



Efficient iterative algorithms for the stochastic finite element method with application to acoustic scattering

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Abstract

In this study, we describe the algebraic computations required to implement the stochastic finite element method for solving problems in which uncertainty is restricted to right-hand side data coming from forcing functions or boundary conditions. We show that the solution can be represented in a compact outer product form which leads to efficiencies in both work and storage, and we demonstrate that block iterative methods for algebraic systems with multiple right-hand sides can be used to advantage to compute this solution. We also show how to generate a variety of statistical quantities from the computed solution. Finally, we examine the behavior of these statistical quantities in one setting derived from a model of acoustic scattering.

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1. Introduction

It is common practice for mathematical models to be studied under the assumption that data defining the models are precisely understood. In reality, however, this simplifying assumption is often not valid, and there is considerable uncertainty in specification of models. Sources of uncertainty include geological properties of transporting media, material properties of structures, and unknown aspects of boundary conditions.

One approach for addressing this issue is to treat poorly specified data as random variables having some given statistical properties such as means and higher order moments, and then to determine analogous statistical properties of solutions. For boundary value problems with uncertain data (*stochastic partial differential equations*), a methodology known as the *stochastic finite element method* has generated considerable attention in the last decade [7,8,12,13,16]. This approach starts with a boundary value problem in d -dimensional physical space. The stochastic component of the problem statement is then specified approximately using an m -dimensional auxiliary space which is derived from an underlying probability space associated with the data. The result is a $(d + m)$ -dimensional model, which can be stated in a weak form on a suitable function space using a combination of standard variational constructions for the physical component of the problem together with averaging for the stochastic component. We will outline the details of this methodology in Section 2.

Once this weak formulation is specified, a numerical solution of the stochastic partial differential equation can be computed in essentially the same manner as for deterministic problems. In particular, the introduction of finite dimensional subspaces leads to an algebraic system of equations whose solution can be used to approximate statistical properties of the physical solution, such as its mean, variance and covariances. Our concern in this paper is to explore the computational costs of solving the systems in question and of generating statistical analyses of the solution.

We will focus on problems where randomness only affects the right-hand sides of the algebraic systems, that is, where the forcing terms or boundary data are random functions. A natural example of this arises in models of acoustic or electromagnetic scattering, where lack of information about the material properties of scatterers or the shape and structure of boundaries such as ocean bottoms leads to uncertainty in boundary conditions. We will use this model, specifically, the numerical solution of the Helmholtz equation, as a benchmark problem, and in our assessment we will explore computational issues associated with quantities such as moments and probability distributions of acoustic pressures, and how these are affected by characteristics of the problem such as wave numbers.

One of the computational tasks required is the solution of algebraic systems of equations with multiple right-hand sides. In the case of uncertain boundary data, the costs of this component of the computation can be kept low using the fact that the solution has a Kronecker product structure. For our scattering example, the systems can be solved efficiently with a multigrid algorithm for the discrete Helmholtz equation [10], and we show that efficiency can be enhanced in some cases using block iterative methods for systems with multiple right-hand sides [5,11,18,23]. With this strategy for solving the algebraic systems, the dominant cost of the computation is that of computing statistical quantities. We also show that the Kronecker product structure of the solution allows storage costs to be kept relatively low, and moreover it enables the statistical computations to be performed using efficient matrix-oriented operations that are trivially parallelizable and amenable to implementation using Level 3 Basic Linear Algebra Subprograms (BLAS3) [9]. This means that it is possible to handle relatively fine “discretization” in the stochastic domain that would otherwise not be possible.

We note that an alternative approach for handling random right-hand sides has been developed in Schwab and Todor [21], where it is shown that the mean and second moment of the solution can be computed directly. The latter entails the solution of a fourth-order equation derived for this quantity. It is shown in [21] that if the underlying differential operator is coercive, then so is the associated fourth order

system, and efficient multilevel algorithms (but dependent on special sparse grids) can be developed to solve it. The approach under consideration here has the advantage that it readily yields more general statistical information such as higher order moments and probability distributions. It is also relatively straightforward to implement, essentially only requiring algorithm technology for second-order problems. In particular, if, as in the example considered here, the underlying problem is *not* coercive, it is still possible to take advantage of efficient algorithms for that problem.

A summary of the contents of the paper is as follows. Section 2 contains a description of the stochastic finite element methodology and identifies the structure of the algebraic systems derived from discretization. Section 3 describes the iterative algorithms that we consider for solving the discrete Helmholtz equation and the block versions designed to handle multiple right-hand sides, and then it presents some experimental results demonstrating the performance of these solvers. Section 4 then outlines the costs of computing statistical quantities associated with the solution and shows the results of these computations. Finally, Section 5 contains some concluding remarks.

2. The stochastic finite element method

We briefly describe the general methodology with an eye towards showing the structure of the algebraic systems. For our description we use the problem that we will study in experiments, the Helmholtz equation; it will be obvious that the approach is general. See [13] for complete descriptions of this methodology.

2.1. Introduction: weak formulation

A model of acoustic scattering from a bounded obstacle is given by the Helmholtz equation

$$\begin{aligned} -\Delta u - k^2 u &= f \quad \text{in } \mathcal{D} \\ B(u) &= g \quad \text{on } \Gamma \\ \frac{\partial u}{\partial n} &= L(u) \quad \text{on } \Gamma_\infty, \end{aligned} \tag{2.1}$$

where the solution domain $\mathcal{D} \subset \mathbb{R}^d$ is bounded internally by the obstacle boundary $\Gamma \subset \partial\mathcal{D}$ and externally by an artificial boundary Γ_∞ . The boundary differential operator B is such that Dirichlet, Neumann or Robin boundary conditions result along Γ , and L is the Dirichlet-to-Neumann operator [14] or a suitable approximation thereof.

The weak form of this problem is to find $u \in V_g$ such that

$$a(u, v) = \ell(v) \quad \forall v \in V, \tag{2.2}$$

where V and V_g denote the linear and affine subspaces of $H^1(\mathcal{D})$ of functions satisfying any homogeneous, respectively, inhomogeneous essential boundary conditions along Γ . In the simplest case of Dirichlet boundary data along all of Γ , the sesquilinear form $a : H^1(\mathcal{D}) \times H^1(\mathcal{D}) \rightarrow \mathbb{C}$ is

$$a(u, v) = \int_{\mathcal{D}} (\nabla u \cdot \nabla \bar{v} - k^2 u \bar{v}) \, dx - \int_{\Gamma_\infty} \bar{v} L(u) \, ds$$

and the right-hand side functional $\ell : H^1(\mathcal{D}) \rightarrow \mathbb{C}$ is

$$\ell(v) = \int_{\mathcal{D}} f \bar{v} \, dx.$$

To introduce randomness into this formulation, let (Ω, \mathcal{A}, P) denote a probability space with sample space Ω , σ -algebra \mathcal{A} and probability measure P . Let $\zeta : \Omega \rightarrow \mathbb{C}$ be a complex-valued random variable with $\zeta \in L^1(\Omega)$. The mean or expected value of ζ is

$$\langle \zeta \rangle = \int_{\Omega} \zeta(\omega) dP(\omega) = \int_{\mathbb{C}} z d\mu(z),$$

where μ is the distribution probability measure associated with ζ and defined on the Borel sets B in the complex plane by $\mu(B) = P(\zeta^{-1}(B))$. Given a bounded domain $\mathcal{D} \subset \mathbb{R}^d$ as above, a random function

$$u : \mathcal{D} \times \Omega \rightarrow \mathbb{C}, \quad (x, \omega) \mapsto u(x, \omega)$$

is one that is jointly measurable with respect to Lebesgue measure on \mathcal{D} and the measure P on Ω and for which

$$\langle \|u(\cdot, \omega)\|_{L^2(\mathcal{D})} \rangle < \infty.$$

The space of random functions is a Hilbert space $\tilde{L}^2(\mathcal{D} \times \Omega)$ with respect to the inner product

$$(u, v)_{\tilde{L}^2} = \langle (u(x, \cdot), v(x, \cdot))_{L^2(\mathcal{D})} \rangle.$$

The stochastic Sobolev spaces $\tilde{H}^k(\mathcal{D} \times \Omega)$ are defined analogously.

If any of the data in the Helmholtz equation (2.1) is random (e.g., the wave number k , forcing function f , or Dirichlet boundary data g), then the solution u will be a random function. The weak form of the stochastic problem is then to find $u \in \tilde{V}_g$ such that

$$\tilde{a}(u, v) = \tilde{\ell}(v) \quad \forall v \in \tilde{V}, \tag{2.3}$$

where $\tilde{a}(u, v) = \langle a(u, v) \rangle$, $\tilde{\ell}(v) = \langle \ell(v) \rangle$, and \tilde{V}_g and \tilde{V} are the stochastic Sobolev spaces analogous to V_g and V .

2.2. Series representation of randomness and derived weak form

We consider the development of the stochastic finite element method using series representations of random functions in which the deterministic and stochastic components are separated. For the sake of concreteness, we describe the methodology under the assumption that the forcing function f of (2.1) is random. Thus, suppose f is represented in series form as

$$f(x, \omega) = f_0(x) + \sum_{r=1}^{\infty} \sqrt{\lambda_r} f_r(x) \zeta_r(\omega), \tag{2.4}$$

where the equality is to be interpreted in the least squares sense. We will discuss situations where other components of the problem are random in Section 2.3.

An example of a series (2.4) is the Karhunen–Loève expansion [24, pp. 447ff], where the random variables $\{\zeta_r\}$ are uncorrelated and the orthogonal functions $\{f_r\}$ are the eigenfunctions of the covariance function

$$c(x, y) = \langle f(x)f(y) \rangle - \langle f(x) \rangle \langle f(y) \rangle$$

associated with f . That is, the integral equation

$$(\mathcal{C}\gamma)(x) = \lambda\gamma(x), \quad \text{where } (\mathcal{C}\gamma)(x) = \int_{\mathcal{D}} c(x, y)\gamma(y) dy, \tag{2.5}$$

is a linear integral eigenvalue problem in which, by definition, the kernel is symmetric and positive-semidefinite. From the general theory of integral equations [6, Chapter 3], \mathcal{C} is a compact operator and there exists a countable sequence of eigenpairs $\{(\lambda_r, f_r)\}$ where the eigenvalues $\{\lambda_r | \lambda_1 \geq \lambda_2 \geq \dots\}$ are nonnegative and

the eigenfunctions $\{f_r\}$ are orthogonal in $L^2(\mathcal{D})$. Alternative series representations could be obtained by other means, for example, through experiment or properties of the model; see [16] for some discussion of this point. It is also shown in [2] that if the covariance function is not known, it is possible to compute reasonable estimates for the required eigenvalues and eigenvectors from samples of the random function.

For computation, the infinite series (2.4) is approximated by a finite one with, say, m terms. Following [2,8], we will also assume that the random variables $\{\xi_r\}$ are *independent*. In general, the more localized the covariance kernel of f (the smaller the correlation length), the slower the decay of its eigenvalues and the more terms need be retained in the expansion to achieve good accuracy. Thus, the utility of this approach depends on the assumption that the properties of physical systems under consideration vary smoothly, i.e., there are significant correlations in the random data. In this case, it is expected that a truncated version of (2.4) with a small number m of terms in the sum is sufficient to capture the randomness in the system.

Assume now that the random function is given by a finite-term expansion

$$f(x, \xi) = f_0(x) + \sum_{r=1}^m \sqrt{\lambda_r} f_r(x) \xi_r, \tag{2.6}$$

where the random variables $\{\xi_r\}$ are pairwise mutually independent with mean 0 and variance 1. Let $\mathcal{I}_r = \xi_r(\Omega)$ denote the image of ξ_r , and let $\mathcal{I} = \mathcal{I}_1 \times \dots \times \mathcal{I}_m$. Collecting these variables into the random vector $\xi = (\xi_1, \dots, \xi_m)$, we have $\xi(\Omega) \subset \mathcal{I}$. Assume that ξ_r possesses the probability density function $\rho_r : \mathcal{I}_r \rightarrow \mathbb{R}$, which gives rise to the joint density function

$$\rho(\xi) = \rho_1(\xi_1) \rho_2(\xi_2) \dots \rho_m(\xi_m).$$

The stochastic variational formulation of the Helmholtz equation (2.1) uses as test functions random functions in the space

$$\widetilde{\mathcal{V}} = \left\{ u(x, \xi) : \int_{\mathcal{I}} (\|u\|_{H^1(\mathcal{D})}) \rho(\xi) \, d\xi < \infty \right\}, \tag{2.7}$$

satisfying homogeneous boundary conditions $u = 0$ on Γ , with trial functions in the space $\widetilde{\mathcal{V}}_g$ defined analogously. The stochastic variational problem is then specified as in (2.3) with

$$\begin{aligned} \tilde{a}(u, v) &= \int_{\mathcal{I}} \left(\int_{\mathcal{D}} (\nabla u \cdot \nabla v - k^2 uv) \, dx - \int_{\Gamma_\infty} \bar{v} L(u) \, ds \right) \rho(\xi) \, d\xi \\ \tilde{\ell}(v) &= \int_{\mathcal{I}} \left(\int_{\mathcal{D}} f \bar{v} \, dx \right) \rho(\xi) \, d\xi. \end{aligned} \tag{2.8}$$

The weak solution u can be viewed as defined on a $(d + m)$ -dimensional domain $\mathcal{D} \times \mathcal{I}$.

2.3. Discretization and the stochastic system

In order to establish notation, we briefly discuss the discretization of the deterministic problem (2.1), assuming Dirichlet boundary conditions $u = g$ hold on the obstacle boundary Γ . Let $\mathcal{V}^h = \text{span}\{\phi_1, \dots, \phi_{N_x}\}$ denote a finite dimensional subspace of $H_0^1(\mathcal{D})$, and let \mathcal{V}_g^h denote the affine space obtained by adding basis functions $\{\phi_{N_x+1}, \dots, \phi_{N_x+N_E}\}$ to handle degrees of freedom on the boundary. As is well known, the discrete weak formulation entails finding

$$u_h = \sum_{j=1}^{N_x} \mathbf{u}_j \phi_j + \sum_{j=N_x+1}^{N_x+N_E} g(x_j) \phi_j$$

such that

$$\sum_{j=1}^{N_x} a(\phi_j, \phi_i) \mathbf{u}_j = \int_{\mathcal{D}} f \phi_i \, dx - \sum_{j=N_x+1}^{N_x+N_E} a(\phi_j, \phi_i) g(x_j) \quad \forall i = 1, \dots, N_x.$$

This is a linear system of equations $\mathbf{A}\mathbf{u} = \mathbf{f}$ where

$$\mathbf{f} = [(f, \phi_i)]_{i=1}^{N_x} - A_{UE} \mathbf{g}; \tag{2.9}$$

A_{UE} represents the coupling between degrees of freedom constrained by Dirichlet boundary conditions and other unknowns, and $\mathbf{g} = [g(x_j)]_{j=N_x+1}^{N_x+N_E}$ is the vector of nodal boundary values.

Now consider the stochastic problem defined by (2.3) and (2.8). For the discretization, let

$$\widetilde{\mathcal{V}}^h = \text{span}\{\chi_{jq}(x, \xi) = \phi_j(x) \psi_q(\xi) : j = 1, \dots, N_x, q = 1, \dots, N_\xi\},$$

denote a finite-dimensional subspace of $\widetilde{\mathcal{V}}$ of (2.7), where $\{\psi_1, \dots, \psi_{N_\xi}\}$ is a basis for a finite-dimensional subspace of $L^2(\mathcal{I})$. Let $\widetilde{\mathcal{V}}_g^h$ denote the affine space satisfying inhomogeneous essential boundary conditions.

The discrete stochastic problem is then to find $u_h \in \widetilde{\mathcal{V}}_g^h$,

$$u_h(x, \xi) = \sum_{q=1}^{N_\xi} \sum_{j=1}^{N_x} \phi_j(x) \psi_q(\xi) \mathbf{u}_{jq} + \sum_{j=N_x+1}^{N_x+N_E} \phi_j(x) g(x_j) \tag{2.10}$$

such that

$$\langle a(u_h, v_h) \rangle = \langle \ell(v_h) \rangle \quad \forall v_h \in \widetilde{\mathcal{V}}^h.$$

The result is a linear system of equations, the *stochastic system*

$$\mathbf{A}\mathbf{u} = \mathbf{f} \tag{2.11}$$

of order $N_x \times N_\xi$, for unknowns

$$(\mathbf{u}_{11}, \mathbf{u}_{21}, \dots, \mathbf{u}_{N_x-1, N_\xi}, \mathbf{u}_{N_x, N_\xi})^T.$$

Once \mathbf{u} is obtained, statistical properties of the associated random function u_h can be computed easily, see Section 4.

As we have noted, this study concerns the case where randomness only affects the right-hand side of the algebraic systems generated, i.e., where the source term or boundary data is random. Let us consider the structure of the discrete problem (2.11) in this case. The entries of the finite element system matrix \mathbf{A} are

$$\begin{aligned} \langle a(\chi_{jq}, \chi_{ip}) \rangle &= \int_{\mathcal{I}} a(\phi_j \psi_q, \phi_i \psi_p) \rho(\xi) \, d\xi \\ &= \left(\int_{\Gamma} \psi_q \psi_p \rho(\xi) \, d\xi \right) \left(\int_{\mathcal{D}} \nabla \phi_j \cdot \nabla \phi_i - k^2 \phi_j \phi_i \, dx - \int_{\Gamma_\infty} \phi_i L \phi_j \, ds \right) \\ &= \langle \psi_q \psi_p \rangle a(\phi_j, \phi_i), \end{aligned}$$

for $1 \leq i, j \leq N_x, 1 \leq p, q \leq N_\xi$. Denoting by $G \in \mathbb{R}^{N_\xi \times N_\xi}$ the Grammian matrix

$$[G]_{pq} = \langle \psi_q \psi_p \rangle, \quad p, q = 1, \dots, N_\xi, \tag{2.12}$$

and by $A \in \mathbb{C}^{N_x \times N_x}$ the stiffness matrix of the deterministic equation, the coefficient matrix is seen to have the Kronecker structure

$$\mathbf{A} = G \otimes A.$$

Note that this implicitly determines an ordering for the rows and columns of A . The rows are ordered so that for each p , indices $i = 1, \dots, N_x$ are grouped together, and then p is ordered from 1 to N_ξ ; the same grouping applies to the columns.

For the right-hand side, assume as in Section 2.2 that the forcing function is random, and also assume for the moment that homogeneous Dirichlet boundary conditions $g = 0$ hold on Γ . It then follows from (2.6) and (2.8) that the entry of \mathbf{f} corresponding to the test function $\chi_{ip} = \phi_i \psi_p$ is

$$\langle \ell(\chi_{ip}) \rangle = \int_{\mathcal{J}} \ell(f, \chi_{ip}) \rho(\xi) d\xi = \ell(f_0, \phi_i) \langle \psi_p \rangle + \sum_{r=1}^m \sqrt{\lambda_r} \ell(f_r, \phi_i) \langle \xi_r \psi_p \rangle. \tag{2.13}$$

Let us define the vectors

$$\begin{aligned} \mathbf{f}_r &= [(f_r, \phi_i)]_{i=1}^{N_x}, \quad r = 0, 1, \dots, m \\ \boldsymbol{\psi}_0 &= [\langle \psi_p \rangle]_{p=1}^{N_\xi}, \\ \boldsymbol{\psi}_r &= [\langle \xi_r \psi_p \rangle]_{p=1}^{N_\xi}, \quad r = 1, \dots, m, \end{aligned} \tag{2.14}$$

where upon the discrete system has the form

$$(G \otimes A) \mathbf{u} = \mathbf{f}, \quad \mathbf{f} = \boldsymbol{\psi}_0 \otimes \mathbf{f}_0 + \sum_{r=1}^m \sqrt{\lambda_r} (\boldsymbol{\psi}_r \otimes \mathbf{f}_r).$$

That is, the right-hand side lies in an $(m + 1)$ -dimensional subspace of $\mathbb{R}^{N_\xi \times N_x}$. The solution is then

$$\mathbf{u} = (G \otimes A)^{-1} \mathbf{f} = (G^{-1} \otimes A^{-1}) \mathbf{f} = (G^{-1} \boldsymbol{\psi}_0) \otimes (A^{-1} \mathbf{f}_0) + \sum_{r=1}^m \sqrt{\lambda_r} (G^{-1} \boldsymbol{\psi}_r) \otimes (A^{-1} \mathbf{f}_r). \tag{2.15}$$

This entails the solution of $m + 1$ systems of size N_ξ with coefficient matrix G , and $m + 1$ systems of size N_x with coefficient matrix A . In practice, the basis functions $\{\psi_p\}$ for the stochastic component are often chosen to be orthogonal with respect to the probability measure [8,13], in which case G is a diagonal matrix. Thus, the main computational requirement is solution of the $m + 1$ systems with coefficient matrix A .

Although the derivation above is for the case of stochastic forcing function and homogeneous boundary conditions, the conclusion reached is general. For example, if a nonzero Dirichlet condition holds on Γ , then the construction is identical except \mathbf{f}_0 has the form (cf. (2.9))

$$\mathbf{f}_0 = [(f_0, \phi_i)]_{i=1}^{N_x} - A_{UE} \mathbf{g}.$$

More generally, if it is Dirichlet boundary conditions that are random (we will explore this in experiments described in Section 3), then terms of the form

$$\boldsymbol{\psi}_0 \otimes (A_{UE} \mathbf{g}_0) + \sum_r \sqrt{\lambda_r} (\boldsymbol{\psi}_r \otimes (A_{UE} \mathbf{g}_r))$$

will be included in the right-hand side. Similar considerations apply for Neumann conditions on the obstacle boundary.

2.4. Implementation

The notation used in the previous section treats the unknowns \mathbf{u} of (2.11) as a vector. In an implementation, it is in fact more convenient to treat the solution as a two-dimensional array. In particular, consider the matrices

$$F = [\mathbf{f}_0, \mathbf{f}_1, \dots, \mathbf{f}_m], \quad A = \text{diag}(1, \sqrt{\lambda_1}, \dots, \sqrt{\lambda_m}), \quad \Psi = [\boldsymbol{\psi}_0, \boldsymbol{\psi}_1, \dots, \boldsymbol{\psi}_m],$$

where the vectors $\{\mathbf{f}_r\}$ and $\{\boldsymbol{\psi}_r\}$ are defined in (2.14). Then the system (2.11) is essentially of the form

$$AU = B, \quad (2.16)$$

where $B = FW^T$ with $W = G^{-1}(\Psi A)$. The solution can then be represented implicitly in outer-product form as

$$U = VW^T, \quad (2.17)$$

where $V = A^{-1}F$ is obtained by solving the system of equations $AV = F$ with $m + 1$ right-hand sides.

3. Iterative solution of the stochastic system

For the problem under consideration, the coefficient matrix of (2.16) is a discrete Helmholtz operator, which is complex, symmetric and indefinite. In this section, we describe an iterative algorithm that can be used to solve this system and demonstrate its effectiveness on a set of benchmark problems.

3.1. Solution algorithm

The basic solution algorithm we use is a multigrid method designed for the Helmholtz equation, adapted to handle multiple right-hand sides. As is well known, the principle behind multigrid is to combine *smoothers* to eliminate oscillatory components of the error on fine grids, together with *coarse grid corrections* to eliminate smooth components. For the Helmholtz equation, standard multigrid approaches are not effective. There are two difficulties:

1. Standard smoothers such as the Jacobi and Gauss–Seidel methods do not work because certain smooth modes are amplified by these operations.
2. The eigenvalues associated with some smooth modes change signs during the grid coarsening process, which causes the coarse grid correction to also amplify some smooth modes rather than eliminate them from the error.

These difficulties derive from the indefiniteness of the system. In [10], we developed a method that addresses them. The first difficulty is handled by replacing standard smoothers with Krylov subspace methods, i.e., GMRES iteration [20] is used as the smoother. The second one is handled by using the multigrid operation as a preconditioner for an outer Krylov subspace iteration, so that components of the error not treated correctly by the multigrid coarse grid computations are eliminated. Because the multigrid smoother is no longer a linear operator, the outer iteration must handle this via a so-called “flexible” GMRES algorithm [19]. A complete description and analysis of the preconditioning strategy is given in [10], where it is demonstrated that the algorithm exhibits “textbook multigrid” convergence behavior, that is, convergence rates that are independent of the discretization parameter; there is some dependence on the wave number k .

We also adapt this approach to handle the system (2.16) with multiple right-hand sides, the number of which is denoted by m within this section. Recall that Krylov subspace methods generate an iterate at step s using a certain subspace of dimension s . Two types of Krylov subspace algorithms have been proposed for problems with multiple right-hand sides:

- *Block algorithms* [3,18] construct a subspace of dimension ms formed by the union of the s -dimensional subspaces for each right-hand side. Then, for each right-hand side, they find the best solution within that subspace. Deflation is used to remove vectors that become linearly dependent.

- *Seed algorithms* [4,22] form a Krylov subspace using one of the right-hand sides and then find the best solution for each of the m problems within that subspace. If the seed problem converges before the others, then a different right-hand side is chosen as the seed and the algorithm is repeated.

Each of these approaches has its advantages. Seed methods tend to perform best when the right-hand sides are related to each other, for example, if they arise from functions evaluated at nearby points [4]. This approach requires less storage than block methods: for systems of order N , the seed GMRES method requires storage proportional to sN , compared to smN for a block GMRES solver. On the other hand, block algorithms tend to converge more rapidly for more general right-hand sides, or when a small number of eigenvalues are well-separated from the others [18]. The block algorithm also makes much better use of computer memory traffic, since each access to the coefficient matrix is used for m matrix–vector products.

In our application, the right-hand side vectors (columns of F in (2.16)) derive from the orthogonal eigenvectors of the covariance matrix, and we found the seed method to be ineffective. Therefore, we restrict our attention to a block method. The idea of block iterative methods is due to O’Leary and Underwood. The block biconjugate gradient algorithm was described in [18], and a block quasi-minimum residual method in [11]. Algorithms for altering the block size adaptively were given in [1]. A block GMRES algorithm was presented by Vital [23].

We also need to modify the algorithm to handle the nonlinear preconditioner, as described in [5]. To present this *block flexible GMRES method* for (2.16), we use the generic notation $Ax = b$ for the linear system, and $w = M(v)$ to represent a generic *preconditioning* operation. This may be a linear operation derived from a matrix, or (as in the present setting) a nonlinear operation. Let $x_j^{(s)}$ denote an approximation of the solution to the j th equation (of m) computed at iteration s . The block flexible GMRES algorithm generates a sequence of matrices $\{V_j\}$ of dimensions $N \times m$ that together form a matrix $V = [V_1, \dots, V_s]$, and a set of matrices $Z_j = M(V_j)$, and $Z = [Z_1, \dots, Z_s]$. The block-Hessenberg matrix H has block entries H_{ij} . Each component $\{x_j^{(s)} : j = 1, \dots, m\}$ of the iterative solution has the form

$$x_j^{(s)} = x_j^{(0)} + Z_j y_j^{(s)} \tag{3.1}$$

such that the norm of the residual $\|b_j - Ax_j^{(s)}\|$ is minimal, where $x_j^{(0)}$ is a possibly arbitrary initial value. If $M(v)$ is a linear operation (say Mv where M is a matrix), then \mathcal{V}_s , the span of the columns of V , is the same as

$$\text{span} \{Mr_1, (MA)Mr_1, \dots, (MA)^{s-1}Mr_1, \dots, Mr_m, (MA)Mr_m, \dots, (MA)^{s-1}Mr_m\},$$

the span of the Krylov vectors. In this case, it is not necessary to store the auxiliary matrix Z . If M is a nonlinear operator, then \mathcal{V}_s will not be a Krylov subspace, but the s th iterate is still optimal among all possibilities of the form (3.1), which corresponds to an affine subspace of \mathbb{C}_N of dimension ms . The algorithm is as follows:

```

Compute the residual  $r = b - Ax$  of dimension  $N \times m$ .
until ( $\|r_\ell\| \leq \delta$ ,  $\ell = 1, \dots, m$ ),
% Generate a subspace of dimension  $ms$  from the residual  $r$ .
Define  $V_1$  to be the orthogonal factor in the  $QR$  factorization of  $r$ .
for  $i = 2, \dots, s + 1$ ,
    % Generate the directions defining the new basis vectors.
     $Z_{i-1} = M(V_{i-1})$ 
     $W = AZ_{i-1}$ 
    % Orthogonalize these directions against the previous ones.
    for  $j = 1, \dots, i-1$ ,
```

```


$$H_{j,i-1} = V_j^* W$$


$$W = W - V_j H_{j,i-1}$$

end for j.
Perform a QR factorization of W, obtaining the upper triangular factor  $H_{i,i-1}$  and the orthogonal matrix  $V_i$ .
end for i.
% Update each of the solutions.
for j = 1, ..., m,
     $c = V_j^* r_j$ 
    Solve the least squares problem  $\min_y \|c - Hy\|$ .
     $x_j = x_j + Zy$ 
end for j.
 $r = b - Ax$ 
end until

```

The loop on i can break down if the matrix W becomes rank-deficient. In this case, we reduce the size of the block by dropping the dependent columns, updating the solutions, and continuing with the residuals that have not converged.

3.2. Experimental results

We tested the performance of the block flexible GMRES algorithm for solving the stochastic Helmholtz equation on the two-dimensional domain \mathcal{D} consisting of the complement within the unit circle of a scatterer taken to be a semi-open cavity. Dirichlet-to-Neumann conditions are specified on the external boundary Γ_∞ . The discretization in space consists of piecewise linear elements on triangles. Fig. 1 shows the scatterer and the initial mesh used in all tests. For each wave number k , this mesh is refined until kh_{\max} was on the order of $\pi/5 \approx .63$, so that there are approximately ten points per wavelength. All computations were done using MATLAB. Mesh construction was done using the MATLAB PDE TOOLBOX routines `initmesh` and `refinemesh`, which performs a uniform mesh refinement.

All uncertainty in the specification of the boundary value problem occurs in the statement of boundary conditions on Γ , the boundary of the scatterer, where Dirichlet boundary conditions $u = g$ are such that g is

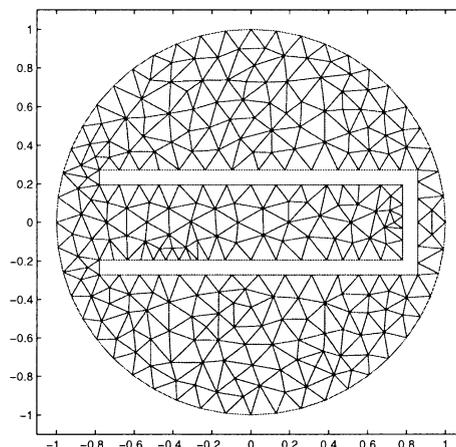


Fig. 1. Spatial domain and initial mesh used in spatial discretization.

a random function as in (2.6), with mean determined by an incident plane wave $g(x_1, x_2) = -e^{ik(x_1 \cos \theta + x_2 \sin \theta)}$ at angle $\theta = \pi/4$. We assume as in [7,8] that the random variables $\{\xi_r\}$ making up the series expansion of the random function are independent and uniformly distributed on an interval $\mathcal{I}_r = [-\alpha, \alpha]$, giving rise to the joint uniform distribution on $\mathcal{I} = [-\alpha, \alpha]^m$ with joint density function

$$\rho(\xi) = \left(\frac{1}{2\alpha}\right)^m. \tag{3.2}$$

The convention that $\langle \xi_i \xi_j \rangle = \delta_{ij}$ leads to the condition $\alpha = \sqrt{3}$. The stochastic domain $\mathcal{I} = [-\alpha, \alpha]^m$ is discretized using a uniform mesh: each of the m coordinate intervals is subdivided into n_ξ equal subintervals, resulting in $N_\xi = n_\xi^m$ elements, each of which is an m -dimensional cube with side length $h_\xi = 2\alpha/n_\xi$. Since there are no continuity requirements on the probability space, the basis functions are taken to be piecewise constants; that is, the basis function ψ_q has the value one on the cube with index q and zero elsewhere. This leads to the particularly simple diagonal structure for the Grammian matrix G of (2.12), $G = \frac{1}{N_\xi} I$. (Recall that this is always the case if the basis functions for the stochastic discretization are orthogonal with respect to the probability measure.)

To define the series expansion (2.6), we assume that the covariance function associated with g has the form

$$c(x_1, x_2) = e^{-[|(x_1)_1 - (x_2)_1|/2 + |(x_1)_2 - (x_2)_2|]}, \quad x_1, x_2 \in \Gamma.$$

In the integral eigenvalue problem (2.5), the eigenfunctions become more oscillatory as their indices increase, and they can be interpreted as simulating the results of a series of independent experiments on different scales. A general requirement of this methodology is that the first m eigenfunctions and eigenvalues of the covariance operator, or discrete approximations to them, be available. In some circumstances, these may be obtained in closed form [13, pp. 27ff], or, alternatively, they may be approximated using a Galerkin discretization of the integral equation (2.5). For the problems we are considering, the domain of g is one-dimensional, and the Galerkin computation is inexpensive. The discrete eigenvalues and eigenvectors are computed directly from the Galerkin approximation. If the domain in question is of higher dimension, this computation can be done efficiently using sparse eigenvalue methods and fast summation techniques [15].

Table 1 examines the performance of the block flexible GMRES algorithm and compares it with that of the flexible GMRES algorithm (FGMRES) applied to each right-hand side separately. In these tests, the stopping criterion for the solvers was for the Euclidean norm of each component of the residual to satisfy

$$\|r_j\|/\|b_j\| < 10^{-6}, \quad j = 1, \dots, m + 1.$$

For the block method, the iteration was stopped when the maximal individual residual component meets this criterion. The table shows the total number of matrix–vector products performed during the course of the computation, and in parentheses, the number of iterations required for convergence. For FGMRES, the latter number is the average for $m + 1$ right-hand sides; for block FGMRES, it is the number of block iterations. Note that the dimension of the spaces constructed by the block FGMRES method depends on m , the number of right-hand sides, but not on the discretization parameter n_ξ associated with the stochastic domain, since F does not depend on n_ξ in (2.16).

It is clear from these results that the block methods require fewer matrix–vector products in all cases, and the difference in the number of these products becomes more dramatic as the number of right-hand sides increases and also as the wave number k grows, i.e., as the problem becomes more difficult. The results provide further evidence of the mesh independent performance of the multigrid algorithm. We note, however, that as the number of steps s increases, the advantages of the block FGMRES method become less pronounced, since the overhead in generating the Krylov space grow like $m^2 s^2 N_x$, compared to $m s^2 N_x$ when the right-hand sides are processed separately. A block Krylov subspace method such as QMR [11] would not suffer from this drawback, although it is not clear that this approach can be adapted to handle

Table 1

Number of matrix–vector products required to solve $m + 1$ systems arising from m -term series expansion, using preconditioned FGMRES

| | $m = 4$ | $m = 6$ | $m = 8$ |
|--------------------------|------------|------------|------------|
| $k = 5\pi$ | | | |
| $kh = .72, N_x = 4170$ | | | |
| Block FGMRES | 35 (7) | 49 (7) | 54 (6) |
| FGMRES | 37 (7.4) | 52 (7.4) | 67 (7.4) |
| $kh = .36, N_x = 16,196$ | | | |
| Block FGMRES | 40 (8) | 56 (8) | 72 (8) |
| FGMRES | 45 (9.0) | 63 (9.0) | 81 (9.0) |
| $k = 10\pi$ | | | |
| $kh = .72, N_x = 16,196$ | | | |
| Block FGMRES | 85 (17) | 105 (15) | 135 (15) |
| FGMRES | 153 (30.6) | 214 (30.6) | 276 (30.7) |
| $kh = .36, N_x = 63,816$ | | | |
| Block FGMRES | 90 (18) | 119 (17) | 162 (18) |
| FGMRES | 157 (31.4) | 220 (31.4) | 282 (31.3) |
| $k = 20\pi$ | | | |
| $kh = .72, N_x = 63,816$ | | | |
| Block FGMRES | 200 (40) | 245 (35) | 288 (32) |
| FGMRES | 360 (72.0) | 495 (70.7) | 636 (70.7) |

Numbers in parentheses are average iteration counts or number of block iterations.

a nonlinear preconditioner. Because solution of the linear systems step represents a low order cost for the complete construction of statistical data (see the next section), we have not explored this issue.

4. Computation of statistical data

Once the random function u_h of (2.10) is available, we are interested in statistical properties such as moments and probability distributions associated with it. In the case of time-harmonic wave propagation, an important quantity is the modulus $|u_h|$, which indicates the significance of the component with wave number k in the wave field. In this section, we describe the computations required to generate statistical data associated with the random function $|u_h(x, \xi)|$.

4.1. Computation of a distribution function

Consider the construction of the probability distribution function for the maximum modulus

$$F(a) \equiv \Pr(\max_x |u_h(x, \xi)| \leq a). \quad (4.1)$$

Let

$$\mathcal{S}_a = \{\xi \in \mathcal{S} : \max_x |u_h(x, \xi)| \leq a\}.$$

Using the definition of the joint density function (3.2), we have

$$F(a) = \int_{\mathcal{S}_a} \rho(\xi) d\xi = |\mathcal{S}_a| \left(\frac{1}{2\alpha}\right)^m.$$

To determine the volume of \mathcal{S}_a , let $\xi \in \mathcal{I}$ be given, and let $q = q(\xi)$ be the index of the stochastic element $\mathcal{I}_q \subset \mathcal{I}$ containing ξ . Then

$$\max_x |u_h(x, \xi)| = \max_x \left| \sum_{j=1}^{N_\xi} \mathbf{u}_{jq} \phi_j(x) \right| = \max_j |\mathbf{u}_{jq}|,$$

where the latter equality follows from the linearity of u_h in space. Letting

$$s_a = \left| \left\{ q : \max_j |\mathbf{u}_{jq}| \leq a \right\} \right|$$

it follows that $|\mathcal{S}_a| = \frac{(2\alpha)^m}{N_\xi} s_a$, and therefore $F(a) = s_a/N_\xi$. This construction requires $\max_j |\mathbf{u}_{jq}|$ for each q . Once these maxima are computed, they can then be used to compute $F(a)$ for any a .

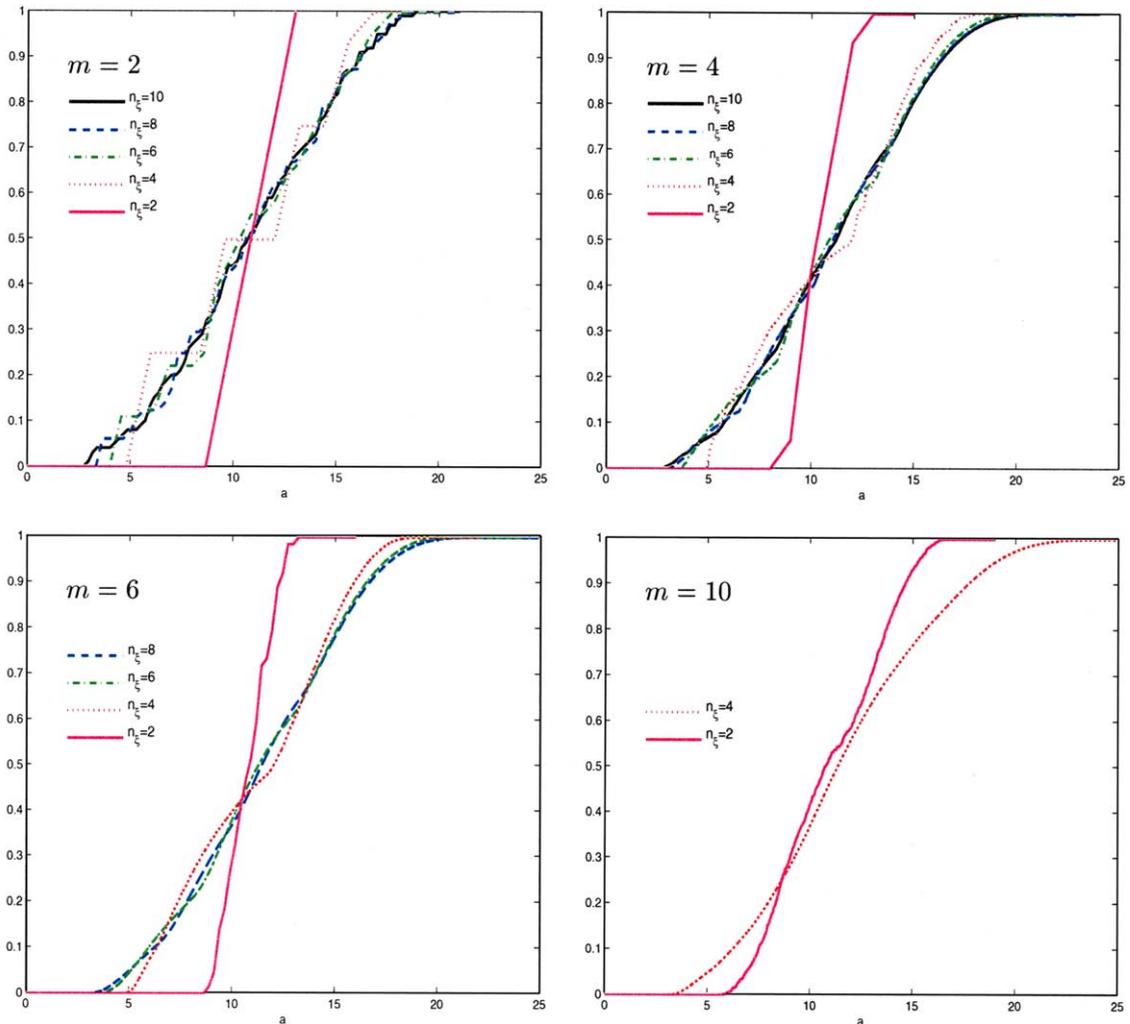


Fig. 2. Effect of the stochastic discretization parameter n_ξ on the estimate of $Pr(\max_x |u_h(x, \xi)| \leq a)$, for $k = 5\pi$ and various choices of m , the number of terms in the expansion (2.6).

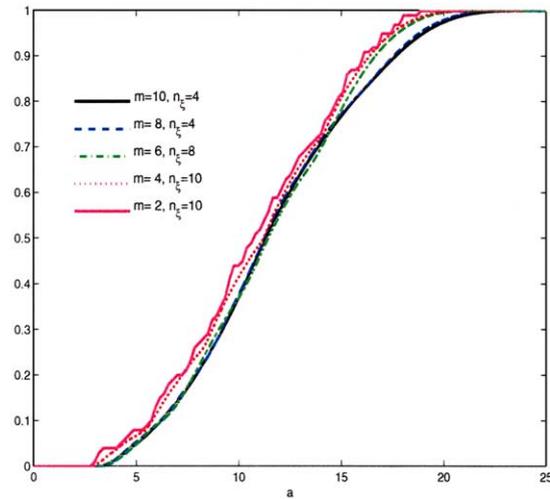


Fig. 3. Effect of the number of terms m in the expansion (2.6) on the estimate of $\Pr(\max_x |u_h(x, \xi)| \leq a)$, for $k = 5\pi$.

We show in Figs. 2 and 3 the results of computing the distribution function (4.1) for various parameter values. There is no analytic expression for this quantity, so it is difficult to make a rigorous assessment of the accuracy of the computations. Nevertheless, it is possible to identify certain qualitative aspects of the results as well as to place the costs of producing them in context. First, note that discretization of the random component of the problem can be viewed from two perspectives, derived from the number of terms m used in the finite expansion (2.6), and, once m is fixed, from the value n_ξ of the discretization parameter in \mathcal{S} . Convergence of the series expansion depends on the correlations within the process; when the finite expansion is fixed, it is shown in [8] that the error in the stochastic discretization (assuming an accurate spatial discretization) is proportional to n_ξ^{-1} . Since the number of stochastic degrees of freedom is proportional to $N_\xi = n_\xi^m$, it would be desirable to keep m as small as possible.

In Fig. 2, we consider the impact of the two parameters m and n_ξ , for a fixed wave number $k = 5\pi$. (The spatial discretization was such that $kh = .36$.) Each subplot in this figure corresponds to a fixed value of m for which n_ξ is allowed to vary. Each plot shows convergence to a fixed curve with refinement in n_ξ , as expected. It is also noteworthy that as m is increased, the quality of the solution for fixed n_ξ appears to improve. (For example, the solution for $n_\xi = 4$ is closer to the limiting value for each successive choice $m = 2, 4, 6$.) This indicates that the constants associated with the error bounds are smaller as m increases. Fig. 3 explores the impact of m more closely. For this example, the results suggest that $m = 8$ is an appropriate limiting value for the number of terms in the expansion. With $n_\xi = 4$, this yields 65,536 spatial degrees of freedom.⁴ The combination of smaller values of m together with large n_ξ (e.g., $m = 2$ and $n_\xi = 10$, yielding 1024 stochastic degrees of freedom) is able convey the qualitative structure of the distribution at significantly smaller cost.

Finally, Fig. 4 shows the estimated distribution function (4.1) for different values of k . These results suggest that this probability distribution function does not vary dramatically as the wave number increases.

⁴ We also remark that for $m = 10$, $n_\xi = 4$ was the largest discretization parameter we could use in our MATLAB environment. This led to $N_\xi = 1,048,576$ stochastic degrees of freedom.

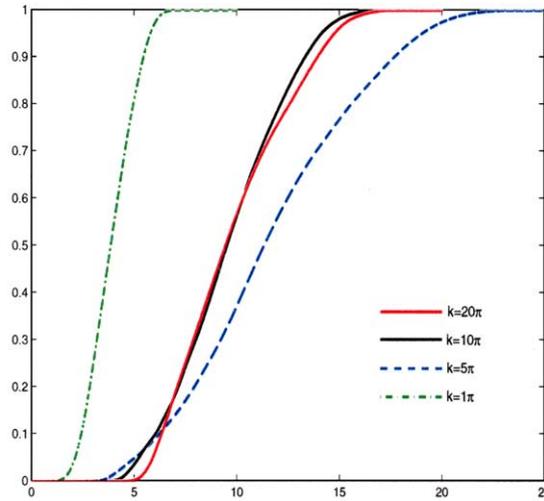


Fig. 4. Estimated probability distribution function $\Pr(\max_x |u_h(x, \xi)| \leq a)$, for various values of the wave number k .

4.2. Computation of higher moments

For examples of other statistical data to be computed, consider the moments of the modulus of u_h . Let $a_h(x, \xi) \equiv |u_h(x, \xi)|$, and let

$$a_h^{(v)}(x) \equiv \langle a_h(x, \cdot)^v \rangle, \quad v = 1, 2, 3, \dots$$

denote the moments of a_h . We have

$$\begin{aligned} a_h^{(v)}(x) &= \int_{\mathcal{J}} |u_h(x, \xi)|^v \rho(\xi) \, d\xi \\ &= \sum_{p=1}^{N_\xi} \int_{\mathcal{J}_p} \left| \sum_{q=1}^{N_\xi} \left(\sum_{j=1}^{N_x} \mathbf{u}_{jq} \phi_j(x) \right) \psi_q(\xi) \right|^v \rho(\xi) \, d\xi \\ &= \sum_{p=1}^{N_\xi} \left(\int_{\mathcal{J}_p} \rho(\xi) \, d\xi \right) \left| \sum_{j=1}^{N_x} \mathbf{u}_{jp} \phi_j(x) \right|^v \\ &= \frac{1}{N_\xi} \sum_{p=1}^{N_\xi} \left| \sum_{j=1}^{N_x} \mathbf{u}_{jp} \phi_j(x) \right|^v. \end{aligned}$$

This is straightforward to evaluate once the coefficients $\{\mathbf{u}_{jp}\}$ are available. In particular, the nodal values are

$$\mathbf{a}_i^{(v)} \equiv a_h^{(v)}(x_i) = \frac{1}{N_\xi} \sum_{p=1}^{N_\xi} |\mathbf{u}_{ip}|^v,$$

giving the piecewise linear interpolant of the v 'th moment,

$$\hat{a}_h^{(v)}(x) = \sum_{j=1}^{N_x} \mathbf{a}_j^{(v)} \phi_j(x).$$

The computations required for central moments $\langle (a_h - a_h^{(1)})^v \rangle$ are identical in structure.

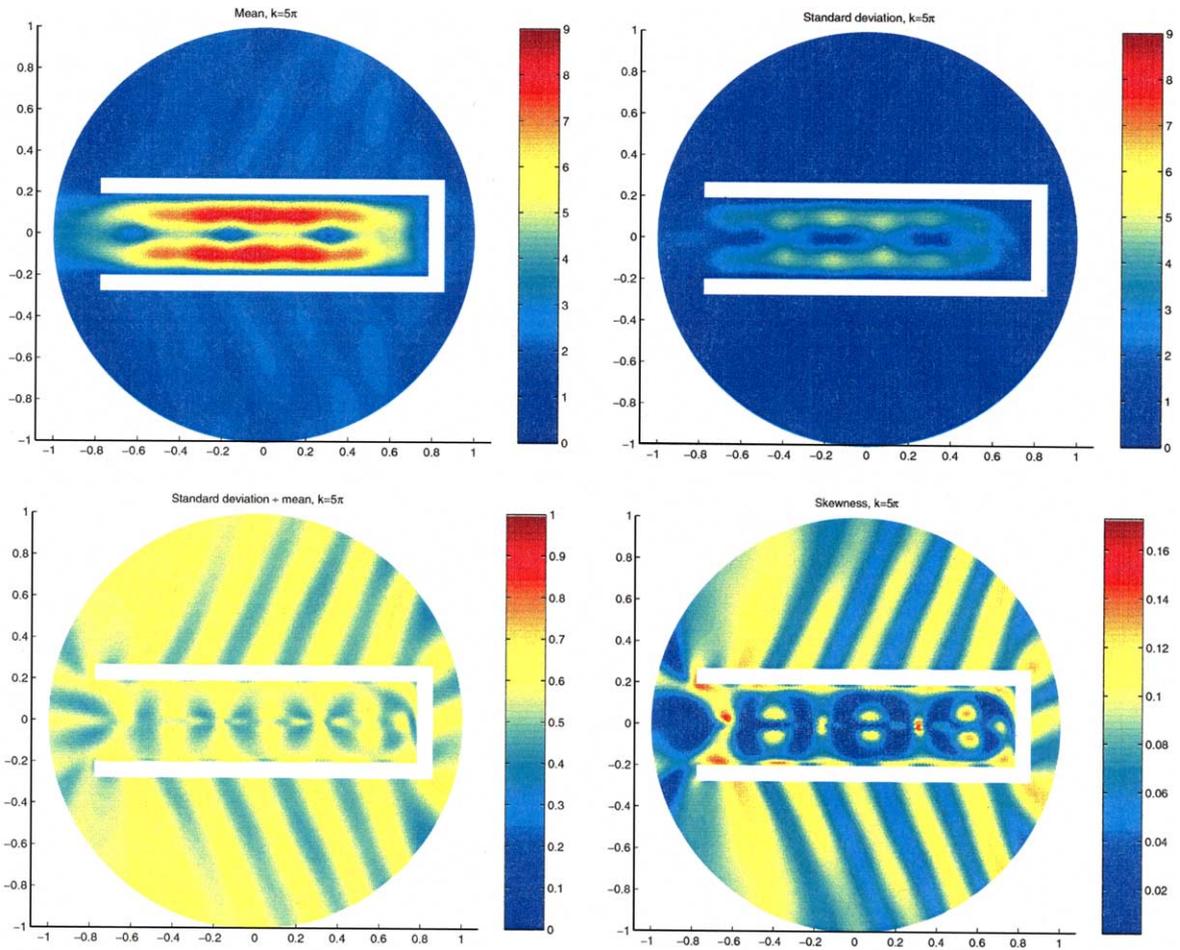


Fig. 5. Estimated mean μ_h , standard deviation σ_h , ratio σ_h/μ_h , and skewness of $|u_h|$, for $k = 5\pi$.

We examine some of these quantities in Figs. 5–7. Four things are shown: the mean μ_h , standard deviation σ_h , ratio of standard deviation to mean, and scaled third central moment (the *coefficient of skewness* [17])

$$\frac{1}{\sigma_h^3} \left\langle \left(a_h - a_h^{(1)} \right)^3 \right\rangle$$

of a_h . The data used for these plots come from the parameter choices $m = 8$ for the truncated expansion (2.6), stochastic discretization parameter $n_\xi = 4$ and spatial discretization satisfying $kh = .36$ for both $k = 5\pi$ and 10π and $kh = .72$ for $k = 20\pi$. Within each figure, the means and standard deviations are displayed using the same scalings. The magnitudes of the standard deviations largely mirror those of the means, and there is virtually no difference in the relative sizes of these quantities for different wave numbers. This indicates that size of the wave number k will not have a significant impact on the confidence that can be attributed to computed mean solutions. The depictions of skewness indicate that near the corner singu-

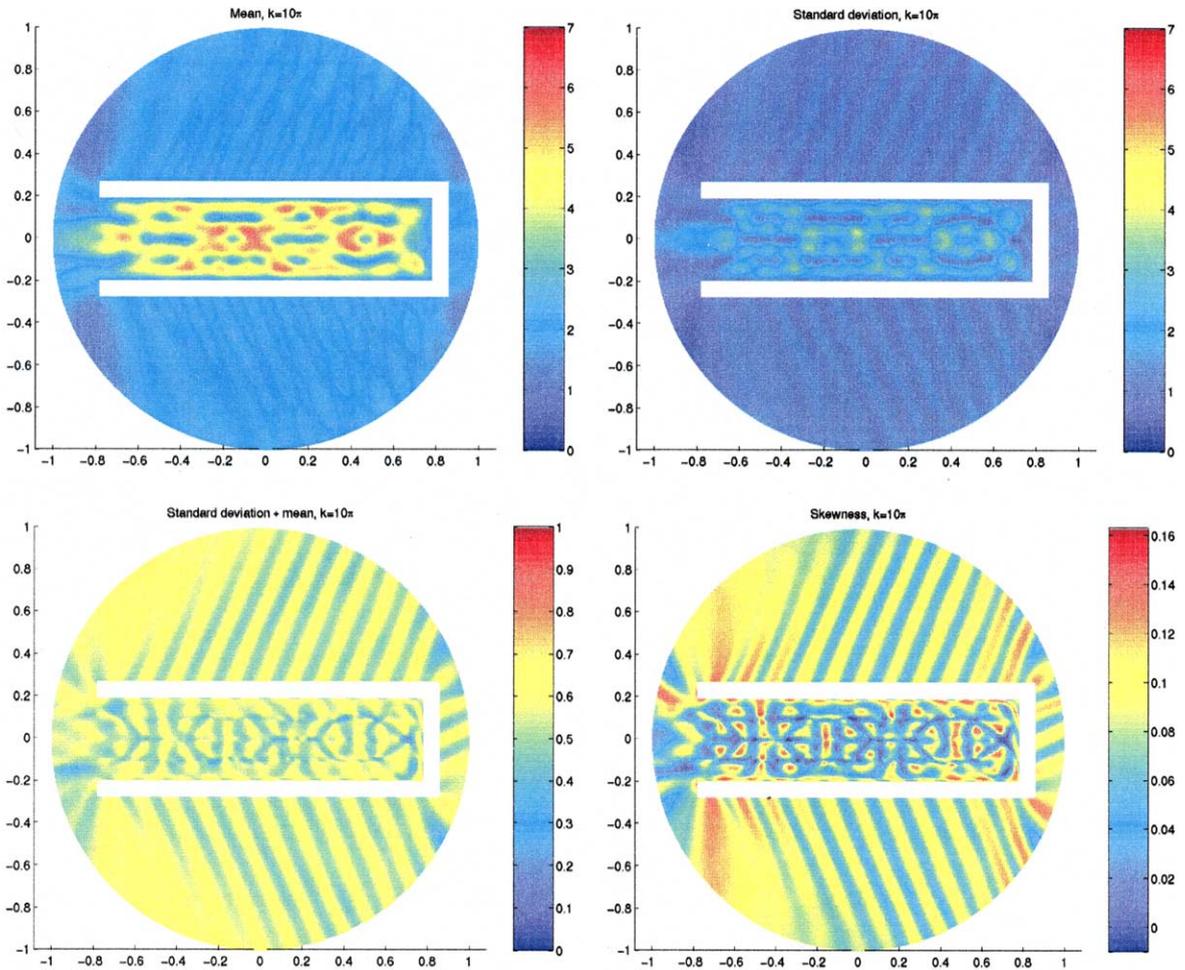


Fig. 6. Estimated mean μ_h , standard deviation σ_h , ratio σ_h/μ_h , and skewness of $|u_h|$, for $k = 10\pi$.

larities, the distributions tend to be more skewed toward the right (positive direction) with respect to the mean, and inside the cavity they are skewed more toward the left; this may be of use in identifying the shape of scatterers.

Note that all these computations require the complete set of values $\{\mathbf{u}_{jq} : j = 1, \dots, N_x, q = 1, \dots, N_\xi\}$, which are obtained from (2.17) as

$$\mathbf{u}_{jq} = \sum_{r=0}^m v_{jr} \mathbf{w}_{qr}.$$

Consequently, the cost is of order $O(N_x N_\xi)$ and these computations represent the dominant expense of the process. The storage costs are also of this magnitude but can be reduced to order $m \max(N_x, N_\xi)$ by taking advantage of the outer-product representation (2.17) and recomputing \mathbf{u}_{jq} whenever it is needed. The trade-off here is a (small) additional computational expense of magnitude $O(m N_x N_\xi)$. This makes it feasible to handle large values of m or n_ξ that storage restrictions would otherwise prevent.

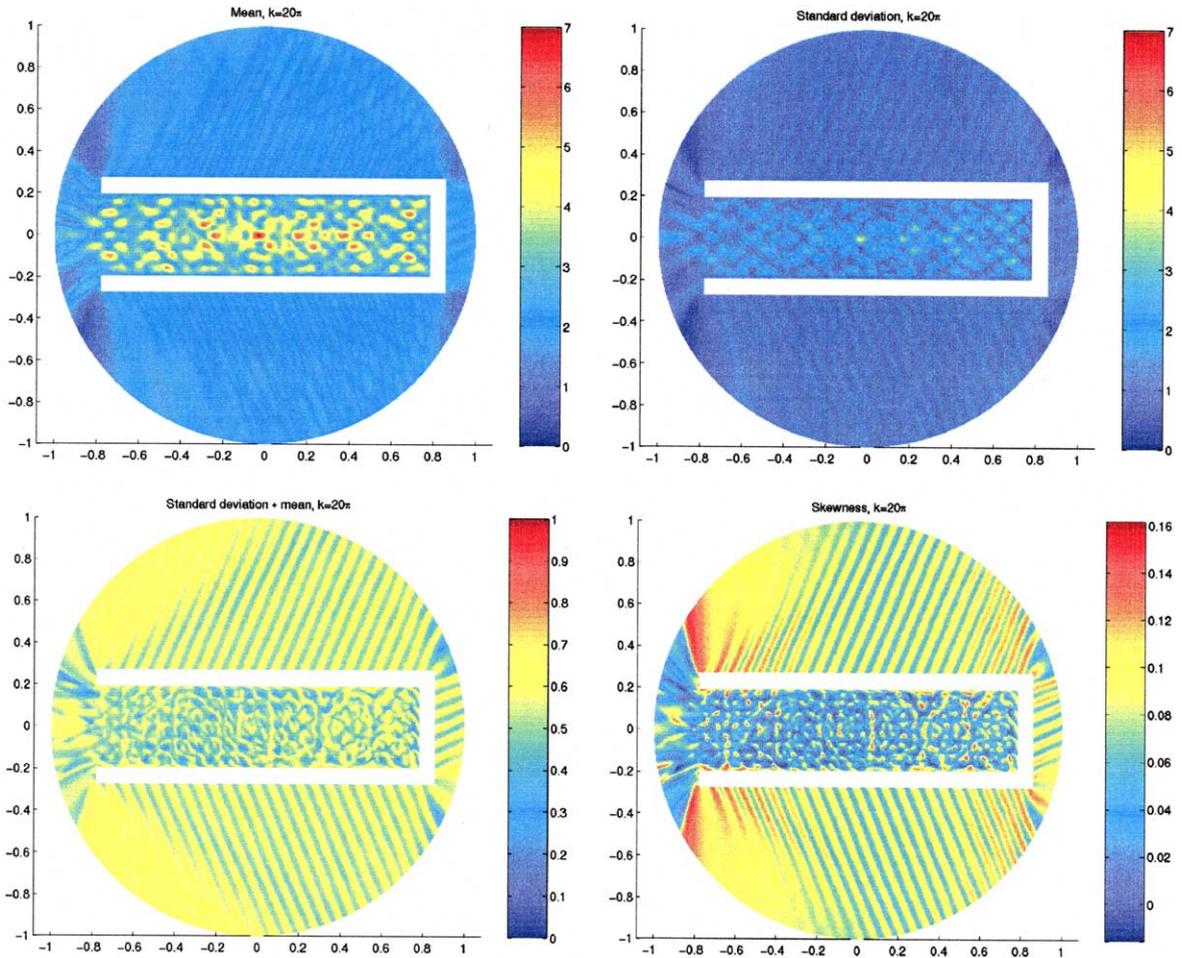


Fig. 7. Estimated mean μ_h , standard deviation σ_h , ratio σ_h/μ_h , and skewness of $|u_h|$, for $k = 20\pi$.

5. Concluding remarks

Our aim in this work was to carefully outline the computational issues associated with implementing the stochastic finite element method and processing the results for a model of acoustic scattering, where uncertainty is restricted to boundary data. We have shown that a representation of the solution in outer product form leads to significant savings in storage and also enables the relatively inexpensive computation of the random solution. The dominant cost comes from postprocessing of the solution to compute statistical data, although the outer product form in this setting reduces the storage overhead of these computations. Finally, we note that if uncertainty appears in the differential operator instead of the right-hand side, then the outer product formulation of the stochastic system is not available, and this problem would be more costly to solve.

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