_u^2 = E[(u_{hp} - \mu)^2], \quad (3.30)
$$

where $E[\cdot] = \int_{\Gamma} \cdot \rho(\xi) d\xi$ refers to the expectation. Let $u_{hp}^{(\text{full})}$ refer to the discrete solution (of form (2.11)) obtained from a full-rank solution of (3.2) (i.e., with no truncation), and let $u_{hp}^{(\text{low})}$ refer to that obtained using Algorithm 3. We will examine the accuracy of $u_{hp}^{(\text{low})}$ by comparing its mean and variance to those of a reference
solution \(u_{hp}^{(\text{ref})}\) as follows:

\[
\eta_\mu \equiv \|\mu_{\text{ref}} - \mu_{\text{low}}\|_2 \leq \|\mu_{\text{ref}} - \mu_{\text{full}}\|_2 + \|\mu_{\text{full}} - \mu_{\text{low}}\|_2, \tag{3.31}
\]

\[
\eta_\sigma \equiv \|\sigma^2_{u,\text{ref}} - \sigma^2_{u,\text{low}}\|_2 \leq \|\sigma^2_{u,\text{ref}} - \sigma^2_{u,\text{full}}\|_2 + \|\sigma^2_{u,\text{full}} - \sigma^2_{u,\text{low}}\|_2, \tag{3.32}
\]

where the norm in (3.31)–(3.32) is the \(\ell_2\)-norm (e.g., \(\|\mu\|_2 = (\int_D \mu(x)^2 dx)^{1/2}\)). For these tests, \(u_{hp}^{(\text{full})}\) and \(u_{hp}^{(\text{low})}\) were computed using a fixed discretization on a spatial grid (\(\ell = 7\)) and polynomial degree \(p = 3\) for the stochastic discretization, and \(u_{hp}^{(\text{ref})}\) was computed using the larger polynomial degree \(p = 5\).\(^2\) Thus, for the means in (3.31), \(\mu_{\text{ref}} - \mu_{\text{full}}\) represents an approximate to the discretization error and \(\mu_{\text{full}} - \mu_{\text{low}}\) is the error caused by the low-rank approximation, which we refer to as the bias. Note that the mean and the variance of the stochastic Galerkin solution (2.11) can be computed easily by exploiting the orthonormality of the basis functions (i.e., for 

\[u(\xi) = \sum_{i=1}^n u_i \psi_i(\xi), \mu = u_1 E[\psi_1] = u_1\text{ and }\sigma^2_u = \sum_{i=2}^n u_i^2 E[\psi_i^2] = \sum_{i=2}^n u_i^2\]).

Figure 3.3 shows the results for various tolerances \(\epsilon^c\) and two examples of the diffusion problem (2.1) (with \(M = 5\) and \(M = 7\) in (2.7)) and one example of the convection-diffusion problem (3.23) with \(M = 5\). In all cases, it can be seen that the error for the low-rank solution is somewhat larger than the discretization error for large \(\epsilon^c\) (and this is caused by the bias), but the bias is significantly smaller than the tolerance \(\epsilon^c\). The bias is negligible for \(\epsilon^c = 10^{-7}\).

\(^2\)We also computed a more accurate reference solution with \(p = 7\) for the moderate-dimensional problem (i.e., the diffusion problem (2.1) with \(M = 5\)) and found the results to be virtually identical.
Figure 3.3: Errors in the mean and the variance of the low-rank approximate solutions shown in (3.31) and (3.32) for the stochastic diffusion problem (a)-(d) and the stochastic convection-diffusion problem (e)-(f).

3.7 Conclusion

We have studied iterative solvers for low-rank solutions of stochastic Galerkin systems of stochastic partial differential equations. In particular, we have explored
low-rank projection methods in tensor format for linear systems of Kronecker-product structure. For the computational efficiency of the projection methods, basis vectors and iterates in the projection methods are forced to have low rank, which is achieved by a coarse-grid rank-reduction strategy. We have examined the performance of this strategy with two benchmark problems: stochastic diffusion problems and stochastic convection-diffusion problems. For both problem classes, the rank structure of the solution can be identified by an inexpensive coarse-grid computation, and with the resulting coarse-grid rank-reduction strategy, the low-rank projection method is more efficient than methods for which the truncation operator is based on singular values.
4.1 Introduction

In this chapter, we present a low-rank approximation method for the steady-state Navier–Stokes equations with uncertain viscosity. Such uncertainty may arise from measurement error or uncertain ratios of multiple phases in porous media. The uncertain viscosity can be modeled as a positive random field parameterized by a set of random variables \([82, 99, 104]\) and, consequently, the solution of the stochastic Navier–Stokes equations also can be modeled as a random vector field depending on the parameters associated with the viscosity (i.e., a function of the same set of random variables). As a solution method, we consider the stochastic Galerkin method combined with the generalized polynomial chaos (gPC) expansion, which provides a spectral approximation of the solution function. The stochastic Galerkin method results in a coupled algebraic system of equations, for which computational costs may be high when the global system becomes large.

One way to address this issue is thorough use of tensor \textit{Krylov subspace} methods, which operate in tensor format and reduce the costs of matrix operations by
exploiting a Kronecker-product structure of system matrices. Variants of this approach have been developed for the Richardson iteration [61,70], the conjugate gradient method [61], the BiCGstab method [61], the minimum residual method [103], and the general minimum residual (GMRES) method [6]. Efficiencies are also obtained from the fact that solutions can often be well approximated by low-rank objects. These ideas have been shown to reduce costs for solving steady [65,70] and unsteady stochastic diffusion equations [10].

In this study, we adapt the low-rank approximation scheme to a solver for the systems of nonlinear equations obtained from the stochastic Galerkin discretization of the stochastic Navier–Stokes equations. In particular, we consider a low-rank variant of linearization schemes based on Picard and Newton iteration, where the solution of the nonlinear system is computed by solving a sequence of linearized systems using a low-rank variant of the GMRES method (lrGMRES) [6] in combination with inexact nonlinear iteration [30].

We base our development of the stochastic Galerkin formulation of the stochastic Navier–Stokes equations on ideas from [82,99]. In particular, we consider a random viscosity affinely dependent on a set of random variables as suggested in [82] (and in [99], which considers a gPC approximation of the lognormally distributed viscosity). The stochastic Galerkin formulation of the stochastic Navier–Stokes equations is also considered in [9], which studies an optimal control problem constrained by the stochastic Navier–Stokes problem and computes an approximate solution using a low-rank tensor-train decomposition [77]. Related work [104] extends a Proper Generalized Decomposition method [75] for the stochastic Navier–Stokes equations,
where a low-rank approximate solution is built from successively computing rank-one approximations. See the book [63] for an overview and other spectral approximation approaches for models of computational fluid dynamics.

An outline of the chapter is as follows. In section 4.2, we review the stochastic Navier–Stokes equations and their discrete Galerkin formulations. In section 4.3, we present an iterative low-rank approximation method for solutions of the discretized stochastic Navier–Stokes problems. In section 4.4, we introduce an efficient variant of the inexact Newton method, which solves linear systems arising in nonlinear iteration using low-rank format. We follow a hybrid approach, which employs several steps of Picard iteration followed by Newton iteration. In section 4.5, we examine the performance of the proposed method on a set of benchmark problems that model the flow over an obstacle. Finally, in section 4.6, we draw some conclusions.

4.2 Stochastic Navier–Stokes equations

Consider the stochastic Navier–Stokes equations: Find velocity \( \vec{u}(x, \xi) \) and pressure \( p(x, \xi) \) such that

\[
-\nu(x, \xi) \nabla^2 \vec{u}(x, \xi) + (\vec{u}(x, \xi) \cdot \nabla) \vec{u}(x, \xi) + \nabla p(x, \xi) = \vec{f}(x, \xi),
\]

\[
\nabla \cdot \vec{u}(x, \xi) = 0,
\]

(4.1)
in $D \times \Gamma$, with a boundary conditions

$$
\vec{u}(x, \xi) = \vec{g}(x, \xi), \quad \text{on } \partial \Omega_{\text{Dir}},
$$

$$
\nu(x, \xi) \nabla \vec{u}(x, \xi) \cdot \vec{n} - p(x, \xi) \vec{n}(x, \xi) = 0, \quad \text{on } \partial \Omega_{\text{Neu}},
$$

where $\partial \Omega = \partial \Omega_{\text{Dir}} \cup \partial \Omega_{\text{Neu}}$. The stochasticity of the equation (4.1) stems from the random viscosity $\nu(x, \xi)$, which is modeled as a positive random field parameterized by a set of independent, identically distributed random variables $\xi = \{\xi_1, \ldots, \xi_{n_{\nu}}\}$. The random variables comprising $\xi$ are defined on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ such that $\xi : \Omega \rightarrow \Gamma \subset \mathbb{R}^{n_{\nu}}$, where $\Omega$ is a sample space, $\mathcal{F}$ is a $\sigma$-algebra on $\Omega$, and $\mathbb{P}$ is a probability measure on $\Omega$. The joint probability density function of $\xi$ is denoted by $\rho(\xi)$ and the expected value of a random function $\nu(\xi)$ on $\Gamma$ is then

$$
\langle \nu \rangle_\rho = \mathbb{E}[\nu] \equiv \int_\Gamma \nu(\xi) \rho(\xi) d\xi.
$$

For the random viscosity, we consider a random field that has affine dependence on the random variables $\xi$,

$$
\nu(x, \xi) \equiv \nu_0 + \sigma_\nu \sum_{k=1}^{n_{\nu}} \nu_k(x) \xi_k, \quad (4.2)
$$

where $\{\nu_0, \sigma_\nu^2\}$ are the mean and the variance of the random field $\nu(x, \xi)$. We will also refer to the coefficient of variation ($CoV$), the relative size of the standard deviation with respect to the mean,

$$
CoV \equiv \frac{\sigma_\nu}{\nu_0}. \quad (4.3)
$$
The random viscosity leads to the random Reynolds number

\[ \text{Re}(\xi) \equiv \frac{UL}{\nu(\xi)}, \quad (4.4) \]

where \( U \) is the characteristic velocity and \( L \) is the characteristic length. We denote the Reynolds number associated with the mean viscosity by \( \text{Re}_0 = \frac{UL}{\nu_0} \). In this study, we ensure that the viscosity (4.2) has positive values by controlling \( \text{CoV} \) and only consider small enough \( \text{Re}_0 \) so that the flow problem has a unique solution.

4.2.1 Stochastic Galerkin method

In the stochastic Galerkin method, a mixed variational formulation of (4.1) can be obtained by employing Galerkin orthogonality: Find \((\bar{u}, p) \in (V_E, Q_D) \otimes L^2(\Gamma)\) such that

\[ \left\langle \int_D \nu \nabla \bar{u} : \nabla \tilde{v} + (\bar{u} \cdot \nabla \bar{u})\tilde{v} - p(\nabla \cdot \bar{v}) \right\rangle_{\rho} = \left\langle \int_D \bar{f} \cdot \tilde{v} \right\rangle_{\rho}, \quad \forall \tilde{v} \in V_D \otimes L^2(\Gamma), \quad (4.5) \]

\[ \left\langle \int_D q(\nabla \cdot \bar{u}) \right\rangle_{\rho} = 0, \quad \forall q \in Q_D \otimes L^2(\Gamma). \quad (4.6) \]

The velocity solution and test spaces are \( V_E = \{ \bar{u} \in \mathcal{H}^1(D)^2| \bar{u} = \bar{g} \text{ on } \partial D_{\text{Dir}} \} \) and \( V_D = \{ \tilde{v} \in \mathcal{H}^1(D)^2| \tilde{v} = 0 \text{ on } \partial D_{\text{Dir}} \} \), where \( \mathcal{H}^1(D) \) refers to the Sobolev space of functions with derivatives in \( L^2(D) \), for the pressure solution, \( Q_D = L^2(D) \), and \( L^2(\Gamma) \) is a Hilbert space equipped with an inner product

\[ \langle u, v \rangle_{\rho} \equiv \int_{\Gamma} u(\xi)v(\xi)\rho(\xi)d\xi. \]

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The solution of the variational formulation (4.5)–(4.6) satisfies

\[ R(\vec{u}, p; \vec{v}, q) = 0, \quad \forall \vec{v} \in V_D \otimes L^2(\Gamma), \forall q \in Q_D \otimes L^2(\Gamma), \quad (4.7) \]

where \( R(\vec{u}, p; \vec{v}, q) \) is a nonlinear residual

\[ R(\vec{u}, p; \vec{v}, q) \equiv \begin{bmatrix} \left\langle \int_D \vec{f} \cdot \vec{v} - \nu \nabla \vec{u} : \nabla \vec{v} + (\vec{u} \cdot \nabla \vec{u}) \vec{v} - \int_D p (\nabla \cdot \vec{v}) \right\rangle_\rho \\ \left\langle - \int_D q (\nabla \cdot \vec{u}) \right\rangle_\rho \end{bmatrix}. \quad (4.8) \]

To compute the solution of the nonlinear equation (4.7), we employ linearization techniques based on either Picard iteration or Newton iteration [39]. Replacing \((\vec{u}, p)\) of (4.5)–(4.6) with \((\vec{u} + \delta\vec{u}, p + \delta p)\) and neglecting the quadratic term \(c(\delta\vec{u}; \delta\vec{u}, \vec{v})\), where \(c(\vec{z}; \vec{u}, \vec{v}) \equiv \int_D (\vec{z} \cdot \nabla \vec{u}) \cdot \vec{v}, \) gives

\[ \begin{bmatrix} \left\langle \int_D \nu \nabla \delta\vec{u} : \nabla \vec{v} + c(\delta\vec{u}; \vec{u}, \vec{v}) + c(\vec{u}; \delta\vec{u}, \vec{v}) - \int_D \delta p (\nabla \cdot \vec{v}) \right\rangle_\rho \\ \left\langle \int_D q (\nabla \cdot \delta\vec{u}) \right\rangle_\rho \end{bmatrix} = R(\vec{u}, p; \vec{v}, q). \quad (4.9) \]

In Newton iteration, the \((n+1)\)st iterate \((\vec{u}^{n+1}, p^{n+1})\) is computed by taking \(\vec{u} = \vec{u}^n, p = p^n\) in (4.9), solving (4.9) for \((\delta\vec{u}^n, \delta p^n)\), and updating

\[ \vec{u}^{n+1} \coloneqq \vec{u}^n + \delta\vec{u}^n, \quad p^{n+1} \coloneqq p^n + \delta p^n. \]

In Picard iteration, the term \(c(\delta\vec{u}; \vec{u}, \vec{v})\) is omitted from the linearized form (4.9).
4.2.2 Discrete stochastic Galerkin system

To obtain a discrete system, the velocity $\vec{u}(x,\xi)$ and the pressure $p(x,\xi)$ are approximated by a generalized polynomial chaos expansion [112]:

$$\vec{u}(x,\xi) \equiv \sum_{i=1}^{n_\xi} \vec{u}_i(x)\psi_i(\xi), \quad p(x,\xi) \equiv \sum_{i=1}^{n_\xi} p_i(x)\psi_i(\xi),$$  \hspace{1cm} (4.10)

where $\{\psi_i(\xi)\}_{i=1}^{n_\xi}$ is a set of $n_\nu$-variate orthogonal polynomials (i.e., $\langle \psi_i \psi_j \rangle_\rho = 0$ if $i \neq j$). This set of orthogonal polynomials gives rise to a finite-dimensional approximation space $S = \text{span}(\{\psi_i(\xi)\}_{i=1}^{n_\xi}) \subset L^2(\Gamma)$. For spatial discretization, a div-stable mixed finite element method [39] is considered, the Taylor-Hood element consisting of biquadratic velocities and bilinear pressure. Basis sets for the velocity space $V_h^E$ and the pressure space $Q_h^D$ are denoted by

$$\phi_i(x) = \begin{bmatrix} \phi_i(x) \\ 0 \end{bmatrix}, \quad \phi_i(x) = \begin{bmatrix} 0 \\ \phi_i(x) \end{bmatrix},$$

and $\{\varphi_i(x)\}_{i=1}^{n_p}$, respectively. Then the fully discrete version of (4.10) can be written as

$$\vec{u}(x,\xi) = \begin{bmatrix} \vec{u}_x(x,\xi) \\ \vec{u}_y(x,\xi) \end{bmatrix} = \begin{bmatrix} \sum_{i=1}^{n_\xi} \sum_{j=1}^{n_u} u_{ij}^x \phi_j(x)\psi_i(\xi) \\ \sum_{i=1}^{n_\xi} \sum_{j=1}^{n_u} u_{ij}^y \phi_j(x)\psi_i(\xi) \end{bmatrix}, \quad p(x,\xi) = \sum_{i=1}^{n_\xi} \sum_{j=1}^{n_p} p_{ij} \varphi_j(x)\psi_i(\xi).$$  \hspace{1cm} (4.11)

Let us introduce a vector notation for the coefficients, $\vec{u}_i^x \equiv [u_{i1}^x, \ldots, u_{in_u}^x]^T \in \mathbb{R}^{n_u}$, $\vec{u}_i^y \equiv [u_{i1}^y, \ldots, u_{in_u}^y]^T \in \mathbb{R}^{n_u}$, and $p_i \equiv [p_{i1}, \ldots, p_{in_p}]^T \in \mathbb{R}^{n_p}$ for $i = 1, \ldots, n_\xi$, which, for each gPC index $i$, groups the horizontal velocity coefficients together followed by the vertical velocity coefficients, and then by the pressure coefficients, giving a
vector
\[ \bar{u}_i = [\bar{u}_i^x, \bar{u}_i^y, p_i]^T. \quad (4.12) \]

Taking \( \nu(x, \xi) \) from (4.2) and replacing \( \bar{u}(x, \xi), p(x, \xi) \) in (4.9) with their discrete approximations (4.11) yields a system of linear equations of order \((2n_u + n_p)n_\xi \). The coefficient matrix has a Kronecker-product structure,

\[ J \equiv G_1 \otimes F_1 + \sum_{l=2}^{n_\xi} G_l \otimes F_l, \quad (4.13) \]

where \( G_l \) refers to the \( l \)th “stochastic matrix”

\[ [G_l]_{ij} = \langle \psi_l \psi_i \psi_j \rangle_{\rho}, \quad l = 1, \ldots, n_\xi \]

with \( \psi_1(\xi) = 1, \psi_i(\xi) = \xi_{i-1} \) for \( i = 2, \ldots, n_\nu + 1 \) and

\[ F_1 = \begin{bmatrix} F_1 & B^T \\ B & 0 \end{bmatrix}, \quad F_l = \begin{bmatrix} F_l & 0 \\ 0 & 0 \end{bmatrix}, \quad l = 2, \ldots, n_\xi \]

with \( F_l \equiv A_l + N_l + W_l \) for the Newton iteration and \( F_l \equiv A_l + N_l \) for the Picard iteration. We refer to the matrix of (4.13) derived from the Newton iteration as the \textit{Jacobian} matrix, and that derived from the Picard iteration as the \textit{Oseen} matrix, denoted by \( J_N \) and \( J_P \), respectively. Here, \( A_l \) is the \( l \)th symmetric matrix defined
\[ [A_l]_{ij} \equiv \int_D \nu_{l-1}(x)(\nabla \phi_i : \nabla \phi_j), \quad l = 1, \ldots, n_\nu + 1, \quad (4.14) \]

\[ N_l = N(\bar{u}_l(x)) \quad \text{and} \quad W_l = W(\bar{u}_l(x)) \quad \text{are, respectively, the} \]
\[ l\text{th vector-convection matrix and the} \]
\[ l\text{th Newton derivative matrix with} \]
\[ \bar{u}_l^n(x) \quad \text{from the} \]
\[ l\text{th term of (4.10)}, \]

\[ [N_l]_{ij} = [N(\bar{u}_l(x))]_{ij} \equiv \int_D (\bar{u}_l(x) \cdot \nabla \phi_j(x)) \cdot \phi_i(x), \quad l = 1, \ldots, n_\xi, \]
\[ [W_l]_{ij} = [W(\bar{u}_l(x))]_{ij} \equiv \int_D (\phi_j(x) \cdot \nabla \bar{u}_l(x)) \cdot \phi_i(x), \quad l = 1, \ldots, n_\xi, \]

and \( B \) is the divergence matrix,

\[ [B]_{ij} \equiv \int_D \phi_j(\nabla \cdot \phi_i). \quad (4.15) \]

If the number of gPC polynomial terms in (4.11) is larger than the number of terms

in (4.2) (i.e., \( n_\xi > n_\nu + 1 \)), we simply set \( \{A_l\}_{l=n_\nu + 2}^{n_\xi} \) as matrices containing only

zeros so that \( \mathcal{F}_l = N_l + W_l \) for \( l = n_\nu + 2, \ldots, n_\xi. \)

A discrete version of (4.8) can be derived in a similar way,

\[ \bar{r} := \bar{y} - \left( G_1 \otimes \mathcal{P}_1 + \sum_{l=2}^{n_\xi} G_l \otimes \mathcal{P}_l \right) \bar{u} \quad (4.16) \]

where \( \bar{u} := [\bar{u}_1^T \ldots \bar{u}_{n_\xi}^T]^T \in \mathbb{R}^{(2n_u + n_p)n_\xi} \) with \( \bar{u}_i \) as in (4.12), \( \bar{y} \) is the right-hand side.
determined from the forcing function and Dirichlet boundary data, and

\[
P_1 \equiv \begin{bmatrix} A_1 + N_1 & B^T \\ B & 0 \end{bmatrix}, \quad P_l \equiv \begin{bmatrix} A_l + N_l & 0 \\ 0 & 0 \end{bmatrix} \quad l = 2, \ldots, n_\xi.
\]

The system of linear equations arising at the \(n\)th nonlinear iteration is

\[
J^n \delta \bar{u}^n = -\bar{r}^n, \quad (4.17)
\]

where the matrix \(J^n\) from (4.13) and the residual \(\bar{r}^n\) from (4.16) each evaluated at the \(n\)th iterate \(\bar{u}^n\), and the update \(\delta \bar{u}^n\) is computed by solving (4.17). The order of the system \((2n_u + n_p)n_\xi\) grows fast as the number of random variables used to parameterize the random viscosity increases. Even for a moderate-dimensional stochastic Navier–Stokes problem, solving a sequence of linear systems of order \((2n_u + n_p)n_\xi\) can be computationally prohibitive. To address this issue, we present an efficient variant of Newton–Krylov methods in the following sections.

4.3 Low-rank Newton–Krylov method

In this section, we outline the formalism in which the solutions to (4.16) and (4.17) can be efficiently approximated by low-rank objects while not losing much accuracy and we show how solvers are adjusted within this formalism.

Before presenting these ideas, we describe the nonlinear iteration. We consider a hybrid strategy. An initial approximation for the nonlinear solution is computed
by solving the parameterized Stokes equations,

\[-\nu(x,\xi)\nabla^2\bar{u}(x,\xi) + \nabla p(x,\xi) = \bar{f}(x,\xi),\]

\[\nabla \cdot \bar{u}(x,\xi) = 0.\]

The discrete Stokes operator, which is obtained from the stochastic Galerkin discretization as shown in Section 4.2.2, is

\[
\left( G_1 \otimes S_1 + \sum_{l=2}^{n_\nu+1} G_l \otimes S_l \right) \bar{u}_{st} = b_{st}, \tag{4.18}
\]

where

\[
S_1 = \begin{bmatrix}
A_1 & B^T \\
B & 0
\end{bmatrix}, \quad S_l = \begin{bmatrix}
A_l & 0 \\
0 & 0
\end{bmatrix}, \quad l = 2, \ldots, n_\nu + 1,
\]

with \(\{A_l\}_{l=1}^{n_\nu+1}\) defined in (4.14) and \(B\) defined in (4.15). After this initial computation, updates to the solution are computed by first solving \(m_p\) Picard systems with coefficient matrix \(J_P\) and then using Newton’s method with coefficient matrix \(J_N\) to compute the solution.

Algorithm 4 Solution methods

1: compute an approximate solution of \(A_{st}\bar{u}_{st} = b_{st}\) in (4.18)
2: set an initial guess for the Navier–Stokes problem \(\bar{u}^0 := \bar{u}_{st}\)
3: for \(k = 0, \ldots, m_p - 1\) do
4: \hspace{1cm} solve \(J_P^k \delta \bar{u}^k = -\bar{r}^k\)
5: \hspace{1cm} update \(\bar{u}^{k+1} := \bar{u}^k + \delta \bar{u}^k\)
6: end for
7: while \(k < m_n\) and \(\|\bar{r}^k\|_2 > \epsilon_n\|\bar{r}^0\|_2\) do
8: \hspace{1cm} solve \(J_N^k \delta \bar{u}^k = -\bar{r}^k\)
9: \hspace{1cm} update \(\bar{u}^{k+1} := \bar{u}^k + \delta \bar{u}^k\)
10: end while

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4.3.1 Approximation in low rank

We now develop a low-rank variant of Algorithm 4. Let us begin by introducing some concepts to define the rank of computed quantities. Let \( X = [\bar{x}_1, \ldots, \bar{x}_{n_2}] \in \mathbb{R}^{n_1 \times n_2} \) and \( \bar{x} = [\bar{x}_1^T, \ldots, \bar{x}_{n_2}^T]^T \in \mathbb{R}^{n_1 n_2} \), where \( \bar{x}_i \in \mathbb{R}^{n_1} \) for \( i = 1, \ldots, n_2 \). That is, \( \bar{x} \) can be constructed by rearranging the elements of \( X \), and vice versa. Suppose \( X \) has rank \( \alpha_x \). Then two mathematically equivalent expressions for \( X \) and \( \bar{x} \) are given by

\[
X = YZ^T = \sum_{i=1}^{\alpha_x} \bar{y}_i \bar{z}_i^T \iff \bar{x} = \sum_{i=1}^{\alpha_x} \bar{z}_i \otimes \bar{y}_i,
\]

where \( Y \equiv [\bar{y}_1, \ldots, \bar{y}_{\alpha_x}] \in \mathbb{R}^{n_1 \times \alpha_x} \), \( Z \equiv [\bar{z}_1, \ldots, \bar{z}_{\alpha_x}] \in \mathbb{R}^{n_2 \times \alpha_x} \) with \( \bar{y}_i \in \mathbb{R}^{n_1} \), \( \bar{z}_i \in \mathbb{R}^{n_2} \) for \( i = 1, \ldots, \alpha_x \). The representation of \( X \) and its rank is standard matrix notation; we also use \( \alpha_x \) to refer to the rank of the corresponding vector \( \bar{x} \).

With this definition of rank, our goal is to inexpensively find a low-rank approximate solution \( \bar{u}^k \) satisfying \( \|\bar{r}^k\|_2 \leq \epsilon_{nl} \|\bar{r}^0\|_2 \) for small enough \( \epsilon_{nl} \). To achieve this goal, we approximate updates \( \{\delta \bar{u}^k\} \) in low-rank using a low-rank variant of GMRES method, which exploits the Kronecker product structure in the system matrix as in (4.13) and (4.18). In the following section, we present the solutions \( \bar{u} \) (and \( \delta \bar{u} \)) in the formats of (4.19) together with matrix and vector operations that are essential for developing the low-rank GMRES method.
4.3.2 Solution coefficients in Kronecker-product form

We seek separate low-rank approximations of the horizontal and vertical velocity solutions and the pressure solution. With the representation shown in (4.19), the solution coefficient vector \( \bar{u} \in \mathbb{R}^{(2n_u+n_p)n_\xi} \), which consists of the coefficients of the velocity solution and the pressure solution (4.11), has an equivalent representation \( U \in \mathbb{R}^{(2n_u+n_p)\times n_\xi} \). The matricized solution coefficients \( U = [U^x T, U^y T, P^T]^T \) where \( U^x = [\bar{u}_1^x, \ldots, \bar{u}_{n_\xi}^x] \), \( U^y = [\bar{u}_1^y, \ldots, \bar{u}_{n_\xi}^y] \) \( \in \mathbb{R}^{n_u \times n_\xi} \) and the pressure solution \( P = [\bar{p}_1, \ldots, \bar{p}_{n_\xi}] \) \( \in \mathbb{R}^{n_p \times n_\xi} \). The components admit the following representations:

\[
U^x = \sum_{i=1}^{\alpha_{\bar{u}^x}} \bar{v}_i^x (\bar{w}_i^x)^T = V^x (W^x)^T \quad \Leftrightarrow \quad \bar{u}^x = \sum_{i=1}^{\alpha_{\bar{u}^x}} \bar{w}_i^x \otimes \bar{v}_i^x, \quad (4.20)
\]

\[
U^y = \sum_{i=1}^{\alpha_{\bar{u}^y}} \bar{v}_i^y (\bar{w}_i^y)^T = V^y (W^y)^T \quad \Leftrightarrow \quad \bar{u}^y = \sum_{i=1}^{\alpha_{\bar{u}^y}} \bar{w}_i^y \otimes \bar{v}_i^y, \quad (4.21)
\]

\[
P = \sum_{i=1}^{\alpha_{\bar{p}}} \bar{v}_i^p (\bar{w}_i^p)^T = V^p (W^p)^T \quad \Leftrightarrow \quad \bar{p} = \sum_{i=1}^{\alpha_{\bar{p}}} \bar{w}_i^p \otimes \bar{v}_i^p, \quad (4.22)
\]

where \( V^x = [\bar{v}_1^x \ldots \bar{v}_{\alpha_{\bar{u}^x}}^x] \), \( W^x = [\bar{w}_1^x \ldots \bar{w}_{\alpha_{\bar{u}^x}}^x] \), \( \alpha_{\bar{u}^x} \) is the rank of \( \bar{u}^x \) and \( U^x \), and the same interpretation can be applied to \( \bar{u}^y \) and \( \bar{p} \).

4.3.2.1 Matrix operations

In this section, we introduce essential matrix operations used by the low-rank GMRES methods, using the representations shown in (4.20)–(4.22). First, consider the matrix-vector product with the Jacobian system matrix (4.13) and
vectors (4.20)-(4.22),

\[ J^n \bar{u}^n = \left( \sum_{l=1}^{n} G_l \otimes \mathcal{F}^n_l \right) \bar{u}^n, \tag{4.23} \]

where

\[
\mathcal{F}^n_l = \begin{bmatrix}
A^n_l & W^n_{xy,l} & B^x \\
W^n_{yx,l} & A^n_l & B^y \\
B^x & B^y & 0
\end{bmatrix}
= \begin{bmatrix}
\mathcal{F}^{xx,n}_l & \mathcal{F}^{xy,n}_l & B^x \\
\mathcal{F}^{yx,n}_l & \mathcal{F}^{yy,n}_l & B^y \\
B^x & B^y & 0
\end{bmatrix}
\]

with \( \mathcal{F}^{xx,n}_l, \mathcal{F}^{xy,n}_l, \mathcal{F}^{yx,n}_l, \mathcal{F}^{yy,n}_l \in \mathbb{R}^{n_u \times n_u} \) and \( B^x, B^y \in \mathbb{R}^{n_p \times n_u} \). The expression (4.23) has the equivalent matricized form \( \sum_{l=1}^{n} \mathcal{F}^n_l U^n G^T_l \) where the \( l \)th-term is evaluated as

\[
\mathcal{F}^n_l U^n G^T_l = \begin{bmatrix}
\mathcal{F}^{xx,n}_l V^{x,n}_n(G_l W^{x,n}_x)^T + \mathcal{F}^{xy,n}_l V^{y,n}_n(G_l W^{y,n}_y)^T + B^x V^{p,n}_n(G_l W^{p,n}_p)^T \\
\mathcal{F}^{yx,n}_l V^{x,n}_n(G_l W^{x,n}_x)^T + \mathcal{F}^{yy,n}_l V^{y,n}_n(G_l W^{y,n}_y)^T + B^y V^{p,n}_n(G_l W^{p,n}_p)^T \\
B^x V^{x,n}_n(G_l W^{x,n}_x)^T + B^y V^{y,n}_n(G_l W^{y,n}_y)^T
\end{bmatrix}.
\tag{4.24} \]
Equivalently, in the Kronecker-product structure, the matrix-vector product (4.24) updates each set of solution coefficients as follows:

\[
\sum_{l=1}^{n_\xi}(G_l \otimes F_l^{x,n})\vec{u}_x^n + (G_l \otimes F_l^{y,n})\vec{u}_y^n + (G_l \otimes B^{xT})\vec{p}_n, \quad (x\text{-velocity}), \quad (4.25)
\]

\[
\sum_{l=1}^{n_\xi}(G_l \otimes F_l^{yx,n})\vec{u}_x^n + (G_l \otimes F_l^{yy,n})\vec{u}_y^n + (G_l \otimes B^{yT})\vec{p}_n, \quad (y\text{-velocity}) \quad (4.26)
\]

\[
\sum_{l=1}^{n_\xi}(G_l \otimes B^x)\vec{u}_x^n + (G_l \otimes B^y)\vec{u}_y^n, \quad \text{(pressure)} \quad (4.27)
\]

where each matrix-vector product can be performed by exploiting the Kronecker-product structure, for example,

\[
\sum_{l=1}^{n_\xi}(G_l \otimes F_l^{x,n})\vec{u}_x^n = \sum_{l=1}^{n_\xi} G_l \otimes F_l^{x,n} \sum_{i=1}^{\alpha_{ux}} w_i^x \otimes v_i^x = \sum_{l=1}^{n_\xi} \sum_{i=1}^{\alpha_{ux}} G_l w_i^x \otimes F_l^{x,x_i} v_i^x.
\quad (4.28)
\]

The matrix-vector product shown in (4.25)–(4.27) requires \(O(2n_u + n_p + n_\xi)\) flops, whereas (4.23) requires \(O((2n_u + n_p)n_\xi)\) flops. Thus, as the problem size grows, the additive form of the latter count grows much less rapidly than the multiplicative form for (4.23).

The addition of two vectors \(\vec{u}^x\) and \(\vec{u}^y\) can also be efficiently performed in the Kronecker-product structure,

\[
\vec{u}^x + \vec{u}^y = \sum_{i=1}^{\alpha_{ux}} w_i^x \otimes v_i^x + \sum_{i=1}^{\alpha_{uy}} w_i^y \otimes v_i^y = \sum_{i=1}^{\alpha_{ux} + \alpha_{uy}} \hat{w}_i + \hat{v}_i, \quad (4.29)
\]

where \(\hat{w}_i = v_i^x, \hat{w}_i = w_i^x\) for \(i = 1, \ldots, \alpha_{ux}\), and \(\hat{v}_i = v_i^y, \hat{w}_i = w_i^y\) for \(i = \alpha_{ux} + \)
1, \ldots, \alpha_{\bar{a}x} + \alpha_{\bar{a}y}.

Inner products can be performed with similar efficiencies. Consider two vectors \( \bar{x}_1 \) and \( \bar{x}_2 \), whose matricized representations are

\[
X_1 = \begin{bmatrix}
Y_{11}Z_{11}^T \\
Y_{12}Z_{12}^T \\
Y_{13}Z_{13}^T
\end{bmatrix}, \quad X_2 = \begin{bmatrix}
Y_{21}Z_{21}^T \\
Y_{22}Z_{22}^T \\
Y_{23}Z_{23}^T
\end{bmatrix}.
\] (4.30)

Then the Euclidean inner product between \( \bar{x}_1 \) and \( \bar{x}_2 \) can be evaluated as

\[
\bar{x}_1^T \bar{x}_2 = \text{trace}((Y_{11}Z_{11}^T)^TY_{21}Z_{21}^T) + \text{trace}((Y_{12}Z_{12}^T)^TY_{22}Z_{22}^T) + \text{trace}((Y_{13}Z_{13}^T)^TY_{23}Z_{23}^T),
\]

where \( \text{trace}(X) \) is defined as a sum of the diagonal entries of the matrix \( X \).

Although the matrix-vector product and the sum, as described in (4.28) and (4.29), can be performed efficiently, the results of (4.28) and (4.29) are represented by \( n_{\xi}\alpha_{a^x} \) terms and \( \alpha_{\bar{a}x} + \alpha_{\bar{a}y} \) terms, respectively, which typically causes the ranks of the computed quantities to be higher than the inputs for the computations and potentially undermines the efficiency of the solution method. To resolve this issue, a truncation operator will be used to modify the result of matrix-vector products and sums and force the ranks of quantities used to be small.

4.3.2.2 Truncation of \( U^{x,n} \), \( U^{y,n} \) and \( P^n \)

We now explain the details of the truncation. Consider the velocity and the pressure represented in a matrix form as in (4.20)–(4.22). The best \( \alpha \)-rank ap-
proximation of a matrix can be found by using the singular value decomposition (SVD) \[61, 70\]. Here, we define a truncation operator for a given matrix \(U = VW^T\) whose rank is \(\alpha_U\),

\[
\mathcal{T}_{\text{trunc}} : U \to \tilde{U},
\]

where the rank of \(U\) is larger than the rank of \(\tilde{U}\) (i.e., \(\alpha_U \gg \alpha_{\tilde{U}}\)). The truncation operator \(\mathcal{T}_{\text{trunc}}\) compresses \(U\) to \(\tilde{U}\) such that \(\|\tilde{U} - U\|_F \leq \epsilon_{\text{trunc}}\|U\|_F\) where \(\|\cdot\|_F\) is the Frobenius norm. To achieve this goal, the singular value decomposition of \(U\) can be computed (i.e., \(U = \hat{V}D\hat{W}^T\) where \(D = \text{diag}(d_1, \ldots, d_n)\) is the diagonal matrix of singular values). Letting \(\{\hat{v}_i\}\) and \(\{\hat{w}_i\}\) denote the singular vectors, the approximation is \(\tilde{U} = \sum_{i=1}^{\alpha_{\tilde{U}}} \hat{v}_i\hat{w}_i^T\) with \(\hat{v}_i = d_i\hat{v}_i\) and the truncation rank \(\alpha_{\tilde{U}}\) is determined by the condition

\[
\sqrt{d_{\alpha_{\tilde{U}}+1}^2 + \cdots + d_n^2} \leq \epsilon_{\text{trunc}}\sqrt{d_1^2 + \cdots + d_n^2}.
\] (4.31)

### 4.3.3 Low-rank GMRES method

We describe the low-rank GMRES method (lrGMRES) with a generic linear system \(Ax = b\). The method follows the standard Arnoldi iteration used by GMRES \[92\]: construct a set of basis vectors \(\{v_i\}_{i=1}^{m_{gm}}\) by applying the linear operator \(A\) to basis vectors, i.e., \(w_j = Av_j\) for \(j = 1, \ldots, m_{gm}\), and orthogonalizing the resulting vector \(w_j\) with respect to previously generated basis vectors \(\{v_i\}_{i=1}^{j-1}\). In the low-rank GMRES method, iterates, basis vectors \(\{v_i\}\) and intermediate quantities \(\{w_i\}\) are
represented in terms of the factors of their matricized representations (so that $X$ in (4.19) would be represented using $Y$ and $Z$ without explicit construction of $X$), and matrix operations such as matrix-vector products are performed as described in Section 4.3.2.1. As pointed out in Section 4.3.2.1, these matrix operations typically tend to increase the rank of the resulting quantity, and this is resolved by interleaving the truncation operator $\mathcal{T}$ with the matrix operations. The low-rank GMRES method computes a new iterate by solving

$$
\min_{\beta \in \mathbb{R}^{m_{gm}}} \|b - A(x_0 + V_{m_{gm}}\bar{\beta})\|_2,
$$

(4.32)

and constructing a new iterate $x_1 = x_0 + V_{m_{gm}}\bar{\beta}$ where $x_0$ is an initial guess. Due to truncation, the basis vectors $\{v_i\}$ are not orthogonal and $\text{span}(V_{m_{gm}})$, where $V_{m_{gm}} = [v_1 \ldots v_{m_{gm}}]$, is not a Krylov subspace, so that (4.32) must be solved explicitly rather than exploiting Hessenberg structure as in standard GMRES. Algorithm 5 summarizes the lrGMRES. We will use this method to solve the linear system of (4.17).

4.3.4 Preconditioning

We also use preconditioning to speed convergence of the low-rank GMRES method. For this, we consider a right-preconditioned system

$$
J^n(M^n)^{-1}\tilde{u}^n = \tilde{r}^n,
$$
Algorithm 5 Restarted low-rank GMRES method in tensor format

1: set the initial solution $\bar{u}^0_{gm}$
2: for $k = 0, 1, \ldots$ do
3: $r^k_{gm} := f - A\bar{u}^k_{gm}$
4: if $\|r^k_{gm}\|_2/\|f\|_2 < \epsilon_{gmres}$ or $\|r^k_{gm}\|_2 \geq \|r^{k-1}_{gm}\|_2$ then
5: return $\bar{u}^k_{gm}$
6: end if
7: $\bar{v}_1 := \mathcal{T}_{\text{trunc}}(r^k_{gm})$
8: $v_1 := \bar{v}_1/\|\bar{v}_1\|_2$
9: for $j = 1, \ldots, m_{gm}$ do
10: $w_j := Av_j$
11: solve $(V_j^T V_j)\bar{v}_j = V_j^T w_j$ where $V_j = [v_1, \ldots, v_j]$
12: $\bar{v}_{j+1} := \mathcal{T}_{\text{trunc}}(w_j - \sum_{i=1}^j \alpha_i v_i)$
13: $v_{j+1} := \bar{v}_{j+1}/\|\bar{v}_{j+1}\|_2$
14: end for
15: solve $(W_{mgm}^T A V_{mgm})\bar{\beta} = W_{mgm}^T r^k_{gm}$ where $W_j = [w_1, \ldots, w_j]$
16: $\bar{u}^{k+1}_{gm} := \mathcal{T}_{\text{trunc}}(\bar{u}^k_{gm} + V_{mgm}\bar{\beta})$
17: end for

where $M^n$ is the preconditioner and $M^n\bar{u}^n = \bar{u}^n$ such that $J^n\bar{u}^n = \bar{r}^n$. We consider an approximate mean-based preconditioner [81], which is derived from the matrix $G_1 \otimes \mathcal{F}_1$ associated with the mean $\nu_0$ of the random viscosity (4.2),

$$M^n = G_1 \otimes \begin{bmatrix} M^n_A & B^T \\ 0 & -M^n_s \end{bmatrix}, \quad (4.33)$$

where

$$M^n_A = \begin{bmatrix} A_1^{xx} + N_1^n & 0 \\ 0 & A_1^{yy} + N_1^n \end{bmatrix}, \quad \text{(Picard iteration)},$$

$$M^n_A = \begin{bmatrix} A_1^{xx} + N_1^n + W_1^{xx,n} & 0 \\ 0 & A_1^{yy} + N_1^n + W_1^{yy,n} \end{bmatrix}, \quad \text{(Newton iteration)}.$$
For approximating the action of the inverse, \((M_s^n)^{-1}\), we choose the boundary-adjusted least-squares commutator (LSC) preconditioning scheme \([39]\),

\[
M_s^n = BF_1^{-1}B^T \approx (BH^{-1}B^T)(BM_s^{-1}F_1H^{-1}B^T)^{-1}(BM_s^{-1}B^T),
\]

where \(M_s\) is the diagonal of the velocity mass matrix and \(H = D^{-1/2}M_sD^{-1/2}\), where \(D\) is a diagonal scaling matrix deemphasizing contributions near the boundary. During the low-rank GMRES iteration, the action of the inverse of the preconditioner \((4.33)\) can be applied to a vector in a manner analogous to \((4.25)-(4.27)\).

4.4 Inexact nonlinear iteration

As outlined in Algorithm 4, we use the hybrid approach, employing a few steps of Picard iteration followed by Newton iteration, and the linear systems are solved using lrGMRES (Algorithm 5). We extend the hybrid approach to an inexact variant based on an inexact Newton algorithm, in which the accuracy of the approximate linear system solution is tied to the accuracy of the nonlinear iterate (see e.g., \([57]\) and references therein). That is, when the nonlinear iterate is far from the solution, the linear systems may not have to be solved accurately. Thus, a sequence of iterates

\[
\bar{u}^{n+1} := \bar{u}^n + \delta \bar{u}^n
\]

is computed where \(\delta \bar{u}^n\) satisfies

\[
\|J_N^n \delta \bar{u}^n + \bar{r}^n\|_2 \leq \epsilon_{gmres}^n \|\bar{r}^n\|_2, \quad (J_P \text{ for Picard iteration}),
\]
where the lrGMRES stopping tolerance ($\epsilon_{gmres}^n$ of Algorithm 5) is given by

$$
\epsilon_{gmres}^n := \rho_{gmres} \| \vec{r}^n \|_2,
$$

(4.34)

where $0 < \rho_{gmres} \leq 1$. With this strategy, the Jacobian system is solved with increased accuracy as the error becomes smaller, leading to savings in the average cost per step and, as we will show, with no degradation in the asymptotic convergence rate of the nonlinear iteration.

In addition, in Algorithms 4 and 5, the truncation operator $T_{trunc}$ is used for the low-rank approximation of the nonlinear iterate (i.e., truncating $\bar{u}^x$, $\bar{u}^y$, and $\bar{p}$ at lines 5 and 9 in Algorithm 4) and updates (i.e., truncating $\delta \bar{u}^x$, $\delta \bar{u}^y$, and $\delta \bar{p}$ in Algorithm 5). As the lrGMRES stopping tolerance is adaptively determined by the criterion (4.34), we also choose the value of the truncation tolerances $\epsilon_{trunc,sol}$ and $\epsilon_{trunc,corr}^n$, adaptively. For truncating the nonlinear iterate, the truncation tolerance for the iterate $\{\epsilon_{trunc,sol}^n\}$ is chosen based on the nonlinear iteration stopping tolerance,

$$
\epsilon_{trunc,sol} := \rho_{nl} \epsilon_{nl},
$$

where $0 < \rho_{nl} \leq 1$. For truncating the updates (or corrections), the truncation tolerance for the correction $\{\epsilon_{trunc,corr}^n\}$ is adaptively chosen based on the stopping
tolerance of the linear solver,

\[ \varepsilon_{n,\text{trunc,corr}}^n := \rho_{\text{trunc,P}} \varepsilon_{n,\text{gmres}}^n, \quad \text{(for the } n\text{th Picard step)}, \]

\[ \varepsilon_{n,\text{trunc,corr}}^n := \rho_{\text{trunc,N}} \varepsilon_{n,\text{gmres}}^n, \quad \text{(for the } n\text{th Newton step)}, \]

where \( 0 < \rho_{\text{trunc,P}}, \rho_{\text{trunc,N}} \leq 1 \). Thus, for computing \( n\text{th update } \delta \bar{u}^n \), we set \( \varepsilon_{\text{trunc}} = \varepsilon_{n,\text{trunc,corr}}^n \) in Algorithm 5.

**Algorithm 6** Inexact nonlinear iteration with adaptive tolerances

1: set \( \varepsilon_{\text{trunc,sol}} := \rho_{\text{nl}} \varepsilon_{\text{nl}} \)
2: compute an approximate solution of \( A_{st} \bar{u}_{st} = b_{st} \) using Algorithm 5
3: set an initial guess for the Navier–Stokes problem \( \bar{u}^0 := \bar{u}_{st} \)
4: for \( k = 0, \ldots, m_p - 1 \) do
5: set \( \varepsilon_{\text{gmres}}^k = \rho_{\text{gmres}} \| \bar{r}^k \|_2 \), and \( \varepsilon_{\text{trunc,corr}}^k = \rho_{\text{trunc,P}} \| \bar{r}^k \|_2 \)
6: solve \( J_{P}^k \delta \bar{u}^k = -\bar{r}^k \) using Algorithm 5
7: update \( \bar{u}^{k+1} := T_{\varepsilon_{\text{trunc,sol}}} (\bar{u}^k + \delta \bar{u}^k) \)
8: end for
9: while \( \| \bar{r}^k \|_2 > \varepsilon_{\text{nl}} \| \bar{r}^0 \|_2 \) do
10: set \( \varepsilon_{\text{gmres}}^k = \rho_{\text{gmres}} \| \bar{r}^k \|_2 \), and \( \varepsilon_{\text{trunc,corr}}^k = \rho_{\text{trunc,N}} \| \bar{r}^k \|_2 \)
11: solve \( J_{N}^k \delta \bar{u}^k = -\bar{r}^k \) using Algorithm 5
12: update \( \bar{u}^{k+1} := T_{\varepsilon_{\text{trunc,sol}}} (\bar{u}^k + \delta \bar{u}^k) \)
13: end while

4.5 Numerical results

In this section, we present the results of numerical experiments on a model problem, flow around a square obstacle in a channel, for which the details are depicted in Figure 4.1. The domain has length 12 and height 2, and it contains a square obstacle centered at (2,0) with sides of length .25.

For the numerical experiments, we define the random viscosity (4.2) using the
Karhunen-Loève (KL) expansion \[67\],

\[
\nu(x, \xi) = \nu_0 + \sigma_\nu \sum_{i=1}^{n_\nu} \sqrt{\lambda_i} \nu_i(x) \xi_i, \tag{4.35}
\]

where \(\nu_0\) and \(\sigma^2_\nu\) are the mean and the variance of the viscosity of \(\nu(x, \xi)\), and \(\{ (\lambda_i, \nu_i(x)) \}_{i=1}^{n_\nu} \) are eigenpairs of the eigenvalue problem associated with the covariance kernel \(C(x, y)\) of the random field. We consider two types of covariance kernel: absolute difference exponential (AE) and squared difference exponential (SE), which are defined via

\[
C^{AE}(x, y) = \exp \left( -\sum_{i=1}^{2} \frac{|x_i - y_i|}{l_i} \right), \quad C^{SE}(x, y) = \exp \left( -\sum_{i=1}^{2} \frac{(x_i - y_i)^2}{l_i^2} \right), \tag{4.36}
\]

where \(x = (x_1, x_2)\) and \(y = (y_1, y_2)\) are points in the spatial domain, and \(l_1, l_2\) are correlation lengths. We assume that the random variables \(\{\xi_i\}_{i=1}^{n_\nu}\) are independent and identically distributed and that \(\xi_i\) (for \(i = 1, \ldots, n_\nu\)) follows a uniform distribution over \([-1, 1]\). For the mean of the viscosity, we consider several choices, \(\nu_0 = \{ \frac{1}{50}, \frac{1}{100}, \frac{1}{150} \}\), which corresponds to \(\text{Re}_0 = \{100, 200, 300\}\). In all experiments, we use a finite-term KL-expansion with \(n_\nu = 5\). For constructing the
finite-dimensional approximation space \( S = \text{span}\{\psi_i(\xi)\}_{i=1}^{n_\xi} \) in the parameter domain, we use orthogonal polynomials \( \{\psi_i(\xi)\}_{i=1}^{n_\xi} \) of total degree 3, which results in \( n_\xi = 56 \). The orthogonal polynomials associated with uniform random variables are Legendre polynomials, \( \psi_i(\xi) = \prod_{j=1}^{n_\nu} \ell_{d_j(i)}(\xi_j) \) where \( d(i) = (d_1(i), \ldots, d_{n_\nu}(i)) \) is a multi-index consisting of non-negative integers and \( \ell_{d_j(i)} \) is the \( d_j(i) \)th order Legendre polynomial of \( \xi_j \). For the spatial discretization, Taylor–Hood elements are used on a stretched grid, which results in \( \{6320, 6320, 1640\} \) degrees of freedom in \( \{\vec{u}^x, \vec{u}^y, p\} \), respectively (i.e., \( n_u = 6320 \) and \( n_p = 1640 \).) The implementation is based on the Incompressible Flow and Iterative Solver Software (IFISS) package [38, 98].

4.5.1 Low-rank inexact nonlinear iteration

In this section, we compare the results obtained from the low-rank inexact nonlinear iteration with those obtained from other methods, the exact and the inexact nonlinear iteration with full rank solutions, and the Monte Carlo method. Default parameter settings are listed in Table 4.1, where the truncation tolerances only apply to the low-rank method. Unless otherwise specified, the linear system is solved using a restarted version of low-rank GMRES, lrGMRES(20).

We first examine the convergence behavior of the inexact nonlinear iteration for a model problem characterized by \( \text{Re}_0 = 100, \text{CoV} = 1\% \), and SE covariance kernel in (4.36) with \( l_1 = l_2 = 32 \). We compute a full-rank solution using the exact nonlinear iteration (\( \epsilon_{\text{gmres}} = 10^{-12} \) and no truncation) until the nonlinear iterate reaches
Table 4.1: Tolerances and adaptive parameters.

<table>
<thead>
<tr>
<th>Description</th>
<th>Formula</th>
</tr>
</thead>
<tbody>
<tr>
<td>Nonlinear iteration stopping tolerance</td>
<td>$\epsilon_{\text{nl}} = 10^{-5}$</td>
</tr>
<tr>
<td>GMRES tolerance (Stokes)</td>
<td>$\epsilon_{\text{gmres}} = 10^{-4}$</td>
</tr>
<tr>
<td>GMRES tolerances (Picard and Newton)</td>
<td>$\epsilon_{\text{gmres}}^{n} = \rho_{\text{gmres}} | \bar{r}_n |<em>2$ ($\rho</em>{\text{gmres}} = 10^{-5}$)</td>
</tr>
<tr>
<td>Truncation tolerance for solutions</td>
<td>$\epsilon_{\text{trunc,sol}} = \rho_{\text{nl}} \epsilon_{\text{nl}}$ ($\rho_{\text{nl}} = 10^{-1}$)</td>
</tr>
<tr>
<td>Truncation tolerance for corrections</td>
<td>$\epsilon_{\text{trunc,corr}}^{n} = \rho_{\text{trunc}} \epsilon_{\text{gmres}}$ ($\rho_{\text{trunc}} = 10^{-1}$)</td>
</tr>
</tbody>
</table>

the nonlinear stopping tolerance, $\epsilon_{\text{nl}} = 10^{-8}$. Then we compute another full-rank solution using the inexact nonlinear iteration (i.e., adaptive choice of $\epsilon_{\text{gmres}}^{n}$ as shown in Table 4.1 and no truncation). Lastly, we compute a low-rank approximate solution using the low-rank inexact nonlinear iteration (i.e., adaptive choices of $\epsilon_{\text{gmres}}^{n}$ and $\epsilon_{\text{trunc,corr}}^{n}$ as shown in Table 4.1 and for varying $\epsilon_{\text{trunc,sol}} = \{10^{-5}, 10^{-6}, 10^{-7}, 10^{-8}\}$).

Figure 4.2 shows the convergence behavior of the three methods. In Figure 4.2(a), the hybrid approach is used, in which the first step corresponds to the Stokes problem (line 2 of Algorithm 6), the 2nd–5th steps correspond to the Picard iteration (line 4–8 of Algorithm 6, and $m_p = 4$), and the 6th–7th steps correspond to the Newton iteration (line 9–13 of Algorithm 6). Figure 4.2(a) confirms that the inexact nonlinear iteration is as effective as the exact nonlinear iteration. The low-rank inexact nonlinear iteration behaves similarly up to the 6th nonlinear step but when the truncation tolerances are large $\epsilon_{\text{trunc,sol}} = \{10^{-5}, 10^{-6}\}$, it fails to produce a nonlinear solution satisfying $\epsilon_{\text{nl}} = 10^{-8}$. Similar results can be seen in Figure 4.2(b), where only the Picard iteration is used. As expected, in that case, the relative residual decreases linearly for all solution methods, but the low-rank inexact nonlinear iteration with the mild truncation tolerances also fails to reach the nonlinear
iteration stopping tolerance.

Figure 4.2: Convergence of both exact and inexact nonlinear iterations (full-rank) and the low-rank inexact nonlinear iteration.

Figure 4.3 shows means and variances of the components of the full-rank solution, given by

\[
\begin{align*}
\mu_{u^x} &= \mathbb{E}[\bar{u}^x], & \mu_{u^y} &= \mathbb{E}[\bar{u}^y], & \mu_p &= \mathbb{E}[p], \\
\sigma_{u^x}^2 &= \mathbb{E}[(\bar{u}^x - \mu_{u^x})^2], & \sigma_{u^y}^2 &= \mathbb{E}[(\bar{u}^y - \mu_{u^y})^2], & \sigma_p^2 &= \mathbb{E}[(p - \mu_p)^2].
\end{align*}
\] (4.37) (4.38)

These quantities are easily computed by exploiting the orthogonality of basis functions in the gPC expansion. Figure 4.4 shows the differences in the means and variances of the solutions computed using the full-rank and the low-rank inexact nonlinear iteration. Let us denote the full-rank and low-rank horizontal velocity solutions by \(u^{x,\text{full}}\) and \(u^{x,\text{lr}}\), with analogous notation for the vertical velocity and
Figure 4.3: Mean and variances of full-rank velocity solutions \( \vec{u}_x(x, \xi), \vec{u}_y(x, \xi) \), and pressure solution \( p(x, \xi) \) for \( Re_0 = 100, CoV = 1 \), and \( l_1 = l_2 = 32 \).

the pressure. Thus, the differences in the means and the variances are

\[
\eta_\mu^x = \mu_{u_x, \text{full}} - \mu_{u_x, \text{lr}}, \quad \eta_\mu^y = \mu_{u_y, \text{full}} - \mu_{u_y, \text{lr}}, \quad \eta_\mu^p = \mu_{p, \text{full}} - \mu_{p, \text{lr}}, \\
\eta_\sigma^x = \sigma_{u_x, \text{full}}^2 - \sigma_{u_x, \text{lr}}^2, \quad \eta_\sigma^y = \sigma_{u_y, \text{full}}^2 - \sigma_{u_y, \text{lr}}^2, \quad \eta_\sigma^p = \sigma_{p, \text{full}}^2 - \sigma_{p, \text{lr}}^2.
\]

Figure 4.4 shows these differences, normalized by graph norms \( \| \nabla \vec{u}_{\text{full}} \| + \| \mu_{\text{full}} \| \) for the means and \( \| \nabla \sigma_{\text{full}}^2 \| + \| \sigma_{\text{full}}^2 \| \) for the variances, where \( \| \nabla \vec{u} \| = (\int_D \nabla \vec{u} : \nabla \vec{u} \, dx)^{\frac{1}{2}} \) and \( \| p \| = (\int_D p^2 \, dx)^{\frac{1}{2}} \). Figure 4.4 shows that the normalized differences in the mean and the variance are of order \( 10^{-9} \sim 10^{-10} \) and \( 10^{-10} \sim 10^{-12} \), respectively, i.e., the errors in low-rank solutions are considerably smaller than the magnitude of the truncation tolerances \( \epsilon_{\text{trunc,sol}}, \epsilon_{\text{trunc,corr}} \) (see Table 4.1).

4.5.2 Characteristics of the Galerkin solution

In this section, we examine various properties of the Galerkin solutions, with emphasis on comparison of the low-rank and full-rank versions of these solutions and
development of an enhanced understanding of the relation between the Galerkin solution and the polynomial chaos basis. We use the same experimental setting studied above (SE covariance kernel, $l_1 = l_2 = 32$, $Re_0 = 100$ and $CoV = 1\%$).

We begin by comparing the Galerkin solution with one obtained using Monte Carlo methods. In particular, we estimate a probability density function (pdf) of the velocity solutions $(\vec{u}_x(x, \xi), \vec{u}_y(x, \xi))$ and the pressure solution ($p(x, \xi)$) at a specific point on the spatial domain $D$. In the Monte Carlo method, we solve $n_{MC} = 25000$ deterministic systems, $\mathcal{R}(\vec{u}, p, \vec{v}, q; \xi^{(k)}) = 0$ associated with $n_{MC}$ realizations $\{\xi^{(k)}\}_{k=1}^{n_{MC}}$ in the parameter space. Using the MATLAB function ksdensity, the pdfs of $(\vec{u}_x(x, \xi), \vec{u}_y(x, \xi), p(x, \xi))$ are estimated at the spatial point with coordinates $(3.6436, 0)$, where the variance of $\vec{u}_x(x, \xi)$ is large (see Figure 4.3). The results are shown in Figure 4.5. They indicate that the pdf of the Galerkin solution is virtually identical to that of the Monte Carlo solution, and there is essentially no difference between the low-rank and full-rank results.

Next, we explore some characteristics of the Galerkin solution, focusing on the
horizontal velocity solution; the observations made here also hold for the other components of the solution. Given the coefficients of the velocity solution in matricized form, \( U^x \), the discrete velocity solution is then given by

\[
\vec{u}^x(x, \xi) = \Phi^T(x) U^x \Psi(\xi),
\]

where \( \Phi(x) = [\phi_1(x), \ldots, \phi_{n_u}(x)]^T \) and \( \Psi(\xi) = [\psi_1(\xi), \ldots, \psi_{n_\xi}(\xi)]^T \). Consider in particular the component of this expression corresponding to the jth column of \( U^x \):

\[
\left( \sum_{i=1}^{n_u} \tilde{u}_x^i \phi_i(x) \right) \psi_j(\xi)
\]
so that this (jth) column $\bar{u}_j^x = [U^x]_j$ corresponds to the coefficient of the jth polynomial basis function $\psi_j$. Figure 4.6 plots the values of the coefficients $\|\bar{u}_j^x\|_2$. (This data is computed with $Re_0 = 100$, $CoV = 1\%$, and SE covariance kernel with $l_1 = l_2 = 32$). Note that the gPC indices $\{j\}$ are in one-to-one correspondence with multi-indices $d(j) = (d_1(j), \ldots, d_{n_u}(j))$, where the element of the multi-index indicates the degree of univariate Legendre polynomial. The multi-indices $\{d(i)\}_{i=1}^{n_x}$ are ordered in the lexicographical order, for example, the first eight multi-indices are as $d(1) = (0, 0, 0, 0, 0)$, $d(2) = (1, 0, 0, 0, 0)$, $d(3) = (0, 1, 0, 0, 0)$, $d(4) = (0, 0, 1, 0, 0)$, $d(5) = (0, 0, 0, 1, 0)$, $d(6) = (0, 1, 0, 0, 0)$, $d(7) = (2, 0, 0, 0, 0)$, and $d(8) = (1, 1, 0, 0, 0)$. In Figure 4.6, the blue square is associated with the zeroth-order gPC component ($d(1)$), the red circles are associated with the first-order gPC components ($\{d(i)\}_{1 \leq i \leq 6}$), and so on. Let us focus on three gPC components associated only with $\xi_1$, $\{\psi_2(\xi) = \ell_1(\xi_1), \psi_7(\xi) = \ell_2(\xi_1), \psi_{22}(\xi) = \ell_3(\xi_1)\}$, where, for $i = 2, 7, 22$, the multi-indices are $d(2) = (1, 0, 0, 0, 0)$, $d(7) = (2, 0, 0, 0, 0)$, and $d(22) = (3, 0, 0, 0, 0)$. The figure shows that the coefficients of gPC components $\{\psi_2(\xi), \psi_7(\xi), \psi_{22}(\xi)\}$ decay more slowly than those of gPC components associated with other random variables $\{\xi_i\}_{i=1}^{n_x}$.

We continue the examination of this data in Figure 4.7(a), which shows two-dimensional mesh plots of the 2nd through 7th columns of $U^x$. These images show that these coefficients are either symmetric with respect to the horizontal axis, or reflectionally symmetric (equal in magnitude but of opposite sign), and (as also revealed in Figure 4.6), they tend to have smaller values as the index $j$ is increased.

We now look more closely at features of the factors of the low-rank approximate solution and compare these with those of the (unfactored) full-rank solu-
Figure 4.7: Plots of coefficients of gPC components 2–7 of $\vec{u}^x(x, \xi)$ and coefficients $v_i$ of $\theta_i(\xi)$ for $i = 2, \ldots, 7$ for $Re_0 = 100$, $CoV = 1$, and $l_1 = l_2 = 32$.

In the low-rank format, the solution is represented using factors $\vec{u}^x(x, \xi) = (\Phi^T(x)V^x)(\Psi^T(\xi)W^x)^T$. Let us introduce a concise notation of

$$\vec{u}^x(x, \xi) = Z^x_{\alpha_{\tilde{u}^x}}(x)\Theta^x_{\alpha_{\tilde{u}^x}}(\xi) = \sum_{i=1}^{\alpha_{\tilde{u}^x}} \zeta_i^x(x)\theta_i^x(\xi)$$

where $Z^x_{\alpha_{\tilde{u}^x}}(x) = [\zeta_1^x(x), \ldots, \zeta_{\alpha_{\tilde{u}^x}}^x(x)]$ and $\Theta^x_{\alpha_{\tilde{u}^x}}(\xi) = [\theta_1^x(\xi), \ldots, \theta_{\alpha_{\tilde{u}^x}}^x(\xi)]$ with $\zeta_i^x(x) = [\Phi^T(x)V^x]_i$ and $\theta_i^x(\xi) = [(\Psi^T(\xi)W^x)]_i$ for $i = 1, \ldots, \alpha_{\tilde{u}^x}$. Figure 4.7(b) shows the coefficients of the $i$th random variable $\theta_i(\xi)$. As opposed to the gPC coefficients of the
full-rank solution, the norms of the coefficients of \( \{ \theta_i(\xi) \} \) decrease monotonically as the index \( i \) increases. This is a consequence of the fact that the ordering for \( \{ \theta_i(\xi) \} \) comes from the singular values of \( U^x \). Figure 4.7(b) shows the 2nd-7th columns of \( V^x \). Figures 4.7(a) and 4.7(b) show that the coefficients \( \{ v_i \} \) of \( \{ \theta_i(\xi) \} \) are comparable to the coefficients \( \{ u_i^x \} \) of the gPC components. Each pair of components in the following parentheses is similar to each other: \((u_2, v_2), (u_3, v_3), (u_7, -v_4), (u_4, -v_7), (u_5, v_5), \) and \((u_6, -v_6)\).

While the columns of \( V^x \) show the resemblance to the subset of the columns of \( U^x \), \( W^x \) tends to act as a permutation matrix. Figure 4.8 shows a “heat map” of \((W^x)^T\), where values of the elements in \( W^x \) are represented as colors and the map shows that a very few elements of \( W_i^x \) are dominant and a sum of those elements is close to 1. Recall that \( \theta_i^x(\xi) = (W_i^x)^T \Psi(\xi) \). Many dominant elements are located in the diagonal of \( W^x \), which results in \( \theta_i^x(\xi) \approx \pm \psi_i(\xi) \) (e.g., \( i = 1, 2, 3, 5, \ldots \)). In the case of \( W_4^x \), the most dominant element is the 7th element and has a value close to -1, which results in \( \theta_4^x(\xi) \approx -\psi_7(\xi) \). As observed in Figure 4.6, \( \psi_7(\xi) \) has a larger contribution than other gPC components and, in the new solution representation, \( \theta_4^x(\xi) \), which consists mainly of \( \psi_7(\xi) \), appears earlier in the representation.
4.5.3 Computational costs

In this section, we assess the costs of the low-rank inexact nonlinear iteration under various experimental settings: two types of covariance kernels (4.36), varying CoV (4.3), and varying Re₀. In addition, for various values of these quantities, we investigate the decay of the eigenvalues \{\lambda_i\} used to define the random viscosity (4.35) and their influence on the rank of solutions. All numerical experiments are performed on an Intel 3.1 GHz i7 CPU, 16 GB RAM using MATLAB R2016b and costs are measured in terms of CPU wall time (in seconds). For larger CoV and Re₀, we found the solver to be more effective using the slightly smaller truncation tolerance \rho_{\text{trunc}} = 10^{-1.5} and used this choice for all experiments described below. (Other adaptive tolerances are those shown as in Table 4.1.) This change had little impact on results for small CoV and Re₀.

Figure 4.9 shows the 50 largest eigenvalues \{\lambda_i\} of the eigenvalue problems associated with the SE covariance kernel and the AE covariance kernel (4.36) with \(l_1 = l_2 = 8, \text{CoV} = 1\%\), and Re₀ = 100. The eigenvalues of the SE covariance kernel decay much more rapidly than those of the AE covariance kernel. Because we choose a fixed number of terms \(n_\nu = 5\), the random viscosity with the SE covariance kernel retains a smaller variance.

Figure 4.10(a) shows the computational costs (in seconds) needed for computing the full-rank solutions and the low-rank approximate solutions using the inexact nonlinear iteration for the two covariance kernels and a set of correlation lengths, \(l_1 = l_2 = \{1, 2, 4, 8, 16, 32\}\). Figure 4.10(b) shows the ranks of the low-rank approx-
imate solutions that satisfy the nonlinear stopping tolerance $\epsilon_{nl} = 10^{-5}$. Again, $Re_0 = 100$ and $CoV = 1\%$. For this benchmark problem, 4 Picard iterations and 1 Newton iteration are enough to generate a nonlinear iterate satisfying the stopping tolerance $\epsilon_{nl}$. It can be seen from Figure 4.10(a) that in all cases the use of low rank methods reduces computational cost. Moreover, as the correlation length becomes larger, the ranks of the corrections and the nonlinear iterates become smaller. As a result, the low-rank method achieves greater computational savings for the problems with larger correlation length.

Figure 4.9: Eigenvalue decay of the AE and the SE covariance kernels.

Figure 4.10: Computational costs and ranks for varying correlation lengths with SE and AE covariance kernel.
Next, we examine the performances of the low-rank approximation method for varying CoV, which is defined in (4.3). In this experiment, we fix the value of \( \text{Re}_0 = 100 \) and the variance of the random \( \sigma_\nu \) is controlled. We consider the SE covariance kernel.

![Graphs showing computational costs and ranks for varying correlation lengths and varying CoV with \( \text{Re}_0 = 100 \).]

Figure 4.11: Computational costs and ranks for varying correlation lengths and varying CoV with \( \text{Re}_0 = 100 \).

Figure 4.11 shows the performances of the full-rank and the low-rank methods for varying CoV = \{1%, 5%, 10%\}. We use Algorithm 6 with 4 Picard steps,
followed by several Newton steps until convergence. For \( CoV = \{1\%, 5\%\} \), one Newton step is required for the convergence and, for \( CoV = 10\% \), two Newton steps are required. Figure 4.11(a) shows the computational costs. For \( CoV = \{1\%, 5\%\} \), the computational benefits of using the low-rank approximation methods are pronounced whereas, for \( CoV = 10\% \), the performances of the two approaches are essentially the same for shorter correlation lengths. Indeed, for higher \( CoV \), the ranks of solutions \( \bar{u} \) (see Figures 4.11(b)–4.11(d)) as well as updates \( \delta \bar{u}^k \) at Newton steps become close to the full rank \( (n_\xi = 56) \).

Lastly, we study the benchmark problems with varying mean viscosity with SE covariance kernel and \( CoV = 1\% \). As the mean viscosity decreases, \( Re_0 \) grows, and the nonlinear problem tends to become harder to solve, and for the larger Reynolds numbers \( Re_0 = 200 \) or 300, we use more Picard steps (5 or 6, respectively) before switching to Newton’s method.

Figure 4.12 shows the performances of the low-rank methods for varying Reynolds number, \( Re_0 = \{100, 200, 300\} \). For \( Re_0 = 200 \), after 5 Picard steps, one Newton step leads to convergence (and 6 Picard steps and one Newton step for \( Re_0 = 300 \)). As the figures 4.12(b)–4.12(d) show, the ranks of the solutions increase slightly as the Reynolds number becomes larger and, thus, for all \( Re_0 \) tested here, the low-rank method demonstrates notable computational savings (with \( CoV = 1\% \)). Note that overall computational costs in Figure 4.12(a) increase as the Reynolds number becomes larger because (1) the number of nonlinear steps required to converge increases as the Reynolds number increases and (2) to solve each linearized systems, typically more lrGMRES cycles are required for the problems with higher
4.6 Conclusion

In this study, we have developed the inexact low-rank nonlinear iteration for the solutions of the Navier–Stoke equations with uncertain viscosity in the stochastic Galerkin context. At each step of the nonlinear iteration, the solution of the linear
system is inexpensively approximated in low rank using the tensor variant of the GMRES method. We examined the effect of the truncation to an accuracy of the low-rank approximate solutions by comparing those solutions to the ones computed using exact, inexact nonlinear iterations in full rank and the Monte Carlo method. Then we explored the efficiency of the proposed method with a set of benchmark problems for various settings of uncertain viscosity. The numerical experiments demonstrated that the low-rank nonlinear iteration achieved significant computational savings for the problems with smaller $CoV$ and larger correlation lengths. The experiments also showed that the mean Reynolds number does not significantly affect the rank of the solution and the low-rank nonlinear iteration achieves computational savings for varying Reynolds number for small $CoV$ and large correlation lengths.
Chapter 5: Stochastic Least-Square Petrov Galerkin method

5.1 Introduction

In this chapter, we consider the issues of optimality associated with the stochastic Galerkin method. The stochastic Galerkin method combined with generalized polynomial chaos (gPC) expansions [112] seeks a polynomial approximation of the numerical solution in the stochastic domain by enforcing a Galerkin orthogonality condition, i.e., the residual of the parameterized linear system is forced to be orthogonal to the span of the stochastic polynomial basis with respect to an inner product associated with an underlying probability measure. The Galerkin projection scheme is popular for its simplicity (i.e., the trial and test bases are the same) and its optimality in terms of minimizing an energy norm of solution errors when the underlying PDE operator is symmetric positive definite. In many applications, however, the stochastic Galerkin method does not exhibit any optimality property [73]. That is, it does not produce solutions that minimize any measure of the solution error. In such cases, the stochastic Galerkin method can lead to poor approximations and non-convergent behavior.

To address this issue, we propose a novel optimal projection technique, which we refer to as the stochastic least-squares Petrov–Galerkin (LSPG) method. Inspired
by the successes of LSPG methods in nonlinear model reduction \cite{18-20}, finite element methods \cite{13,14,55}, and iterative linear solvers (e.g., GMRES, GCR) \cite{91}, we propose, as an alternative to enforcing the Galerkin orthogonality condition, to directly minimize the residual of a parameterized linear system over the stochastic domain in a (weighted) $\ell^2$-norm. The stochastic LSPG method produces an optimal solution for a given stochastic subspace and guarantees that the $\ell^2$-norm of the residual monotonically decreases as the stochastic basis is enriched. In addition to producing monotonically convergent approximations as measured in the chosen weighted $\ell^2$-norm, the method can also be adapted to target output quantities of interest (QoI); this can be accomplished by employing a weighted $\ell^2$-norm used for least-squares minimization that coincides with the $\ell^2$-(semi)norm of the error in the chosen QoI.

In addition to proposing the stochastic LSPG method, this study shows that specific choices of weighting functions lead to equivalences between the stochastic LSPG method and both the stochastic Galerkin method and the pseudo-spectral method \cite{109,110}. We demonstrate the effectiveness of the LSPG method with extensive numerical experiments on various SPDEs. The results show that the proposed LSPG technique significantly outperforms the stochastic Galerkin when the solution error is measured in different weighted $\ell^2$-norms. We also show that the proposed method can effectively minimize the error in target QoIs.

An outline of the chapter is as follows. Section 5.2 formulates parameterized linear algebraic systems and reviews conventional spectral approaches for computing numerical solutions. Section 5.3 develops a residual minimization formulation
based on least-squares methods and its adaptation to the stochastic LSPG method. We also provide proofs of optimality and monotonic convergence behavior of the proposed method. Section 5.4 provides error analysis for stochastic LSPG methods. Section 5.5 demonstrates the efficiency and the effectiveness of the proposed methods by testing them on various benchmark problems. Finally, Section 5.6 outlines some conclusions.

5.2 Spectral methods for parameterized linear systems

We begin by introducing a mathematical formulation of parameterized linear systems and briefly reviewing the stochastic Galerkin and the pseudo-spectral methods, which are spectral methods for approximating the numerical solutions of such systems.

5.2.1 Problem formulation

Consider a parameterized linear system

\[ A(\xi)u(\xi) = b(\xi), \]  

(5.1)

where \( A : \Gamma \rightarrow \mathbb{R}^{n \times n} \), and \( u, b : \Gamma \rightarrow \mathbb{R}^{n} \). The system is parameterized by a set of stochastic input parameters \( \xi(\omega) \equiv \{\xi_{1}(\omega), \ldots, \xi_{n_{\xi}}(\omega)\} \). Here, \( \omega \in \Omega \) is an elementary event in a probability space \((\Omega, \mathcal{F}, P)\) and the stochastic domain is denoted by \( \Gamma \equiv \prod_{i=1}^{n_{\xi}} \Gamma_{i} \) where \( \xi_{i} : \Omega \rightarrow \Gamma_{i} \). We are interested in computing solutions in finite-dimensional subspaces of \( L^{2}(\Gamma) \) (defined below) using weak formulations of
(5.1) corresponding to Galerkin and Petrov–Galerkin projections.

Let \( \rho \equiv \rho(\xi) \) be a density function defining an underlying measure of the stochastic space \( \Gamma \) and

\[
\langle g, h \rangle_{\rho} \equiv \int_{\Gamma} g(\xi) h(\xi) \rho(\xi) d\xi, \tag{5.2}
\]

\[
E[g] \equiv \int_{\Gamma} g(\xi) \rho(\xi) d\xi, \tag{5.3}
\]

define an inner product between scalar-valued functions \( g(\xi) \) and \( h(\xi) \) with respect to \( \rho(\xi) \) and the expectation of \( g(\xi) \), respectively. The inner product (5.2) also determines the Hilbert space \( L^2(\Gamma) \). In addition, the \( \ell^2 \)-norm of a vector-valued function \( v(\xi) \in \mathbb{R}^{n_x} \) is defined as

\[
\|v\|^2_2 \equiv \sum_{i=1}^{n_x} \int_{\Gamma} v_i^2(\xi) \rho(\xi) d\xi = E[v^T v]. \tag{5.4}
\]

We are interested in computing approximate solutions to (5.1) using spectral methods, that is, finding solutions in an \( n_\psi \)-dimensional subspace \( S_{n_\psi} \) spanned by a finite set of polynomials \( \{\psi_i(\xi)\}_{i=1}^{n_\psi} \) such that \( S_{n_\psi} \equiv \text{span}\{\psi_i\}_{i=1}^{n_\psi} \subseteq L^2(\Gamma) \). Then

\[
u(\xi) \approx \tilde{u}(\xi) = \sum_{i=1}^{n_\psi} \tilde{u}_i \psi_i(\xi) = (\psi^T(\xi) \otimes I_{n_x}) \tilde{u}, \tag{5.5}
\]

where \( \{\tilde{u}_i\}_{i=1}^{n_\psi} \) with \( \tilde{u}_i \in \mathbb{R}^{n_x} \) are unknown coefficient vectors, \( \tilde{u} \equiv [\tilde{u}_1^T \cdots \tilde{u}_{n_\psi}^T]^T \in \mathbb{R}^{n_x n_\psi} \) is the vertical concatenation of these coefficient vectors, \( \psi \equiv [\psi_1 \cdots \psi_{n_\psi}]^T \in \mathbb{R}^{n_\psi} \) is a concatenation of the polynomial basis, \( \otimes \) denotes the Kronecker prod-
uct, and $I_{n_x}$ denotes the identity matrix of dimension $n_x$. Note that $\tilde{u} \in (S_{n_\psi})^{n_x}$.

Typically, the “stochastic” basis $\{\psi_i\}$ consists of products of univariate polynomials: $\psi_i \equiv \psi_{\alpha(i)} \equiv \prod_{k=1}^{n_\xi} \pi_{\alpha_k(i)}(\xi_k)$ where $\{\pi_{\alpha_k(i)}\}_{k=1}^{n_\xi}$ are univariate polynomials, $\alpha(i) = (\alpha_1(i), \cdots, \alpha_{n_\xi}(i)) \in \mathbb{N}_0^{n_\xi}$ is a multi-index and $\alpha_k$ represents the degree of a polynomial in $\xi_k$. The dimension of the stochastic subspace $n_\psi$ depends on the number of random variables $n_\xi$, the maximum polynomial degree $p$, and a construction of the polynomial space (e.g., a total-degree space that contains polynomials with total degree up to $p$, $\sum_{k=1}^{n_\xi} \alpha_k(i) \leq p$). By substituting $u(\xi)$ with $\tilde{u}(\xi)$ in (5.1), the residual can be defined as

$$r(\tilde{u}; \xi) := b(\xi) - A(\xi) \sum_{i=1}^{n_\psi} \tilde{u}_i \psi_i(\xi) = b(\xi) - (\psi^T(\xi) \otimes A(\xi)) \tilde{u}, \quad (5.6)$$

where $\psi^T(\cdot) \otimes A(\cdot) : \Gamma \rightarrow \mathbb{R}^{n_x \times n_\psi n_x}$.

It follows from (5.5) and (5.6) that our goal now is to compute the unknown coefficients $\{\tilde{u}_i\}_{i=1}^{n_\psi}$ of the solution expansion. We briefly review two conventional approaches for doing so: the stochastic Galerkin method and the pseudo-spectral method. Typically, the polynomial basis is constructed to be orthogonal in the $\langle \cdot, \cdot \rangle_\rho$ inner product, i.e., $\langle \psi_i, \psi_j \rangle_\rho = \prod_{k=1}^{n_\xi} \langle \pi_{\alpha_k(i)}, \pi_{\alpha_k(j)} \rangle_\rho_k = \delta_{ij}$, where $\delta_{ij}$ denotes the Kronecker delta.

### 5.2.2 Stochastic Galerkin method

The stochastic Galerkin method computes the unknown coefficients $\{\tilde{u}_i\}_{i=1}^{n_\psi}$ of $\tilde{u}(\xi)$ in (5.5) by imposing orthogonality of the residual (5.6) with respect to the
inner product $\langle \cdot, \cdot \rangle_\rho$ in the subspace $S_{n_\psi}$. This Galerkin orthogonality condition can be expressed as follows: Find $\bar{u}^{SG} \in \mathbb{R}^{n_x n_\psi}$ such that

$$\langle r_i(\bar{u}^{SG}), \psi_j \rangle_\rho = E[r_i(\bar{u}^{SG})\psi_j] = 0, \quad i = 1, \ldots, n_x, \ j = 1, \ldots, n_\psi, \quad (5.7)$$

where $r \equiv [r_1 \cdots r_{n_x}]^T$. The condition (5.7) can be represented in matrix notation as

$$E[\psi \otimes r(\bar{u}^{SG})] = 0. \quad (5.8)$$

From the definition of the residual (5.6), this gives a system of linear equations

$$E[\psi \psi^T \otimes A]\bar{u}^{SG} = E[\psi \otimes b], \quad (5.9)$$

of dimension $n_x n_\psi$. This yields an algebraic expression for the stochastic-Galerkin approximation

$$\bar{u}^{SG}(\xi) = (\psi(\xi)^T \otimes I_{n_x})E[\psi \psi^T \otimes A]^{-1}E[\psi \otimes Au]. \quad (5.10)$$

If $A(\xi)$ is symmetric positive definite, the solution of linear system (5.9) minimizes the solution error $e(x) \equiv u - x$ in the $A(\xi)$-induced energy norm $\|v\|_A^2 \equiv E[v^T A v]$, i.e.,

$$\bar{u}^{SG}(\xi) = \text{arg min}_{x \in (S_{n_\psi})_{n_x}} \|e(x)\|_A^2. \quad (5.11)$$

In general, however, the stochastic-Galerkin approximation does not minimize any measure of the solution error.
5.2.3 Pseudo-spectral method

The pseudo-spectral method directly approximates the unknown coefficients \( \{\bar{u}_i\}_{i=1}^{n_\psi} \) of \( \tilde{u}(\xi) \) in (5.5) by exploiting orthogonality of the polynomial basis \( \{\psi_i(\xi)\}_{i=1}^{n_\psi} \).

That is, the coefficients \( \bar{u}_i \) can be obtained by projecting the numerical solution \( u(\xi) \) onto the orthogonal polynomial basis as

\[
\bar{u}_i^{PS} = E[u\psi_i], \quad i = 1, \ldots, n_\psi, \tag{5.12}
\]

which can be expressed as

\[
\tilde{u}^{PS} = E[\psi \otimes A^{-1}b], \tag{5.13}
\]

or equivalently

\[
\tilde{u}^{PS}(\xi) = (\psi(\xi)^T \otimes I_{n_x})E[\psi \otimes u]. \tag{5.14}
\]

The associated optimality property of the approximation, which can be derived from optimality of orthogonal projection, is

\[
\tilde{u}^{PS}(\xi) = \arg \min_{x \in (S_{n_\psi})^{n_x}} \|e(x)\|_2^2. \tag{5.15}
\]

In practice, the coefficients \( \{\bar{u}_i^{PS}\}_{i=1}^{n_\psi} \) are approximated via numerical quadrature as

\[
\bar{u}_i^{PS} = E[u\psi_i] = \sum_{k=1}^{n_q} u(\xi^{(k)})\psi_i(\xi^{(k)})w_k = \sum_{k=1}^{n_q} (A^{-1}(\xi^{(k)})b(\xi^{(k)})) \psi_i(\xi^{(k)})w_k, \tag{5.16}
\]

where \( \{(\xi^{(k)}, w_k)\}_{k=1}^{n_q} \) are the quadrature points and weights.
While stochastic Galerkin leads to an optimal approximation \((5.11)\) under certain conditions and pseudo-spectral projection minimizes the \(\ell^2\)-norm of the solution error \((5.15)\), neither approach provides the flexibility to tailor the optimality properties of the approximation. This may be important in applications where, for example, minimizing the error in a quantity of interest is desired. To address this, we propose a general optimization-based framework for spectral methods that enables the choice of a targeted weighted \(\ell^2\)-norm in which the solution error is minimized.

5.3 Stochastic least-squares Petrov–Galerkin method

As a starting point, we propose a residual-minimizing formulation that computes the coefficients \(\tilde{u}\) by directly minimizing the \(\ell^2\)-norm of the residual, i.e.,

\[
\tilde{u}_{\text{LSPG}}(\xi) = \arg \min_{x \in (S_{n_{\psi}})^{nx}} \|b - Ax\|_2^2 = \arg \min_{x \in (S_{n_{\psi}})^{nx}} \|e(x)\|_{A^TA}^2, \quad (5.17)
\]

where \(\|v\|_{A^TA}^2 \equiv E[v^TA^TAv]\). Thus, the \(\ell^2\)-norm of the residual is equivalent to a weighted \(\ell^2\)-norm of the solution error. Using \((5.5)\) and \((5.6)\), we have

\[
\tilde{u}_{\text{LSPG}} = \arg \min_{\bar{x} \in \mathbb{R}^{nx \times n_{\psi}}} \|r(\bar{x})\|_2^2. \quad (5.18)
\]
The definition of the residual (5.6) allows the objective function in (5.18) to be written in quadratic form as

\[ \|r(\bar{x})\|^2 = \|b - (\psi^T \otimes A)\bar{x}\|^2 = \bar{x}^T E[\psi \psi^T \otimes A^T A] \bar{x} - 2E[\psi \otimes A^T b]^T \bar{x} + E[b^T b]. \]  

(5.19)

Noting that the mapping \( \bar{x} \mapsto \|r(\bar{x})\|^2 \) is convex, the (unique) solution \( \tilde{u}_{\text{LSPG}} \) to (5.18) is a stationary point of \( \|r(\bar{x})\|^2 \) and thus satisfies

\[ E[\psi \psi^T \otimes A^T A] \tilde{u}_{\text{LSPG}} = E[\psi \otimes A^T b], \]  

(5.20)

which can be interpreted as the normal-equations form of the linear least-squares problem (5.18).

Consider a generalization of this idea that minimizes the solution error in a targeted weighted \( \ell^2 \) -norm by choosing a specific weighting function. Let us define a weighting function \( M(\xi) \equiv M_\xi(\xi) \otimes M_x(\xi) \), where \( M_\xi : \Gamma \to \mathbb{R} \) and \( M_x : \Gamma \to \mathbb{R}^{n_x \times n_x} \).

Then, the stochastic LSPG method can be written as

\[ \tilde{u}_{\text{LSPG}(M)}(\xi) = \arg \min_{x \in (S_{n_\psi})^{n_x}} \|M(b - Ax)\|^2 = \arg \min_{x \in (S_{n_\psi})^{n_x}} \|e(x)\|^2_{A^T M^T A}, \]  

(5.21)

with \( \|v\|^2_{A^T M^T A} \equiv E[v^T A^T M^T M A v] = E[(M_\xi^T M_x \otimes (M_x A v)^T M_x A v)]. \) Algebraically,
this is equivalent to

\[
\bar{u}^\text{LSPG}(M) = \arg\min_{\bar{x} \in \mathbb{R}^{n_x \times n_x}} \|Mr(\bar{x})\|^2_2 = \arg\min_{\bar{x} \in \mathbb{R}^{n_x \times n_x}} \| (M_\xi \otimes M_x)(1 \otimes b - (\psi^T \otimes A) \bar{x})\|^2_2 \\
= \arg\min_{\bar{x} \in \mathbb{R}^{n_x \times n_x}} \| M_\xi \otimes (M_x b) - ((M_\xi \psi^T) \otimes (M_x A)) \bar{x}\|^2_2.
\] (5.22)

We will restrict our attention to the case \(M_\xi(\xi) = 1\) and denote \(M_x(\xi)\) by \(M(\xi)\) for simplicity. Now, the algebraic stochastic LSPG problem (5.22) simplifies to

\[
\bar{u}^\text{LSPG}(M) = \arg\min_{\bar{x} \in \mathbb{R}^{n_x \times n_x}} \|Mb - (\psi^T \otimes MA)\bar{x}\|^2_2.
\] (5.23)

The objective function in (5.23) can be written in quadratic form as

\[
\|Mr(\bar{x})\|^2_2 = \bar{x}^T E[(\psi^T \otimes A^T M^T M A)] \bar{x} - 2(E[\psi \otimes A^T M^T M f])^T \bar{x} + E[b^T M^T M b].
\] (5.24)

As before, because the mapping \(\bar{x} \mapsto \|Mr(\bar{x})\|^2_2\) is convex, the unique solution \(\bar{u}^\text{LSPG}(M)\) of (5.23) corresponds to a stationary point of \(\|Mr(\bar{x})\|^2_2\) and thus satisfies

\[
E[\psi^T \otimes A^T M^T M A] \bar{u}^\text{LSPG}(M) = E[\psi \otimes A^T M^T M f],
\] (5.25)

which is the normal-equations form of the linear least-squares problem (5.23). This yields the following algebraic expression for the stochastic-LSPG approximation:

\[
\bar{u}^\text{LSPG}(M)(\xi) = (\psi(\xi)^T \otimes I_{n_x}) E[\psi^T \otimes A^T M^T M A]^{-1} E[\psi \otimes A^T M^T M Au].
\] (5.26)
Petrov–Galerkin projection. Another way of interpreting the normal equations \((5.25)\) is that the (weighted) residual \(M(\xi)r(\hat{u}^{\text{LSPG}(M)}; \xi)\) is enforced to be orthogonal to the subspace spanned by the optimal test basis \(\{\phi_i\}_{i=1}^{n_\psi}\) with \(\phi_i(\xi) := \psi_i(\xi) \otimes M(\xi)A(\xi)\) and \(\text{span}\{\phi_i\}_{i=1}^{n_\psi} \subseteq L^2(\Gamma)\). That is, this projection is precisely the (least-squares) Petrov–Galerkin projection,

\[
E[\phi^T (b - (\psi^T \otimes MA)\hat{u}^{\text{LSPG}(M)})] = 0, \tag{5.27}
\]

where \(\phi(\xi) \equiv [\phi_1(\xi) \cdots \phi_{n_\psi}(\xi)]\).

Monotonic Convergence. The stochastic least-squares Petrov-Galerkin is monotonically convergent. That is, as the trial subspace \(S_{n_\psi}\) is enriched (by adding polynomials to the basis), the optimal value of the convex objective function \(\|Mr(\hat{u}^{\text{LSPG}(M)})\|_2\) monotonically decreases. This is apparent from the LSPG optimization problem \((5.21)\): Defining

\[
\hat{u}^{\text{LSPG'}(M)}(\xi) = \arg\min_{x \in (S_{n_\psi+1})^{n_x}} \|M(b - Ax)\|_2^2, \tag{5.28}
\]

we have \(\|M(b - A\hat{u}^{\text{LSPG'}(M)})\|_2^2 \leq \|M(b - A\hat{u}^{\text{LSPG}(M)})\|_2^2\) (and \(\|u - u^{\text{LSPG'}(M)}\|_{A^T M^T M A} \leq \|u - u^{\text{LSPG}(M)}\|_{A^T M^T M A}\)) if \(S_{n_\psi} \subseteq S_{n_\psi+1}\).

Weighting strategies. Different choices of weighting function \(M(\xi)\) allow LSPG to minimize different measures of the error. We focus on four particular choices:

1. \(M(\xi) = C^{-1}(\xi)\), where \(C(\xi)\) is a Cholesky factor of \(A(\xi)\), i.e., \(A(\xi) = C(\xi)C^T(\xi)\).
This decomposition exists if and only if $A$ is symmetric positive semidefinite. In this case, LSPG minimizes the energy norm of the solution error $\|e(x)\|^2_A \equiv \|C^{-1}r(\bar{x})\|^2_2$ ($= \|e((\Psi^T \otimes I_{n_x})\bar{x})\|^2_A$) and is mathematically equivalent to the stochastic Galerkin method described in Section 5.2.2, i.e., $	ilde{u}^{\text{LSPG}(C^{-1})} = \tilde{u}^{\text{SG}}$. This can be seen by comparing (5.11) and (5.21) with $M = C^{-1}$, as $A^T M^T M A = A$ in this case.

2. $M(\xi) = I_{n_x}$, where $I_{n_x}$ is the identity matrix of dimension $n_x$. In this case, LSPG minimizes the $\ell^2$-norm of the residual $\|e(x)\|_{A^T A} \equiv \|r(\bar{x})\|^2_2$.

3. $M(\xi) = A^{-1}(\xi)$. In this case, LSPG minimizes the $\ell^2$-norm of solution error $\|e(x)\|^2_2 \equiv \|A^{-1}r(\bar{x})\|^2_2$. This is mathematically equivalent to the pseudo-spectral method described in Section 5.2.3, i.e., $	ilde{u}^{\text{LSPG}(A^{-1})} = \tilde{u}^{\text{PS}}$, which can be seen by comparing (5.15) and (5.21) with $M = A^{-1}$.

4. $M(\xi) = F(\xi)A^{-1}(\xi)$ where $F : \Gamma \rightarrow \mathbb{R}^{n_o \times n_x}$ is a linear functional of the solution associated with a vector of $n_o$ output quantities of interest. In this case, LSPG minimizes the $\ell^2$-norm of the error in the output quantities of interest $\|Fe(x)\|^2_2 \equiv \|FA^{-1}r(\bar{x})\|^2_2$.

We again emphasize that two particular choices of the weighting function $M(\xi)$ lead to equivalence between LSPG and existing spectral-projection methods (stochastic Galerkin and pseudo-spectral projection), i.e.,

\[ \tilde{u}^{\text{LSPG}(C^{-1})} = \tilde{u}^{\text{SG}}, \quad \tilde{u}^{\text{LSPG}(A^{-1})} = \tilde{u}^{\text{PS}}, \quad \text{(5.29)} \]
where the first equality is valid (i.e., the Cholesky decomposition \( A(\xi) = C(\xi)C^T(\xi) \) can be computed) if and only if \( A \) is symmetric positive semidefinite. Table 5.1 summarizes the target quantities to minimize (i.e., \( \|e(x)\|^2_\Theta \equiv E[e(x)^T \Theta e(x)] \)), the corresponding LSPG weighting functions, and the method names LSPG(\( \Theta \)).

<table>
<thead>
<tr>
<th>Quantity minimized by LSPG</th>
<th>Expression</th>
<th>Weighting function</th>
<th>Method name</th>
</tr>
</thead>
<tbody>
<tr>
<td>Energy norm of error</td>
<td>( |e(x)|^2_A )</td>
<td>( M(\xi) = C^{-1}(\xi) )</td>
<td>LSPG(A)/SG</td>
</tr>
<tr>
<td>( \ell^2 )-norm of residual</td>
<td>( |e(x)|^2_{A^T A} )</td>
<td>( M(\xi) = I_{n_x} )</td>
<td>LSPG(( A^T A ))</td>
</tr>
<tr>
<td>( \ell^2 )-norm of solution error</td>
<td>( |e(x)|^2_2 )</td>
<td>( M(\xi) = A^{-1}(\xi) )</td>
<td>LSPG(2)/PS</td>
</tr>
<tr>
<td>( \ell^2 )-norm of error in quantities of interest</td>
<td>( |Fe(x)|^2_2 )</td>
<td>( M(\xi) = F(\xi)A^{-1}(\xi) )</td>
<td>LSPG(( F^T F ))</td>
</tr>
</tbody>
</table>

### 5.4 Error analysis

If an approximation satisfies an optimal-projection condition

\[
\tilde{u} = \arg\min_{x \in (S_{\psi})^{n_x}} \|e(x)\|^2_\Theta, \tag{5.30}
\]

then

\[
\|e(\tilde{u})\|^2_\Theta = \min_{x \in (S_{\psi})^{n_x}} \|e(x)\|^2_\Theta. \tag{5.31}
\]

Using norm equivalence

\[
\|x\|^2_{\Theta'} \leq C\|x\|^2_\Theta, \tag{5.32}
\]
we can characterize the solution error $e(\tilde{u})$ in any alternative norm $\Theta'$ as

$$\|e(\tilde{u})\|^2_{\Theta'} \leq C \min_{x \in (S_{n\psi})^{n_x}} \|e(x)\|^2_{\Theta}.$$  \hspace{1cm} (5.33)

Thus, the error in an alternative norm $\Theta'$ is controlled by the optimal objective-function value $\min_{x \in (S_{n\psi})^{n_x}} \|e(x)\|^2_{\Theta}$ (which can be made small if the trial space admits accurate solutions) and the stability constant $C$.

Table 5.2 reports norm-equivalence constants for the norms considered in this work. Here, we have defined

$$\sigma_{\text{min}}(M) \equiv \inf_{x \in (L^2(\Gamma))^{n_x}} \|Mx\|_2/\|x\|_2, \quad \sigma_{\text{max}}(M) \equiv \sup_{x \in (L^2(\Gamma))^{n_x}} \|Mx\|_2/\|x\|_2.$$ \hspace{1cm} (5.34)

<table>
<thead>
<tr>
<th>$\Theta'$</th>
<th>$\Theta = A$</th>
<th>$\Theta' = A^T A$</th>
<th>$\Theta' = 2$</th>
<th>$\Theta' = F^TF$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\Theta = A$</td>
<td>1</td>
<td>$\sigma_{\text{max}}(A)$</td>
<td>$\frac{1}{\sigma_{\text{min}}(A)}$</td>
<td>$\frac{\sigma_{\text{max}}(F)^2}{\sigma_{\text{min}}(A)}$</td>
</tr>
<tr>
<td>$\Theta = A^T A$</td>
<td>$\frac{1}{\sigma_{\text{min}}(A)}$</td>
<td>$\sigma_{\text{max}}(A)$</td>
<td>$\frac{1}{\sigma_{\text{max}}(A)}$</td>
<td>$\frac{\sigma_{\text{max}}(F)^2}{\sigma_{\text{max}}(A)}$</td>
</tr>
<tr>
<td>$\Theta = 2$</td>
<td>$\sigma_{\text{max}}(A)^2$</td>
<td>$\frac{\sigma_{\text{max}}(A)^2}{\sigma_{\text{min}}(F)^2}$</td>
<td>$\frac{1}{\sigma_{\text{max}}(A)^2}$</td>
<td>$\frac{1}{\sigma_{\text{max}}(F)^2}$</td>
</tr>
<tr>
<td>$\Theta = F^TF$</td>
<td>$\sigma_{\text{min}}(F)^2$</td>
<td>$\frac{\sigma_{\text{max}}(F)^2}{\sigma_{\text{min}}(F)^2}$</td>
<td>$\frac{1}{\sigma_{\text{min}}(F)^2}$</td>
<td>1</td>
</tr>
</tbody>
</table>

This exposes several interesting conclusions. First, if the number of output quantities of interest $n_o$ is less than $n_x$, then the null space of $F$ is nontrivial and so $\sigma_{\text{min}}(F) = 0$. This implies that LSPG($F^TF$), for which $\Theta = F^TF$, will have an undefined value of $C$ when the solution error is measured in other norms, i.e., for $\Theta' = A$, $\Theta' = A^T A$, and $\Theta' = 2$. It will have controlled errors only for $\Theta' = F^TF$, in
which case $C = 1$. Second, note that for problems with small $\sigma_{\text{min}}(A)$, the $\ell^2$ norm in the quantities of interest may be large for the LSPG($A$)/SG, or LSPG($A^T A$), while it will remain well behaved for LSPG(2)/PS and LSPG($F^T F$).

5.5 Numerical experiments

This section explores the performance of the LSPG methods for solving elliptic SPDEs parameterized by one random variable (i.e., $n_\xi = 1$). The maximum polynomial degree used in the stochastic space $S_{n_\psi}$ is $p$; thus, the dimension of $S_{n_\psi}$ is $n_\psi = p + 1$. In physical space, the SPDE is defined over a two-dimensional rectangular bounded domain $D$, and it is discretized using the finite element method with bilinear ($Q_1$) elements as implemented in the Incompressible Flow and Iterative Solver Software (IFISS) package [98]. Sixteen elements are employed in each dimension, leading to $n_x = 225 = 15^2$ degrees of freedom excluding boundary nodes. All numerical experiments are performed on an Intel 3.1 GHz i7 CPU, 16 GB RAM using MATLAB R2015a.

**Measuring weighted $\ell^2$-norms.** For all LSPG methods, the weighted $\ell^2$-norms can be measured by evaluating the expectations in the quadratic form of the objective function shown in (5.24). This requires evaluation of three expectations

$$
\|Mr(\bar{x})\|_2^2 := \bar{x}^T T_1 \bar{x} - 2T_2^T \bar{x} + T_3,
$$

(5.35)
with

\[ T_1 := E[(\psi \psi^T \otimes A^T M^T M^T A)] \in \mathbb{R}^{n_x n_\psi \times n_x n_\psi}, \quad (5.36) \]

\[ T_2 := E[\psi \otimes A^T M^T M b] \in \mathbb{R}^{n_x n_\psi}, \quad (5.37) \]

\[ T_3 := E[b^T M^T M b] \in \mathbb{R}. \quad (5.38) \]

Note that \( T_3 \) does not depend on the stochastic-space dimension \( n_\psi \). These quantities can be evaluated by numerical quadrature or analytically if closed-form expressions for those expectations exist. Unless otherwise specified, we compute these quantities using the \texttt{integral} function in \textsc{matlab}, which performs adaptive numerical quadrature based on the 15-point Gauss–Kronrod quadrature formula \cite{95}.

**Error measures.** In the experiments, we assess the error in approximate solutions computed using various spectral-projection techniques using four relative error measures (see Table 5.1):

\[ \eta_c(x) := \frac{\|e(x)\|_A^2}{\|b\|_2^2}, \quad \eta_e(x) := \frac{\|e(x)\|_2^2}{\|u\|_2^2}, \quad \eta_A(x) := \frac{\|e(x)\|_A^2}{\|u\|_A^2}, \quad \eta_Q(x) := \frac{\|F e(x)\|_2^2}{\|F u\|_2^2}. \quad (5.39) \]
5.5.1 Stochastic diffusion problems

Consider the steady-state stochastic diffusion equation with homogeneous boundary conditions,

\[
\begin{aligned}
-\nabla \cdot (a(x, \xi) \nabla u(x, \xi)) &= f(x, \xi) \quad \text{in } D \times \Gamma \\
\quad u(x, \xi) &= 0 \quad \text{on } \partial D \times \Gamma,
\end{aligned}
\]

(5.40)

where the diffusivity \( a(x, \xi) \) is a random field and \( D = [0, 1] \times [0, 1] \). The random field \( a(x, \xi) \) is specified as an exponential of a truncated Karhunen-Loève (KL) expansion \([67]\) with covariance kernel, \( C(x, y) \equiv \sigma^2 \exp \left( -\frac{|x_1-y_1|}{c} - \frac{|x_2-y_2|}{c} \right) \), where \( c \) is the correlation length, i.e.,

\[
a(x, \xi) \equiv \exp(\mu + \sigma a_1(x) \xi),
\]

(5.41)

where \( \{\mu, \sigma^2\} \) are the mean and variance of the KL expansion and \( a_1(x) \) is the first eigenfunction in the KL expansion. After applying the spatial (finite-element) discretization, the problem can be reformulated as a parameterized linear system of the form (5.1), where \( A(\xi) \) is a parameterized stiffness matrix obtained from the weak form of the problem whose \((i, j)\)-element is \([A(\xi)]_{ij} = \int_D \nabla a(x, \xi) \varphi_i(x) \cdot \varphi_j(x) dx\) (with \( \{\varphi_i\} \) standard finite element basis functions) and \( b(\xi) \) is a parameterized right-hand side whose \(i\)th element is \([b(\xi)]_i = \int_D f(x, \xi) \varphi_i(x) dx\). Note that \( A(\xi) \) is symmetric positive definite for this problem; thus LSPG(A)/SG is a valid projection scheme (the Cholesky factorization \( A(\xi) = C(\xi)C(\xi)^T \) exists) and is equal to
stochastic Galerkin projection.

**Output quantities of interest.** We consider \( n_o \) output quantities of interest \((F(\xi)u(\xi) \in \mathbb{R}^{n_o})\) that are random linear functionals of the solution and \( F(\xi) \) is of dimension \( n_o \times n_x \) having the form:

(1) \( F_1(\xi) := g(\xi) \times G \) with \( G \in [0,1]^{n_o \times n_x} \) a constant matrix: The elements of \( G \) are drawn from a uniform distribution (note that this is independent of the distribution \( \rho(\xi) \)) and \( g(\xi) \) is a scalar-valued function of \( \xi \). The resulting output QoI, \( F_1(\xi)u(\xi) \), is a vector-valued function of dimension \( n_o \).

(2) \( F_2(\xi) := b(\xi)^T\bar{M} \): \( \bar{M} \) is a mass matrix defined via \( \bar{M}_{ij} \equiv \int_D \varphi_i(x)\varphi_j(x)dx \).

The output QoI is a scalar-valued function \( F_2(\xi)u(\xi) = b(\xi)^T\bar{M}u(\xi) \), which approximates a spatial average \( \frac{1}{|D|} \int_D f(x, \xi)u(x, \xi)dx \).

5.5.1.1 Diffusion problem 1: Lognormal random coefficient and deterministic forcing

In this example, we take \( \xi \) in (5.41) to follow a standard normal distribution (i.e., \( \rho(\xi) = \frac{1}{\sqrt{2\pi}} \exp\left(-\frac{\xi^2}{2}\right) \) and \( \xi \in (-\infty, \infty) \)) and \( f(x, \xi) = 1 \) is deterministic. Because \( \xi \) is normally distributed, normalized Hermite polynomials (orthogonal with respect to \( \langle \cdot, \cdot \rangle_\rho \)) are used as polynomial basis \( \{\psi_i(\xi)\}_{i=1}^{n_\psi} \).

Figure 5.1 reports the relative errors (5.39) associated with solutions computed using four LSPG methods (LSPG(\( A \))/SG, LSPG(\( A^TA \)), LSPG(2)/PS, and LSPG(\( F^TF \))) for varying polynomial degree \( p \). Here, we consider the random output QoI, i.e., \( F = F_1, n_o = 100, \) and \( g(\xi) = \xi \). This result shows that three methods
Figure 5.1: Relative error measures versus polynomial degree for diffusion problem 1: lognormal random coefficient and deterministic forcing. Note that each LSPG method performs best in the error measure it minimizes.

(LSPG(A)/SG, LSPG($A^T A$), and LSPG(2)/PS) monotonically converge in all four error measures, whereas LSPG($F^T F$) does not. This is an artifact of rank deficiency in $F_1$, which leads to $\sigma_{\text{min}}(F_1) = 0$; as a result, all stability constants $C$ for which $\Theta = F^T F$ in Table 5.2 are unbounded, implying lack of error control. Figure 5.1 also shows that each LSPG method minimizes its targeted error measure for a given stochastic-subspace dimension (e.g., LSPG minimizes the $\ell^2$-norm of the residual);
Figure 5.2: Pareto front of relative error measures versus wall time for varying polynomial degree $p$ ($p$ varies from 1 to 10 in increments of 1 going from left to right) for diffusion problem 1: lognormal random coefficient and deterministic forcing.

This is also evident from Table 5.2, as the stability constant realizes its minimum value ($C = 1$) for $\Theta = \Theta'$. Table 5.3 shows actual values of the stability constant of this problem and well explains the behaviors of all LSPG methods. For example, the first column of Table 5.3 shows that the stability constant is increasing in the order (LSPG$(A)/SG$, LSPG$(A^T A)$, LSPG$(2)/PS$, and LSPG$(F^T F)$), which is represented in Figure 5.1(a).
The results in Figure 5.1 do not account for computational costs. This point is addressed in Figure 5.2, which shows the relative errors as a function of CPU time. As we would like to devise a method that minimizes both the error and computational time, we examine a Pareto front (black dotted line), that is, a curve identifying the methods that minimize the two competing objectives considered in the figure. This typically corresponds to LSPG(2)/PS. This is because this method does not require solution of a coupled system of linear equations of dimension $n_xn_\psi$, which is required by the other three LSPG methods (LSPG(A)/SG, LSPG($A^TA$), and LSPG($F^TF$)). As a result, pseudo-spectral projection (LSPG(2)/PS) generally yields the best overall performance in practice, even when it produces larger errors than other methods for a fixed value of $p$. Also, for a fixed value of $p$, LSPG(A)/SG is faster than LSPG($A^TA$) because the weighted stiffness matrix $A(\xi)$ obtained from the finite element discretization is sparser than $A^T(\xi)A(\xi)$. That is, the number of nonzero entries to be evaluated for LSPG(A)/SG in numerical quadrature is smaller than the ones for LSPG($A^TA$), and exploiting this sparsity structure in the numerical quadrature causes LSPG(A)/SG to be faster than LSPG($A^TA$). Also, note that there are cases (Figure 5.2(b)) where the Pareto front does not correspond to a single method; this outcome will occur with other benchmark problems considered.

Table 5.3: Stability constant $C$ of Diffusion problem 1.

<table>
<thead>
<tr>
<th>$\Theta' = A$</th>
<th>$\Theta' = A^TA$</th>
<th>$\Theta' = 2$</th>
<th>$\Theta' = F^TF$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\Theta = A$</td>
<td>1</td>
<td>2.06</td>
<td>11644.22</td>
</tr>
<tr>
<td>$\Theta = A^TA$</td>
<td>2.06</td>
<td>1</td>
<td>24013.48</td>
</tr>
<tr>
<td>$\Theta = 1$</td>
<td>26.43</td>
<td>4.25</td>
<td>5646.32</td>
</tr>
<tr>
<td>$\Theta = F^TF$</td>
<td>$\infty$</td>
<td>$\infty$</td>
<td>1</td>
</tr>
</tbody>
</table>
5.5.1.2 Diffusion problem 2: Lognormal random coefficient and random forcing

This example uses the same random field \( a(x, \xi) (5.41) \), but instead employs a random forcing term\(^1\) \( f(x, \xi) = \exp(\xi)|\xi - 1| \). Again, \( \xi \) follows a standard normal distribution and normalized Hermite polynomials are used as polynomial basis. We consider the second output QoI, \( F = F_2 \). As shown in Figure 5.3, the stochastic Galerkin method fails to converge monotonically in three error measures as the stochastic polynomial basis is enriched. In fact, it exhibits monotonic convergence only in the error measure it minimizes (for which monotonic convergence is guaranteed).

![Figure 5.3: Relative errors versus polynomial degree for stochastic Galerkin (i.e., LSPG(A)/SG) for diffusion problem 2: lognormal random coefficient and random forcing. Note that monotonic convergence is observed only in the minimized error measure \( \eta_A \).](image)

\(^1\)In [73], it was shown that stochastic Galerkin solutions of an analytic problem \( a(\xi)u(\xi) = f(\xi) \) with this type of forcing are divergent in the \( \ell^2 \)-norm of solution errors as \( p \) increases.
Figure 5.4: Pareto front of relative error measures versus wall time for varying polynomial degree $p$ ($p$ varies from 1 to 20 in increments of 1 going from left to right) for diffusion problem 2: lognormal random coefficient and random forcing

Figure 5.4 shows that this trend applies to other methods as well when effectiveness is viewed with respect to CPU time; each technique exhibits monotonic convergence in its tailored error measure only. Moreover, the Pareto fronts (black dotted lines) in each subgraph of Figure 5.4 shows that the LSPG method tailored for a particular error measure is Pareto optimal in terms of minimizing the error and computational wall time. In the next experiments, we examine
goal-oriented LSPG($F^TF$) for varying number of output quantities of interest $n_o$ and its effect on the stability constant $C$. Figure 5.5 reports three error measures computed using all four LSPG methods. For LSPG($F^TF$), the first linear function $F = F_1$ is applied with $g(\xi) = \sin(\xi)$ and a varying number of outputs $n_o = \{100, 150, 200, 225\}$. When $n_o = 225$, LSPG($F^TF$) and LSPG(2)/PS behave similarly in all three weighted $\ell^2$-norms. This is because when $n_0 = 225 = n_x$, then $\sigma_{\min}(F) > 0$, so the stability constants $C$ for $\Theta = F^TF$ in Table 5.2 are bounded.

Figure 5.6 reports relative errors in the quantity of interest $\eta_Q$ associated with linear functionals $F = F_1$ for two different functions $g(\xi)$, $g_1(\xi) = \sin(\xi)$ and $g_2(\xi) = \xi$. Note that LSPG($A$)/SG and LSPG($A^TA$) fail to converge, whereas LSPG(2)/PS and LSPG($F^TF$) converge, which can be explained by the stability constant $C$ in Table 5.2 where $\sigma_{\max}(A) = 26.43$ and $\sigma_{\min}(A) = 0.48$ for the linear operator $A(\xi)$ of this problem. LSPG($F^TF$) converges monotonically and produces the smallest error (for a fixed polynomial degree $p$) of all the methods as expected.

5.5.1.3 Diffusion problem 3: Gamma random coefficient and random forcing

This section considers a stochastic diffusion problem parameterized by a random variable that has a Gamma distribution, where $a(x, \xi) \equiv \exp(1 + 0.25a_1(x)\xi + 0.01\sin(\xi))$ with density $\rho(\xi) \equiv \frac{e^{\xi} \exp(-\xi)}{\Gamma(\alpha+1)}$, $\Gamma$ is the Gamma function, $\xi \in [0, \infty)$, and $\alpha = 0.5$. Normalized Laguerre polynomials (which are orthogonal with respect to $\langle \cdot, \cdot \rangle_{\rho}$) are used as polynomial basis. We consider a random forcing term
Figure 5.5: Relative error measures versus polynomial degree for a varying dimension $n_o$ of the output matrix $F = F_1$ for diffusion problem 2: lognormal random coefficient and random forcing. Note that LSPG($F^T F$) has controlled errors only when $n_o = n_x$, in which case $\sigma_{\text{min}}(F) > 0$.

\[ f(x, \xi) = \log_{10}(\xi)|\xi - 1| \] and the second QoI $F(\xi) = F_2(\xi) = b(\xi)^T \bar{M}$. Note that numerical quadrature is the only option for computing expectations arise in this problem.

Figure 5.7 shows the results of solving the problem with the four different LSPG methods. Again, each version of LSPG monotonically decreases its corresponding target weighted $\ell^2$-norm as the stochastic basis is enriched. Further, each
(a) Relative $\ell^2$-norm of output QoI error $\eta_Q$ with $F = F_1$, $n_o = 100$, and $g(\xi) = g_1(\xi)$ = \sin(\xi).

(b) Relative $\ell^2$-norm of output QoI error $\eta_Q$ with $F = F_1$, $n_o = 100$, and $g(\xi) = g_2(\xi) = \xi$.

Figure 5.6: Plots of the error norm of output QoI for diffusion problem 2: lognormal random coefficient and random forcing when a linear functional is (a) $F(\xi) \equiv \sin(\xi) \times [0,1]^{100 \times n_x}$ and (b) $F(\xi) = \xi \times [0,1]^{100 \times n_x}$ for varying $p$ and varying $n_o$.

LSPG method is Pareto optimal in terms of minimizing its targeted error measure and the computational wall time.

5.5.2 Stochastic convection-diffusion problem: Lognormal random coefficient and deterministic forcing

We now consider a non-self-adjoint example, the steady-state convection-diffusion equation

$$
\begin{align*}
\begin{cases}
-\epsilon \nabla \cdot (a(x, \xi) \nabla u(x, \xi)) + \vec{w} \cdot \nabla u(x, \xi) &= f(x, \xi) \quad \text{in } D \times \Gamma, \\
\quad u(x, \xi) &= g_D(x) \quad \text{on } \partial D \times \Gamma
\end{cases}
\end{align*}
$$

(5.42)
where $D = [-1, 1] \times [-1, 1]$, $\epsilon$ is the viscosity parameter, and $u$ satisfies inhomogeneous Dirichlet boundary conditions

$$g_D(x) = \begin{cases} 
  g_D(x, 1) = 0 & \text{for } [-1, y] \cup [x, 1] \cup [-1 \leq x \leq 0, -1], \\
  g_D(1, y) = 1 & \text{for } [1, y] \cup [0 \leq x \leq 1, -1].
\end{cases}$$

(5.43)
Figure 5.8: Pareto front of relative error measures versus wall time for varying polynomial degree $p$ ($p$ varies from 1 to 10 in increments of 1 going from left to right) for stochastic convection-diffusion problem: lognormal random coefficient and deterministic forcing term.

The inflow boundary consists of the bottom and the right portions of $\partial D$, $[x, -1] \cup [1, y]$ [39]. We consider a zero forcing term $f(x, \xi) = 0$ and a constant convection velocity $\vec{w} \equiv (-\sin \frac{\pi}{6}, \cos \frac{\pi}{6})$. We consider the convection-dominated case (i.e., $\epsilon = \frac{1}{200}$).

For the spatial discretization, we essentially use the same finite element as
above (bilinear $Q_1$ elements) applied to the weak formulation of (5.42). In addition, we use the streamline-diffusion method [17] to stabilize the discretization in elements with large mesh Peclet number. (See [39], Ch. 8 for details.) Such spatial discretization leads to a parameterized linear system of the form (5.1) with

$$A(\xi) = \epsilon D(a(\xi); \xi) + C(\xi) + S(\xi),$$

(5.44)

where $D(a(\xi); \xi)$, $C(\xi)$ and $S(\xi)$ are the diffusion term, the convection term, and the streamline-diffusion term, respectively, and $[b(\xi)]_i = \int_{D} f(x, \xi) \varphi_i(x) dx$. For this numerical experiment, the number of degrees of freedom in spatial domain is $n_x = 225$ (15 nodes in each spatial dimension) excluding boundary nodes. For LSPG$(F^T F)$, the first linear function $F = F_1$ is applied with $n_o = 100$ outputs and $g(\xi) = \exp(\xi)|\xi - 1|$. Figure 5.8 shows the numerical results computed using the stochastic Galerkin method and three LSPG methods (LSPG$(A^T A)$, LSPG$(2)$/PS, LSPG$(F^T F)$). Note that the operator $A(\xi)$ is not symmetric positive-definite in this case; thus LSPG$(A)$ is not a valid projection scheme (the Cholesky factorization $A(\xi) = C(\xi)C(\xi)^T$ does not exist and the energy norm of the solution error $\|e(x)\|_A^2$ cannot be defined) and stochastic Galerkin does not minimize any measure of the solution error. These results show that pseudo-spectral projection is Pareto optimal for achieving relatively larger error measures; this is because of its relatively low cost since, in contrast to the other methods, it does not require the solution of a coupled linear system of dimension $n_x n_\psi$. In addition, the stochastic Galerkin projection is not
Pareto optimal for any of the examples; this is caused by the lack of optimality of stochastic Galerkin in this case and highlights the significant benefit of optimal spectral projection, which is offered by the stochastic LSPG method. In addition, the residual $\eta_r$ and solution error $\eta_e$ incurred by LSPG($F^TF$) are uncontrolled, because $n_o < n_x$ and thus $\sigma_{\text{min}}(F) = 0$. Finally, note that each LSPG method is Pareto optimal for small errors in its targeted error measure.

5.5.3 Numerical experiment with analytic computations

For the results presented above, expected values were computed using numerical quadrature (using the MATLAB function `integral`). This is a practical and general approach for numerically computing the required integrals of (5.36)–(5.38), and is the only option when analytic computations are not available (as in Section 5.5.1.3). In this section, we briefly discuss how the costs change if analytic methods based on closed-form integration exist and are used for these integrals. Note that in general, however, analytic computation are unavailable, for example, if the random variables have a finite support (e.g., truncated Gaussian random variables as shown in [106]).

Computing $T_1$. Analytic computation of $T_1$ is possible if either $E[A^T M M A \psi_l]$ or $E[M A \psi_l]$ can be evaluated analytically. For LSPG($A$)/SG and LSPG($A^T A$), if $E[A \psi_l]$ can be evaluated so that the following gPC expansion can be obtained analytically

$$A(\xi) = \sum_{l=1}^{\infty} A_l \psi_l(\xi), \quad A_l \equiv E[A \psi_l], \quad (5.45)$$
Figure 5.9: Pareto front of relative error measures versus wall time for varying polynomial degree $p$ ($p$ varies from 1 to 20 in increments of 1 going from left to right) for diffusion problem 2: Lognormal random coefficient and random forcing. Analytic computations are used as much as possible to evaluate expectations.

where $A_l \in \mathbb{R}^{n_x \times n_x}$, then $T_1$ can be computed analytically. Replacing $A(\xi)$ with the series of (5.45) for LSPG$(A)/SG$ ($M(\xi) = C^{-1}(\xi)$) and LSPG$(A^T A)$ ($M(\xi) = I_{n_x}$) yields

$$
T_1^{\text{LSPG}(A)} = \sum_{l=1}^{n_x} E \left[ \psi_l \psi_l^T (A_l \psi_l) \right] = \sum_{l=1}^{n_x} E[\psi_l \psi_l^T \psi_l \otimes A_l],
$$

(5.46)
and

\[ T_1^{\text{LSPG}(A^TA)} = E[\psi_j^T \otimes \sum_{k=1}^{n_a} \sum_{l=1}^{n_a} (A_k \psi_k)^T (A_l \psi_l)] = \sum_{k=1}^{n_a} \sum_{l=1}^{n_a} E[\psi_j^T \psi_k \psi_l \otimes A_k^T A_l], \]

(5.47)

where the expectations of triple or quadruple products of the polynomial basis (i.e., 
\( E[\psi_j \psi_j \psi_k] \) and \( E[\psi_i \psi_j \psi_k \psi_l] \)) can be computed analytically. For LSPG(2)/PS, an analytic computation of \( T_1 \) is straightforward because \( M(\xi)A(\xi) = I_{n_x} \) and, thus,

\[ T_1^{\text{LSPG}(2)} = E[\psi_j^T \otimes I_{n_x}] = I_{n_x n_\psi}. \]

(5.48)

Similarly, analytic computation of \( T_1 \) is possible for LSPG(\( F^T F \)) if there exists a closed formulation for \( E[F \psi_l] \) or \( E[F^T F \psi_l] \), which is again in general not available.

**Computing \( T_2 \).** Analytic computation of \( T_2 \) can be performed in a similar way. If the random function \( b(\xi) \) can be represented using a gPC expansion,

\[ b(\xi) = \sum_{l=1}^{n_b} b_l \psi_l(\xi), \quad b_l = E[b \psi_l], \]

(5.49)

then, for LSPG(\( A \))/SG and LSPG(\( A^TA \)), \( T_2 \) can be evaluated analytically by computing expectations of \( b_i \) or triple products of the polynomial bases (i.e., \( E[\psi_i \psi_j] \) and \( E[\psi_i \psi_j \psi_k] \)). For LSPG(2)/PS and LSPG(\( F^T F \)), however, an analytic computation of \( T_2 \) is typically unavailable because a closed-form expression for \( A^{-1}(\xi) \) does not exist.

We examine the impact of these observations on the cost of solution of the
problem studied in Section 5.5.1.2, a the steady-state stochastic diffusion equation (5.40) with lognormal random field \( a(x, \xi) \) as in (5.41), and random forcing \( f(x, \xi) = \exp(\xi) |\xi - 1| \).

Figure 5.9 reports results for this problem for analytic computation of expectations. For LSPG(\( A \))/SG, analytic computation of the expectations \( \{T_i\}_{i=1}^3 \) requires fewer terms than for LSPG(\( A^T A \)). In fact, comparing (5.46) and (5.47) shows that computing \( T_1^{LSPG(A^T A)} \) requires computing and assembling \( n_a^2 \) terms, whereas computing \( T_1^{LSPG(A)} \) involves only \( n_a \) terms. Additionally the quantities \( \{A_k^T A_l\}_{k,l=1}^{n_a} \) appearing in the terms of \( T_1^{LSPG(A^T A)} \) in (5.47) are typically denser than the counterparts \( \{A_k\}_{k=1}^{n_a} \) appearing in (5.46), as the sparsity pattern of \( \{A_k\}_{k=1}^{n_a} \) is identical to that of the finite element stiffness matrices. As a result, LSPG(\( A \))/SG is Pareto optimal for small computational wall times when any error metric is considered.

When the polynomial degree \( p \) is small, LSPG(\( A \))/SG is computationally faster than LSPG(2)/PS, as LSPG(2)/PS requires the solution of \( A(\xi^{(k)}) u(\xi^{(k)}) = f(\xi^{(k)}) \) at each quadrature point and cannot exploit analytic computation. As the stochastic basis is enriched, however, each tailored LSPG method outperforms other LSPG methods in minimizing its corresponding target error measure.

5.6 Conclusion

In this work, we have proposed a general framework for optimal spectral projection wherein the solution error can be minimized in weighted \( \ell^2 \)-norms of interest. In particular, we propose two new methods that minimize the \( \ell^2 \)-norm of the resid-
ual (LSPG($A^T A$)) and the $\ell^2$-norm of the error in an output quantity of interest (LSPG($F^T F$)). Further, we showed that when the linear operator is symmetric positive definite, stochastic Galerkin is a particular instance of the proposed methodology for a specific choice of weighted $\ell^2$-norm. Similarly, pseudo-spectral projection is a particular case of the method for a specific choice of weighted $\ell^2$-norm.

Key results from the numerical experiments include:

- For a fixed stochastic subspace, each LSPG method minimizes its targeted error measure (Figure 5.1).

- For a fixed computational cost, each LSPG method often minimizes its targeted error measure (Figures 5.4, 5.7). However, this does not always hold, especially for smaller computational costs (and smaller stochastic-subspace dimensions) when larger errors are acceptable. In particular pseudo-spectral projection (LSPG(2)/PS) is often significantly less expensive than other methods for a fixed stochastic subspace, as it does not require solving a coupled linear system of dimension $n_x n_\psi$ (Figures 5.2, 5.8). Alternatively, when analytic computations are possible, stochastic Galerkin (LSPG($A$/SG)) may be significantly less expensive than other methods for a fixed stochastic subspace (Figure 5.9).

- Goal-oriented LSPG($F^T F$) can have uncontrolled errors in error measures that deviate from the output-oriented error measure $\eta_Q$ when the linear operator $F$ has more columns $n_x$ than rows $n_o$ (Figure 5.5). This is because the minimum singular value is zero in this case (i.e., $\sigma_{\min}(F) = 0$), which leads to unbounded
stability constants in other error measures (Table 5.2).

- Stochastic Galerkin often leads to divergence in different error measures (Figure 5.3). In this case, applying LSPG with the appropriate targeted error measure can significantly improve accuracy (Figure 5.4).

Future work includes developing efficient sparse solvers for the stochastic LSPG methods and extending the methods to parameterized nonlinear systems.
Chapter 6: Conclusion

In this thesis, we proposed solution algorithms for addressing two difficulties in using the stochastic Galerkin method for solving high-dimensional parameterized PDEs: (1) the solution of the Galerkin systems are computationally expensive and (2) the stochastic Galerkin method does not always guarantee optimality in the solution error. For efficient computations, we proposed the two-level low-rank iterative solver for linear elliptic parameterized PDEs and the low-rank variant of the Newton–Krylov method for nonlinear parameterized PDEs. For optimality, we proposed the stochastic least-squares Petrov–Galerkin method. We examined the efficiency and the optimality of the proposed methods on several benchmark problems.

In Chapter 3, we presented the two-level low-rank iterative solver for linear elliptic parameterized PDEs, which identifies an important low-dimensional subspace with a coarse-grid computation and uses the identified subspace for truncating all intermediate quantities generated during the low-rank GMRES iteration on the fine-grid space. In the low-rank GMRES method, computational efficiency was achieved by using the matrix operations, which exploits the Kronecker-product structure. Numerical experiments on two benchmark problems, a stochastic diffusion prob-
lem and a stochastic convection-diffusion problem, demonstrated that the two-level algorithm achieved significant savings in computational costs.

In Chapter 4, we presented a low-rank variant of the Newton–Krylov method for solving the Navier–Stokes equations with uncertain viscosity. We adapted the hybrid linearization scheme, which employs a few steps of Picard iterations followed by the Newton iterations, to the low-rank variant of the nonlinear iteration. To further achieve computational savings, we consider the inexact version of the nonlinear iteration, which approximately solves the linear system at each nonlinear step. We demonstrated the performance of the proposed method with the set of benchmark problems with various configurations characterizing the statistical features of the uncertain viscosity. The numerical experiments showed that the proposed method achieved significant computational savings for the problems with smaller CoV and larger correlation lengths.

In Chapter 5, we presented the stochastic least-squares Petrov–Galerkin method, which produces an optimal solution in a given finite-dimensional subspace minimizing the solution error in a target norm. We showed that specific choices of the weighting function lead to certain minimization formulations that are mathematically equivalent to the stochastic Galerkin method and the pseudo-spectral method. The method is monotonic convergent in the sense that the method produces monotonically decreasing solution error in a target norm. Using extensive numerical experiments on benchmark problems, we demonstrated that each LSPG method is optimal in minimizing its targeted error measure and is optimal also in terms of computational costs when an accurate solution in a target error measure is sought.
Bibliography


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