

# EFFICIENT INVERSION OF THE GALERKIN MATRIX OF GENERAL SECOND-ORDER ELLIPTIC OPERATORS WITH NONSMOOTH COEFFICIENTS

MARIO BEBENDORF

ABSTRACT. This article deals with the efficient (approximate) inversion of finite element stiffness matrices of general second-order elliptic operators with  $L^\infty$ -coefficients. It will be shown that the inverse stiffness matrix can be approximated by hierarchical matrices ( $\mathcal{H}$ -matrices). Furthermore, numerical results will demonstrate that it is possible to compute an approximate inverse with almost linear complexity.

## 1. INTRODUCTION

We are concerned with the numerical solution of large finite element systems of Dirichlet problems,

$$\begin{aligned} Lu &= f && \text{in } \Omega, \\ u &= g && \text{on } \partial\Omega, \end{aligned}$$

with second-order elliptic operators

$$(1) \quad Lu = -\operatorname{div}[A\nabla u + bu] + c \cdot \nabla u + du$$

on bounded Lipschitz domains  $\Omega \subset \mathbb{R}^n$ . The Galerkin matrix of such operators is sparse but has a bandwidth of order  $N^{1-1/n}$ , where  $N$  is the number of degrees of freedom. Therefore, direct methods are well suited for small problem sizes, but are not competitive if  $N$  is large. In the latter case, iterative methods are usually more efficient. On the other hand, if the coefficient matrix is ill conditioned, these methods suffer from slow convergence.

The aim of this article is to show that hierarchical matrices ( $\mathcal{H}$ -matrices) introduced by Hackbusch et al. [15, 16] can fill this gap. As we will see, they provide a means by which an approximation of the inverse stiffness matrix can be generated and handled with logarithmic-linear complexity. Furthermore, no grid hierarchy is required and  $\mathcal{H}$ -matrices are robust in the sense that their efficiency does not depend on the smoothness and only slightly on the size of the coefficients. The  $\mathcal{H}$ -approximant might be used directly to solve the finite element system or it might be used as a black-box preconditioner in an iterative scheme. If it is used as a preconditioner, there is no need to approximate the inverse with high accuracy, and the complexity can therefore even be reduced. Another application of the

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inverse is Schur complements, which play a central role, for example, in domain decomposition methods.

The structure of  $\mathcal{H}$ -matrices was originally designed to efficiently represent integral operators with asymptotically smooth kernel. For this application the existence of  $\mathcal{H}$ -matrix approximants is well understood. Even efficient algorithms for the generation of the approximants exist; see [1, 3]. Close to the application of integral operators are inverses of elliptic operators since they have an integral representation with the Green function as kernel function:

$$(2) \quad (L^{-1}\varphi)(x) = \int_{\Omega} G(x, y)\varphi(y) \, dy \quad \text{for all } \varphi \in C_0^\infty(\Omega).$$

If  $L$  has smooth coefficients, the Green function is smooth (except for  $x = y$ ) and the mentioned existence theorems apply. However, the algorithms for building the approximants cannot be used, since they are either based on the matrix entries or on the kernel function, neither of which is accessible in general.

In this article the case of  $L^\infty$ -coefficients is treated. In this case it is not obvious that an  $\mathcal{H}$ -matrix approximant exists, since according to the De Giorgi-Nash theorem (see [10]),  $G$  is only locally Hölder continuous. Therefore, proofs cannot rely on the smoothness of the kernel function as they did for integral operators. In [4] we were able to show that the inverse stiffness matrix of the principal parts, i.e.,  $\mathbf{b} = \mathbf{c} = 0$  and  $d = 0$ , can be approximated by  $\mathcal{H}$ -matrices. In the present article this result will be extended to operators (1) with lower-order terms without restrictions on their size. Furthermore, we will present numerical results that, by an  $\mathcal{H}$ -matrix inversion based on the Frobenius formulas, one is able to (approximately) invert the stiffness matrix with logarithmic-linear complexity.

The structure of the rest of this article is as follows. In Section 2 a brief review of the structure of  $\mathcal{H}$ -matrices will be given. All necessary results and notation for the theorems of this article from the field of  $\mathcal{H}$ -matrices will be presented. Section 3 contains the existence theory of degenerate kernel approximants of the Green function, i.e.,

$$G(x, y) \approx \sum_{i=1}^k u_i(x)v_i(y), \quad x \in D_1 \text{ and } y \in D_2,$$

on an appropriate pair of domains  $(D_1, D_2)$ . The usual way to prove existence of a degenerate kernel approximant is to exploit the smoothness of the kernel function. In the case of  $L^\infty$ -coefficients the Green function of the inverse differential operator is not smooth. Therefore, another technique, which is based on the interior regularity of elliptic problems, has to be used. We will show that as long as  $L$  is invertible, the Green function can be approximated by a degenerate kernel even in the case of dominating lower-order terms. From the numerical results it will be seen that these terms enter the constants only in a moderate way.

This result is then employed using (2) to show that the discrete inverse of  $L$  can be approximated by  $\mathcal{H}$ -matrices, which in turn leads to the existence of  $\mathcal{H}$ -matrix approximants to the inverse stiffness matrix. In Section 4 we describe in detail how to implement an efficient  $\mathcal{H}$ -matrix inversion. The  $\mathcal{H}$ -inversion presented is based on the Frobenius formulas and is used to produce numerical results for operators with nonsmooth coefficients. We will see that this algorithm is able to compute an approximate inverse with almost linear complexity.

## 2. HIERARCHICAL MATRICES

This section gives a brief overview over the structure of  $\mathcal{H}$ -matrices originally introduced by Hackbusch et al. [15, 16]. We will describe the two principles on which the efficiency of  $\mathcal{H}$ -matrices is based. These are the hierarchical partitioning of the matrix into blocks and the blockwise restriction to low-rank matrices. These principles were also used in the mosaic-skeleton method [19].

In contrast to other efficient methods like wavelet techniques [6, 7, 8], fast multipole and panel clustering (see [13], [17] and the references therein),  $\mathcal{H}$ -matrices concentrate on the matrix level. They are purely algebraic in the sense that once the  $\mathcal{H}$ -matrix approximant is built, no further information about the underlying problem is needed.

Let us assume that  $M \in \mathbb{R}^{N \times N}$  has indices

$$(3) \quad m_{ij} = a(\varphi_j, \varphi_i),$$

where  $\varphi_i$  are basis functions with support  $X_i := \text{supp } \varphi_i$ ,  $i \in I := \{1, \dots, N\}$ , and  $a$  is a bilinear form. In this section we assume that there is a partition  $P$  of the indices  $I \times I$  of  $M$  such that each block  $b = s \times t$ ,  $s, t \subset I$ , can be approximated by a matrix of low rank, i.e.,

$$M_b \approx UV^T, \quad U \in \mathbb{R}^{s \times k}, V \in \mathbb{R}^{t \times k},$$

where  $k$  is small compared with  $|s|$  and  $|t|$ . Obviously, by  $M_b$  we denote the subblock in the intersection of the rows  $s$  and columns  $t$  of  $M$ . From Example 2.4 it will be seen that the stiffness matrix of operators of type (1) possesses this property. Section 3 will extend this to the inverse stiffness matrix.

**2.1. The cluster tree.** In order to exploit the fact that there is a partition such that each block can be approximated by a matrix of low rank, we first have to find it from the set of possible subsets of  $I \times I$ . This set, however, is too large to be searched for a partition that will satisfy our needs. Therefore, the set of subsets  $b = s \times t$  is restricted to those which consist of index sets  $s$  and  $t$  stemming from a cluster tree  $T_I$ . A tree  $T_I$  satisfying the following conditions is called a *cluster tree*:

- (1)  $I$  is the root of  $T_I$ ,
- (2) if  $t \in T_I$  is not a leaf, then  $t$  has sons  $t_1, t_2 \in T_I$ , so that  $t = t_1 \cup t_2$ .

The set of sons of  $t$  is denoted by  $S(t)$ , while  $\mathcal{L}(T_I)$  stands for the set of leaves of the tree  $T_I$ . The support of a cluster  $t$  is the union of the supports of the basis functions corresponding to the indices in  $t$ :

$$X_t := \bigcup_{i \in t} X_i.$$

A cluster tree is usually generated by recursive subdivision of  $I$  so as to minimize the diameter of each part. For practical purposes the recursion should be stopped if a certain cardinality  $n_{\min}$  of the clusters is reached, rather than subdividing the clusters until only one index is left. The depth of  $T_I$  will be denoted by  $p$ . For reasonable cluster trees one would always expect  $p = \mathcal{O}(\log N)$ . A strategy based on the *principle component analysis* is used in [2]. The complexity of building the cluster tree in the case of quasi-uniform grids can be estimated as  $\mathcal{O}(N \log N)$ .

**Remark 2.1.** Sometimes, the number of sons of a cluster in the previous definition of a cluster tree is not restricted to two. However, this generalization has not proved useful in practice.

**2.2. Admissibility condition.** In order to be able to approximate each block  $b$  of  $M$  by a low-rank matrix,  $b$  has to satisfy a certain condition. This so-called admissibility condition will be the criterion for choosing whether  $b$  belongs to  $P$ . In the field of elliptic partial differential equations the following condition on  $b = s \times t$  has proved useful:

$$(4) \quad \min\{\text{diam } X_s, \text{diam } X_t\} \leq \eta \text{dist}(X_s, X_t),$$

where  $\eta > 0$  is a given real number. We will see that under quite general assumptions this condition allows us to approximate the Green function of  $L$  by a degenerate kernel, i.e., there are functions  $u_i, v_i, i = 1, \dots, k$ , so that

$$(5) \quad G(x, y) \approx \sum_{i=1}^k u_i(x) v_i(y) \quad \text{in } X_s \times X_t,$$

where  $k$  depends only logarithmically on  $N$ . Since by (3) the entries of  $b$  depend only on the values of  $a$  on the domain  $X_s \times X_t$ , the degenerate approximation of  $G$  on  $X_s \times X_t$  will finally lead to a low-rank approximation of the block  $b$ .

Condition (4) was also used to prove convergence of the adaptive cross approximation (ACA) algorithm for the efficient generation of  $\mathcal{H}$ -matrix approximants in the case of integral equations (cf. [1, 3]).

**Remark 2.2.** In the case of unstructured grids the computation of the distance in (4) between two supports  $X_s$  and  $X_t$  is too costly. Therefore, for practical purposes, the supports are enclosed into sets of a simpler structure; e.g., boxes or spheres.

**2.3. Block cluster tree.** Based on a cluster tree  $T_I$  which contains a hierarchy of partitions of  $I$ , we are able to construct the so called *block cluster tree*  $T_{I \times I}$  describing a hierarchy of partitions of  $I \times I$  by the following rule.

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procedure build_block_cluster_tree( $s \times t$ )
begin
  if ( $s, t$ ) does not satisfy (4) and  $s, t \notin \mathcal{L}(T_I)$  then begin
     $S(s \times t) := \{s' \times t' : s' \in S(s), t' \in S(t)\}$ 
    for  $s' \times t' \in S(s \times t)$  do build_block_cluster_tree( $s' \times t'$ )
  end
  else  $S(s \times t) := \emptyset$ 
end

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Applying *build\_block\_cluster\_tree* to  $I \times I$ , we obtain a cluster tree for the index set  $I \times I$ . The set of leaves  $P := \mathcal{L}(T_{I \times I})$  is a partition of  $I \times I$  with blocks  $b = s \times t \in P$  either satisfying (4) or consisting of clusters  $t$  and  $s$ , one of which is a leaf in  $T_I$ . The complexity of building the block cluster tree in the case of quasi-uniform grids can be estimated as  $\mathcal{O}(\eta^{-n} N \log N)$  (cf. [2]).

We are now in a position to define the set of  $\mathcal{H}$ -matrices for a partition  $P$  with blockwise rank  $k$ :

$$\mathcal{H}(P, k) := \{M \in \mathbb{R}^{I \times I} : \text{rank } M_b \leq k \text{ for all } b \in P\}.$$

Note that  $\mathcal{H}(P, k)$  is not a linear space, since in general the sum of two rank  $k$  matrices exceeds rank  $k$ .

**Remark 2.3.** For a block  $B \in \mathbb{R}^{s \times t}$  the low-rank representation  $B = UV^T$ ,  $U \in \mathbb{R}^{s \times k}$ ,  $V \in \mathbb{R}^{t \times k}$ , is only advantageous compared with the entrywise representation if  $k(|s| + |t|) \leq |s||t|$ . For the sake of simplicity in this article, however, we will

assume that each block has the low-rank representation. Employing the entrywise representation for appropriate blocks will accelerate the algorithms.

**Example 2.4.** The stiffness matrix  $S$  of the differential operator  $L$  from (1) is in  $\mathcal{H}(P, n_{\min})$ . If  $b \in P$  satisfies (4), then the supports of the basis functions are pairwise disjoint. Hence, the matrix entries in this block vanish. In the remaining case  $b$  does not satisfy (4). Then the size of one of the clusters is less than or equal to  $n_{\min}$ . In both cases the rank of  $S_b$  does not exceed  $n_{\min}$ .

**2.4. Storage and matrix-vector multiplication.** The cost of multiplying an  $\mathcal{H}$ -matrix  $M \in \mathcal{H}(P, k)$  and its transposed  $M^T$  by a vector  $x \in \mathbb{R}^N$  is inherited from the blockwise matrix-vector multiplication

$$Mx = \sum_{s \times t \in P} M_{s \times t} x_t \quad \text{and} \quad M^T x = \sum_{s \times t \in P} (M_{s \times t})^T x_s.$$

Since each block  $s \times t$  has the representation  $M_{s \times t} = UV^T$ ,  $U \in \mathbb{R}^{s \times k}$ ,  $V \in \mathbb{R}^{t \times k}$  (see Remark 2.3),  $\mathcal{O}(k(|s| + |t|))$  units of memory are needed to store  $M_{s \times t}$  and the matrix-vector products

$$M_{s \times t} x_t = UV^T x_t \quad \text{and} \quad (M_{s \times t})^T x_s = VU^T x_s$$

can be done in  $\mathcal{O}(k(|s| + |t|))$  operations. Exploiting the hierarchical structure of  $M$ , it can therefore be shown that both storing  $M$  and multiplying  $M$  and  $M^T$  by a vector has  $\mathcal{O}(\eta^{-n} k N \log N)$  complexity. For a rigorous analysis the reader is referred to [2]. Therefore,  $\mathcal{H}$ -matrices are well suited for iterative schemes such as Krylov subspace methods.

### 3. APPROXIMATION OF FE INVERSES

In Example 2.4 it was mentioned that the stiffness matrix of a general elliptic operator with  $L^\infty$ -coefficients can be represented as an  $\mathcal{H}$ -matrix. In this section it will be proved, moreover, that its inverse can be approximated by an  $\mathcal{H}$ -matrix. For this purpose it will first be shown that the Green function of  $L$  and the bounded Lipschitz domain  $\Omega \subset \mathbb{R}^n$  can be approximated on a pair  $D_1 \times D_2$  of domains satisfying

$$\text{diam } D_2 \leq \eta \text{dist}(D_1, D_2).$$

Since we consider Dirichlet problems, we assume that  $L : H_0^1(\Omega) \rightarrow H^{-1}(\Omega)$  is an invertible second-order partial differential operator

$$(6) \quad Lu = -\text{div}[A \nabla u + \mathbf{b}u] + \mathbf{c} \cdot \nabla u + du,$$

where  $A(x) \in \mathbb{R}^{n \times n}$  is symmetric with entries  $a_{ij} \in L^\infty(\Omega)$  and

$$(7) \quad 0 < \lambda \leq \lambda(x) \leq \Lambda$$

for all eigenvalues  $\lambda(x)$  of  $A(x)$  and almost all  $x \in \Omega$ . The bound  $\Lambda/\lambda$  on the condition numbers of  $A$  will be denoted by  $\kappa := \Lambda/\lambda$ . Furthermore, let  $\mathbf{b}(x), \mathbf{c}(x) \in \mathbb{R}^n$  and  $d(x) \in \mathbb{R}$  with  $b_i, c_i, d \in L^\infty(\Omega)$ ,  $i = 1, \dots, n$ .

**3.1. Degenerate approximation of the Green function.** In [4] we investigated the principal part  $L_0$  of such operators; i.e., operators with  $\mathbf{b} = \mathbf{c} = 0$  and  $d = 0$ . For these operators it is shown in [12] that in the case  $n \geq 3$  a Green function  $G_0 : \Omega \times \Omega \rightarrow \overline{\mathbb{R}}$  exists with the properties

$$(8a) \quad G_0(x, \cdot) \in H^1(\Omega \setminus B_r(x)) \cap W_0^{1,1}(\Omega) \text{ for all } x \in \Omega \text{ and all } r > 0,$$

$$(8b) \quad a(G_0(x, \cdot), \varphi) = \varphi(x) \text{ for all } \varphi \in C_0^\infty(\Omega) \text{ and } x \in \Omega,$$

where  $B_r(x)$  is the open ball centered at  $x$  with radius  $r$  and

$$(9) \quad a(u, v) = \int_{\Omega} \nabla v \cdot A \nabla u \, dx.$$

Furthermore, for  $x, y \in \Omega$  it holds that

$$(10) \quad |G_0(x, y)| \leq \frac{c_n(\kappa)}{\lambda} |x - y|^{2-n}.$$

In the case of invertible operators  $L$  of type (6) a Green function  $G := (L_0^{-1}L)^{-1}G_0$  satisfying (8) can be defined, where

$$(11) \quad a(u, v) = \int_{\Omega} \nabla v \cdot A \nabla u \, dx + \int_{\Omega} \nabla v \cdot \mathbf{b} u \, dx + \int_{\Omega} \mathbf{c} \cdot \nabla u v \, dx + \int_{\Omega} duv \, dx.$$

Notice that

$$G - G_0 = [(L_0^{-1}L)^{-1} - I]G_0 = -L^{-1}L_1G_0,$$

where  $L_1 := L - L_0$  is the lower-order part of  $L$ . Since  $L^{-1}L_1$  is an operator of order  $-1$ ,  $L^{-1}L_1G_0$  is smoother than  $G_0$ . Hence, the singularity of  $G_0$  at  $x = y$  is carried over to  $G$  and we may assume that there is a constant  $c_n$  such that for  $x, y \in \Omega$  it holds that

$$(12) \quad |G(x, y)| \leq \frac{c_n(\kappa, \mathbf{b}, \mathbf{c}, d)}{\lambda} |x - y|^{2-n}.$$

In the case  $n = 2$ , the existence of a Green function for operators of type (6) has been proved more rigorously in [9]. In this case instead of (10) for  $x, y \in \Omega$ , one has the following bound on the Green function:

$$(13) \quad |G(x, y)| \leq \frac{c(\kappa, \mathbf{b}, \mathbf{c}, d)}{\lambda} \log |x - y|.$$

We will make use of the following characteristic relation between  $L^{-1}$  and  $G$ , which is equivalent to (8b):

$$(14) \quad (L^{-1}\varphi)(x) = \int_{\Omega} G(x, y)\varphi(y) \, dy \quad \text{for all } \varphi \in C_0^\infty(\Omega).$$

In the rest of this paper  $D \subset \mathbb{R}^n$  is a domain. The proof of the following basic lemma is mainly based on the Poincaré inequality (cf. [5]), and can be found in [4].

**Lemma 3.1.** *Let  $D$  be convex and  $X$  a closed subspace of  $L^2(D)$ . Then for any  $k \in \mathbb{N}$  there is a subspace  $V_k \subset X$  satisfying  $\dim V_k \leq k$  so that*

$$(15) \quad \text{dist}_{L^2(D)}(u, V_k) \leq c_A \frac{\text{diam } D}{\sqrt[n]{k}} \|\nabla u\|_{L^2(D)}$$

for each  $u \in X \cap H^1(D)$ , where  $c_A$  depends only on the spatial dimension  $n$ .

The following set will be used in Lemma 3.1 as  $X$ :

$$(16) \quad X(D) = \{u \in H_{\text{loc}}^1(D) : a(u, \varphi) = 0 \ \forall \varphi \in C_0^\infty(D \cap \Omega) \text{ and } u|_{D \setminus \Omega} = 0\},$$

where  $a(\cdot, \cdot)$  is the bilinear form defined in (11). Hence, the set  $X(D)$  consists of  $L$ -harmonic  $H_{\text{loc}}^1$ -functions vanishing outside of  $\Omega$ . A proof for the fact that  $X(D)$  is a closed subspace of  $L^2(D)$  can be found in [4, Lemma 2.2]. We remark that the extension of  $G(x, \cdot)$ ,  $x \in \Omega$ , to  $\mathbb{R}^n$  by zero is in  $X(D)$  for all  $D \subset \mathbb{R}^n$  satisfying  $\text{dist}(x, D) > 0$ .

By the following Caccioppoli inequality we are able to estimate the gradient on a compact subset by the norm of the function on an enclosing domain. This inequality provides a means to overcome the lack of regularity of  $G$ .

**Lemma 3.2.** *Let  $K \subset D$  be a compact subset. There is  $c_R = c_R(\kappa, \lambda, \mathbf{b}, \mathbf{c}, d)$  such that*

$$(17) \quad \|\nabla u\|_{L^2(K)} \leq \frac{c_R}{\text{dist}(K, \partial D)} \|u\|_{L^2(D)}$$

for all  $u \in X(D)$ .

*Proof.* Let  $\eta \in C^1(D)$  satisfy  $0 \leq \eta \leq 1$ ,  $\eta = 1$  in  $K$ ,  $\eta = 0$  in a neighborhood of  $\partial D$  and  $|\nabla \eta| \leq 2/\delta$  in  $D$ , where we set  $\delta = \text{dist}(K, \partial D)$ . Since  $K' := \text{supp } \eta$  is a compact subset of  $D$ , definition (16) of  $X(D)$  implies  $u \in H^1(K')$ . Hence,  $\varphi := \eta^2 u \in H_0^1(D \cap \Omega)$  may be used as a test function in  $a(u, \varphi) = 0$  due to the dense embedding of  $C_0^\infty(D \cap \Omega)$  in  $H_0^1(D \cap \Omega)$ . Since  $\varphi = 0$  in  $D \setminus \Omega$ , we have

$$\begin{aligned} - \int_D (\nabla \eta^2 u) \cdot \mathbf{b} u \, dx - \int_D \eta^2 u \mathbf{c} \cdot \nabla u \, dx - \int_D d \eta^2 |u|^2 \, dx &= \int_D (\nabla \eta^2 u) \cdot A \nabla u \, dx \\ &= 2 \int_D \eta u (\nabla \eta) \cdot A \nabla u \, dx + \int_D \eta^2 (\nabla u) \cdot A \nabla u \, dx. \end{aligned}$$

Hence,

$$\begin{aligned} \int_D \eta^2 |A^{1/2} \nabla u|^2 \, dx &= -2 \int_D \eta u (\nabla \eta) \cdot A \nabla u \, dx - 2 \int_D \eta (\nabla \eta) \cdot \mathbf{b} |u|^2 \, dx \\ &\quad - \int_D \eta^2 (\nabla u) \cdot \mathbf{b} u \, dx - \int_D \eta^2 u \mathbf{c} \cdot \nabla u \, dx - \int_D d \eta^2 |u|^2 \, dx. \end{aligned}$$

For the first integral on the right-hand side of the last equation we obtain

$$\begin{aligned} \left| \int_D \eta u (\nabla \eta) \cdot A \nabla u \, dx \right| &\leq \int_D |A^{1/2} \nabla \eta| |\eta A^{1/2} \nabla u| |u| \, dx \\ &\leq 2 \frac{\sqrt{\Lambda}}{\delta} \|\eta A^{1/2} \nabla u\|_{L^2(D)} \|u\|_{L^2(D)}. \end{aligned}$$

The third integral can be estimated as

$$\begin{aligned} \left| \int_D \eta^2 (\nabla u) \cdot \mathbf{b} u \, dx \right| &\leq \int_D \eta |\mathbf{b}| \eta |\nabla u| |u| \, dx \leq \|\mathbf{b}\|_\infty \int_D \eta |\nabla u| |u| \, dx \\ &\leq \|\mathbf{b}\|_\infty \left( \int_D \eta^2 |\nabla u|^2 \, dx \right)^{1/2} \|u\|_{L^2(D)} \\ &\leq \frac{\|\mathbf{b}\|_\infty}{\sqrt{\lambda}} \|\eta A^{1/2} \nabla u\|_{L^2(D)} \|u\|_{L^2(D)} \end{aligned}$$

and, similarly, one has for the fourth integral

$$\left| \int_D \eta^2 u \mathbf{c} \cdot \nabla u \, dx \right| \leq \frac{\|\mathbf{c}\|_\infty}{\sqrt{\lambda}} \|\eta A^{1/2} \nabla u\|_{L^2(D)} \|u\|_{L^2(D)}.$$

Since

$$2\|\eta A^{1/2} \nabla u\|_{L^2(D)} \|u\|_{L^2(D)} \leq \frac{1}{\epsilon} \|\eta A^{1/2} \nabla u\|_{L^2(D)}^2 + \epsilon \|u\|_{L^2(D)}^2$$

with  $\epsilon := 4\sqrt{\Lambda}/\delta + \lambda^{-1/2}(\|\mathbf{b}\| + \|\mathbf{c}\|_\infty)$ , we end up with

$$\|\eta A^{1/2} \nabla u\|_{L^2(D)}^2 \leq 2 \left( \frac{\epsilon^2}{2} + \frac{4}{\delta} \|\mathbf{b}\|_\infty + \|d\|_\infty \right) \|u\|_{L^2(D)}^2.$$

This leads to

$$\|\eta A^{1/2} \nabla u\|_{L^2(D)} \leq \frac{c}{\delta} \|u\|_{L^2(D)},$$

where

$$c^2 = \left( 4\sqrt{\Lambda} + \frac{\delta}{\sqrt{\lambda}} (\|\mathbf{b}\| + \|\mathbf{c}\|_\infty) \right)^2 + 8\delta \|\mathbf{b}\|_\infty + 2\delta^2 \|d\|_\infty,$$

and hence

$$\|\nabla u\|_{L^2(K)} \leq \|\eta \nabla u\|_{L^2(D)} \leq \lambda^{-1/2} \|\eta A^{1/2} \nabla u\|_{L^2(D)} \leq \frac{c}{\lambda^{1/2} \delta} \|u\|_{L^2(D)}. \quad \square$$

**Remark 3.3.** From the previous proof it can be seen that  $d$  does not enter estimate (17) if  $d \geq 0$ .

**Lemma 3.4.** Assume that  $D_2 \subset D$  is a convex domain such that for some  $\eta > 0$  it holds that

$$0 < \text{diam } D_2 \leq \eta \text{dist}(D_2, \partial D).$$

Then for any  $\varepsilon > 0$  there is a subspace  $W \subset X(D_2)$  so that

$$(18) \quad \text{dist}_{L^2(D_2)}(u, W) \leq \varepsilon \|u\|_{L^2(D)} \quad \text{for all } u \in X(D),$$

and  $\dim W \leq c_\eta^n [\log \varepsilon]^{n+1} + [\log \varepsilon]$ , where  $c_\eta = c_A c_R e(2 + \eta)$ .

*Proof.* Let  $\ell := [\log \varepsilon]$ . We consider a nested sequence of convex domains

$$K_j = \{x \in \mathbb{R}^n : \text{dist}(x, D_2) \leq r_j\}$$

with real numbers  $r_j := (1 - j/\ell) \text{dist}(D_2, \partial D)$ ,  $j = 0, \dots, \ell$ . Notice that  $D_2 = K_\ell \subset K_{\ell-1} \subset \dots \subset K_0 \subset D$ . Using the definition (16) of the space  $X$  we set  $X_j := X(K_j)$ .

Applying Lemma 3.1 to  $K_j$  with the choice  $k := \lceil (c_A c_R (2 + \eta) \ell \varepsilon^{-1/\ell})^n \rceil$  we can find a subspace  $V_j \subset X_j$  satisfying  $\dim V_j \leq k$  and

$$(19) \quad \text{dist}_{L^2(K_j)}(v, V_j) \leq c_A \frac{\text{diam } K_j}{\sqrt[n]{k}} \|\nabla v\|_{L^2(K_j)}$$

for all  $v \in X_j \cap H^1(K_j)$ . From Lemma 3.2 applied to  $(K_j, K_{j-1})$  instead of  $(K, D)$ , we obtain

$$(20) \quad \|\nabla v\|_{L^2(K_j)} \leq \frac{c_R}{\text{dist}(K_j, \partial K_{j-1})} \|v\|_{L^2(K_{j-1})} = c_R \frac{\ell}{r_0} \|v\|_{L^2(K_{j-1})}$$

for all  $v \in X_{j-1}$ . Since any  $v \in X_{j-1}$  also belongs to  $X_j \cap H^1(K_j)$ , the estimates (19) and (20) together with  $\text{diam } K_j \leq (2 + \eta)r_0$  may be combined:

$$(21) \quad \text{dist}_{L^2(K_j)}(v, V_j) \leq \varepsilon^{1/\ell} \|v\|_{L^2(K_{j-1})} \quad \text{for all } v \in X_{j-1}.$$



Let  $u \in X(D)$  and  $v_0 := u|_{K_0} \in X_0$ . By the last estimate we have  $v_0|_{K_1} = u_1 + v_1$  with  $u_1 \in V_1$  and

$$\|v_1\|_{L^2(K_1)} \leq \varepsilon^{1/\ell} \|v_0\|_{L^2(K_0)}.$$

Consequently,  $v_1$  belongs to  $X_1$ . Similarly, for all  $j = 1, \dots, \ell$ , we are able to find an approximant  $u_j \in V_j$  so that  $v_{j-1}|_{K_j} = u_j + v_j$  and  $\|v_j\|_{L^2(K_j)} \leq \varepsilon^{1/\ell} \|v_{j-1}\|_{L^2(K_{j-1})}$ . Using the restrictions of  $V_j$  to the smallest domain  $D_2 = K_\ell$ , let

$$W := \text{span}\{V_j|_{D_2}, j = 1, \dots, \ell\}.$$

Then  $W$  is a subspace of  $X(D_2)$  and, since  $v_0|_{D_2} = v_\ell + \sum_{j=1}^\ell u_j|_{D_2}$ , we are led to

$$\text{dist}_{L^2(D_2)}(v_0, W) \leq \|v_\ell\|_{L^2(D_2)} \leq \left(\varepsilon^{1/\ell}\right)^\ell \|v_0\|_{L^2(K_0)} \leq \varepsilon \|u\|_{L^2(D)},$$

where the last inequality is due to  $K_0 \subset D$ .

The dimension of  $W$  is bounded by  $\sum_{j=1}^\ell \dim V_j \leq \ell k$ . Since  $\varepsilon^{-1/\ell} \leq e$  we obtain  $\dim W \leq (c_{ACRE}(2 + \eta))^n \ell^{n+1} + \ell$ .  $\square$

The previous lemma will now be applied to the Green functions  $G(x, \cdot)$  with  $x \in D_1 \subset \Omega$ . For this purpose let  $g_x$  be the extension of  $G(x, \cdot)$  to  $\mathbb{R}^n \setminus \overline{D_1}$ ; i.e.,

$$(22) \quad g_x(y) := \begin{cases} G(x, y), & y \in \Omega \setminus \overline{D_1}, \\ 0, & y \in \mathbb{R}^n \setminus \Omega. \end{cases}$$

Then  $g_x$  is in  $X(\mathbb{R}^n \setminus \overline{D_1})$ . Note that its approximant  $G_k(x, \cdot)$  from the following theorem is of the desired form (5).

**Theorem 3.5.** *Let  $D_1 \subset \Omega$  and  $D_2 \subset \mathbb{R}^n$  convex. Assume that there is  $\eta > 0$  such that*

$$0 < \text{diam } D_2 \leq \eta \text{dist}(D_1, D_2).$$

*Then for any  $\varepsilon > 0$  there is a separable approximation*

$$G_k(x, y) = \sum_{i=1}^k u_i(x) v_i(y) \quad \text{with } k \leq k_\varepsilon := c_\eta^n \lceil \log \varepsilon \rceil^{n+1} + \lceil \log \varepsilon \rceil,$$

*so that for all  $x \in D_1$ ,*

$$(23) \quad \|G(x, \cdot) - G_k(x, \cdot)\|_{L^2(D_2 \cap \Omega)} \leq \varepsilon \|G(x, \cdot)\|_{L^2(\hat{D}_2)},$$

*where  $\hat{D}_2 := \{y \in \Omega : 2\eta \text{dist}(y, D_2) < \text{diam } D_2\}$  and*

$$c_\eta = 2c_A e(1 + \eta) \sqrt{\left(4\sqrt{\kappa} + \frac{\delta}{\lambda} \|\mathbf{b}\| + \|\mathbf{c}\|_\infty\right)^2 + 8\frac{\delta}{\lambda} \|\mathbf{b}\|_\infty + 2\frac{\delta^2}{\lambda} \|d\|_\infty}, \quad \delta := \frac{\text{diam } D_2}{2\eta}.$$

*Proof.* Let  $D = \{y \in \mathbb{R}^n : 2\eta \text{dist}(y, D_2) < \text{diam } D_2\}$ . Note that because of  $\text{dist}(D_1, D) > 0$ , we have  $g_x \in X(D)$  for all  $x \in D_1$ . Since in addition  $\text{diam } D_2 \leq 2\eta \text{dist}(D_2, \partial D)$ , Lemma 3.4 can be applied with  $\eta$  replaced by  $2\eta$ . Let  $\{v_1, \dots, v_k\}$  be a basis of the subspace  $W \subset X(D_2)$  with  $k = \dim W \leq c_{2\eta}^n \lceil \log \varepsilon \rceil^{n+1} + \lceil \log \varepsilon \rceil$ . By means of (18)  $g_x$  can be decomposed into  $g_x = \hat{g}_x + r_x$  with  $\hat{g}_x \in W$  and  $\|r_x\|_{L^2(D_2)} \leq \varepsilon \|g_x\|_{L^2(D)}$ . Since  $g_x$  and  $\hat{g}_x$  vanish outside of  $\Omega$ , we actually have  $\|r_x\|_{L^2(D_2 \cap \Omega)} \leq \varepsilon \|G(x, \cdot)\|_{L^2(\hat{D}_2)}$ . Expressing  $\hat{g}_x$  by means of the basis of  $W$ , we obtain

$$\hat{g}_x = \sum_{i=1}^k u_i(x) v_i$$

with coefficients  $u_i(x)$  depending on the index  $x \in D_1$ . The function  $G_k(x, y) := \sum_{i=1}^k u_i(x) v_i(y)$  satisfies estimate (23).  $\square$

Note that since the constant  $c_A$  does not depend on  $\Omega$ , the geometry enters  $c_\eta$  only through the diameter. Hence, the shape of the domain does not influence our approximation result.

The existence of degenerate approximants to the Green function will now be used to prove existence of  $\mathcal{H}$ -matrix approximants to the discrete inverse of  $L$  and the inverse stiffness matrix.

**3.2.  $\mathcal{H}(P, k)$ -approximation of discrete operators.** Using a finite element discretization,  $H_0^1(\Omega)$  is approximated by  $V_h \subset H_0^1(\Omega)$ ; i.e., for all  $v \in H_0^1(\Omega)$

$$(24) \quad \inf_{v_h \in V_h} \|v - v_h\|_{H^1} \rightarrow 0 \quad \text{for } h \rightarrow 0.$$

In agreement with the assumptions of Section 2, let  $N = \dim V_h$  be the dimension and  $\{\varphi_i\}_{i \in I}$  a basis of  $V_h$ , where  $I := \{1, \dots, N\}$  is used as an index set. The notation for the support of the finite element basis function is generalized to subsets  $t \subset I$  as

$$(25) \quad X_i := \text{supp } \varphi_i \subset \Omega \quad \text{for } i \in I, \quad X_t := \bigcup_{i \in t} X_i \quad \text{for } t \subset I.$$

In order to avoid technical complications, we consider a *quasi-uniform* and *shape-regular* triangulation. Hence, the step size  $h := \max_{i \in I} \text{diam } X_i$  fulfills

$$(26) \quad \text{vol } X_i \geq c_v h^n.$$

The supports  $X_i$  may overlap. In accordance with the standard finite element discretization, we require that each triangle belongs to the support of a bounded number of basis functions; i.e., there is a constant  $c_M > 0$  so that

$$(27) \quad c_M \text{vol } X_t \geq \sum_{i \in t} \text{vol } X_i.$$

We use the notation  $J$  for the natural bijection  $J : \mathbb{R}^N \rightarrow V_h$  defined by  $Jx = \sum_{i \in I} x_i \varphi_i$ . For quasi-uniform and shape-regular triangulations it is known (see [14, Theorem 8.8.1]) that there are constants  $0 < c_{J,1} \leq c_{J,2}$  (independent of  $h$  and  $N$ ) such that

$$(28) \quad c_{J,1} \|x\|_h \leq \|Jx\|_{L^2(\Omega)} \leq c_{J,2} \|x\|_h \quad \text{for all } x \in \mathbb{R}^N,$$

where  $\|\cdot\|_h$  is the naturally scaled Euclidean norm induced by the scalar product  $\langle x, y \rangle_h = h^n \sum_{i \in I} x_i y_i$ . Since  $J$  is also a function from  $\mathbb{R}^N$  to  $H_0^1(\Omega)$ , the adjoint  $J^* \in L(H^{-1}(\Omega), \mathbb{R}^N)$  with respect to  $\langle \cdot, \cdot \rangle_h$  is defined. We define the following three  $N \times N$  matrices,

$$S = J^* L J, \quad B = J^* L^{-1} J, \quad \text{and} \quad M = J^* J.$$

$S$  is the stiffness matrix,  $B$  the Galerkin discretization of the inverse of  $L$ , and  $M$  is the mass matrix. The matrices  $S$  and  $M$  are sparse, while  $B$  as well as  $S^{-1}$  and  $M^{-1}$  are dense.

**Remark 3.6.**  $M$  is positive definite and  $S$  is invertible for sufficiently small  $h$ . Since the principal part of  $L$  is coercive and the lower-order terms constitute a compact operator due to the compact embedding of  $L^2$  in  $H^{-1}$ ,  $L$  satisfies Gårding's inequality

$$(Lu, u)_{L^2(\Omega)} \geq \gamma \|u\|_{H^1(\Omega)}^2 - c \|u\|_{L^2(\Omega)}^2 \quad \text{for all } u \in H^1(\Omega).$$

From the Céa-Polski Lemma (cf. [18]),  $S$  is invertible if  $h$  is sufficiently small.

We need the following lemma [11] by which the spectral norm of an  $\mathcal{H}$ -matrix can be estimated by its blockwise norms.  $P$  is again assumed to be generated as in subsection 2.3.

**Lemma 3.7.** *There is a constant  $c_{\text{sp}}$  such that for any matrix  $M \in \mathcal{H}(P, k)$  the following inequality holds:*

$$\|M\|_2 \leq c_{\text{sp}} p \max_{b \in P} \|M_b\|_2.$$

**Theorem 3.8.** *Let  $X_t$  be convex for all  $t \in T_I$ . For any  $\varepsilon > 0$ , let  $k_\varepsilon \in \mathbb{N}$  be chosen as in Theorem 3.5. Then for  $k \geq \max\{k_\varepsilon, n_{\min}\}$ , there is  $B_{\mathcal{H}} \in \mathcal{H}(P, k)$  such that*

$$(29) \quad \|B - B_{\mathcal{H}}\|_2 \leq c_n \frac{\varepsilon}{\lambda} p,$$

where  $c_n = c_n(\kappa, \mathbf{b}, \mathbf{c}, d, \eta, \Omega)$  depends on  $\eta$  from (4) and  $\text{diam } \Omega$ .  $p$  is the depth of the cluster tree  $T_I$  defined in subsection 2.1.

*Proof.* Let  $b = s \times t \in P$  with  $\min\{\#s, \#t\} \leq n_{\min}$ . In this case we simply set

$$(B_{\mathcal{H}})_b := B_b = (J^* L^{-1} J)_b.$$

Since the block  $(B_{\mathcal{H}})_b$  has at most  $n_{\min}$  columns or rows,  $\text{rank } (B_{\mathcal{H}})_b \leq k$  holds.

If  $b = s \times t \in P$  with  $\min\{\#s, \#t\} > n_{\min}$ , then  $b$  satisfies (4). Applying Theorem 3.5 with  $D_1 = X_s$ ,  $D_2 = X_t$  there is  $\tilde{G}_b(x, y) = \sum_{i=1}^{k_\varepsilon} u_i^b(x) v_i^b(y)$  such that

$$\|G - \tilde{G}_b\|_{L^2(X_s \times X_t)} \leq \varepsilon \|G\|_{L^2(X_s \times \hat{X}_t)},$$

where  $\hat{X}_t := \{x \in \Omega : 2\eta \text{dist}(x, X_t) \leq \text{diam } X_t\}$ . Let the functions  $u_i^b$  and  $v_i^b$  be extended to  $\Omega$  by zero. We define the integral operator

$$K_b \varphi = \int_{\Omega} \tilde{G}_b(\cdot, y) \varphi(y) dy \quad \text{for } \text{supp } \varphi \subset \bar{\Omega}$$

and set  $(B_{\mathcal{H}})_b = (J^* K_b J)_b$ . The rank of  $(B_{\mathcal{H}})_b$  is bounded by  $k_\varepsilon$  since each term  $u_i^b(x) v_i^b(y)$  in  $\tilde{G}_b$  produces one rank 1 matrix in  $(J^* K_b J)_b$ .

Let  $x \in \mathbb{R}^t$  and  $y \in \mathbb{R}^s$ . To see that  $(B_{\mathcal{H}})_b$  approximates the block  $B_b$ , remember the representation (14) of  $L^{-1}$  and use (28). The estimate

$$\begin{aligned} \langle (B - B_{\mathcal{H}})_b x, y \rangle_h &= \langle J^* (L^{-1} - K_b) J x, y \rangle_h = \langle (L^{-1} - K_b) J x, J y \rangle_{L^2} \\ &\leq \|G - \tilde{G}_b\|_{L^2(X_s \times X_t)} \|J x\|_{L^2(X_t)} \|J y\|_{L^2(X_s)} \\ &\leq \varepsilon \|G\|_{L^2(X_s \times \hat{X}_t)} \|J x\|_{L^2(\Omega)} \|J y\|_{L^2(\Omega)} \\ &\leq \varepsilon c_{J,2}^2 \|G\|_{L^2(X_s \times \hat{X}_t)} \|x\|_h \|y\|_h \end{aligned}$$

proves  $\|(B - B_{\mathcal{H}})_b\|_2 \leq \varepsilon c_{J,2}^2 \|G\|_{L^2(X_s \times \hat{X}_t)}$ .

Although  $G(x, \cdot) \in W^{1,1}(\Omega)$  for all  $x \in \Omega$ ,  $G(\cdot, \cdot)$  does not belong to  $L^2(\Omega \times \Omega)$  as soon as  $n \geq 4$ . From (12) it can be seen that  $\|G\|_{L^2(X_s \times \hat{X}_t)}$  may increase when the sets  $X_s$ ,  $\hat{X}_t$  are approaching each other. The construction of  $\hat{X}_t$ , however, ensures

$$\delta := \text{dist}(X_s, \hat{X}_t) \geq \frac{1}{2} \text{dist}(X_s, X_t) \geq \frac{1}{2\eta} \text{diam } X_s$$

as well as  $2\eta\delta \geq \text{diam } X_t$  due to (4). Hence (12) implies, for the case  $n \geq 3$ ,

$$\|G\|_{L^2(X_s \times \hat{X}_t)} \leq \frac{c_n(\kappa, \mathbf{b}, \mathbf{c}, d)}{\lambda} \delta^{2-n} \sqrt{(\text{vol } X_s)(\text{vol } \hat{X}_t)}.$$

Using  $\text{vol } \hat{X}_t \leq \omega_n(\frac{1}{2}\text{diam } \hat{X}_t)^n \leq \omega_n(\eta + 1/2)^n \delta^n$  and  $\text{vol } X_s \leq \omega_n(\eta\delta)^n$ , where  $\omega_n$  is volume of the unit ball in  $\mathbb{R}^n$ , we see that

$$\|G\|_{L^2(X_s \times \hat{X}_t)} \leq \bar{c}_\eta \frac{c_n(\kappa, \mathbf{b}, \mathbf{c}, d)}{\lambda} \delta^2 \quad \text{with } \bar{c}_\eta := \omega_n(\eta(\eta + 1/2))^{n/2}.$$

The rough estimate  $\delta \leq \text{diam } \Omega$  together with Lemma 3.7 yields (29). Using (13), the case  $d = 2$  can be treated in a similar way.  $\square$

**Remark 3.9.** Assume that each (possibly nonconvex) set  $X_t$  has a convex superset  $Y_t$  satisfying the admissibility condition (4). Then Theorem 3.8 remains valid for  $X_s \times X_t$ . Therefore, according to Remark 2.2, the assumption on the convexity of  $X_t$  in Theorem 3.8 is reasonable even for practical purposes.

The previous theorem shows that we are able to approximate the discrete inverse of  $L$  by  $\mathcal{H}$ -matrices. Our aim however is to prove that the inverse of the stiffness matrix  $S$  possesses this property. For this purpose we use the fact that  $S^{-1}$  can be approximated by  $M^{-1}BM^{-1}$ . The last product, in turn, can be approximated by an  $\mathcal{H}$ -matrix. In [4] we have already presented the details of the above arguments. Since they are quite technical, we just give the main results without proofs.

The finite element approximation is connected with the Ritz projection  $P_h = JA^{-1}J^*L : H_0^1(\Omega) \rightarrow V_h$ . If  $u \in H_0^1(\Omega)$  is the solution of the variational problem  $a(u, v) = f(v)$ ,  $u_h = P_h u$  is its finite element solution. The FE error is then given by

$$e_h(u) := \|u - P_h u\|_{L^2(\Omega)},$$

and the weakest form of the finite element convergence is described by

$$(30) \quad e_h(u) \leq \varepsilon_h \|f\|_{L^2(\Omega)} \quad \text{for all } u = L^{-1}f, \quad f \in L^2(\Omega),$$

where  $\varepsilon_h \rightarrow 0$  as  $h \rightarrow 0$ .

**Remark 3.10.** Due to our quite weak assumptions on the smoothness of the coefficients in (6), one cannot specify the behavior of  $\varepsilon_h$  for  $h \rightarrow 0$ .

**Lemma 3.11.** *It holds that  $\|S^{-1} - M^{-1}BM^{-1}\|_2 \leq 2c_{J,1}^{-4}c_{J,2}^2\varepsilon_h$ .*

Since the product of two  $\mathcal{H}$ -matrices is an  $\mathcal{H}$ -matrix with augmented rank (cf. [11]), it remains to show that  $M^{-1}$  can be approximated by an  $\mathcal{H}$ -matrix  $N_{\mathcal{H}}$ . Then,  $C_{\mathcal{H}} := N_{\mathcal{H}}B_{\mathcal{H}}N_{\mathcal{H}}$  approximates  $S^{-1}$ .

**Lemma 3.12.** *For any  $\varepsilon > 0$ , there is  $N_{\mathcal{H}} \in \mathcal{H}(P, k_\varepsilon)$  satisfying*

$$\|M^{-1} - N_{\mathcal{H}}\|_2 \leq \varepsilon \|M^{-1}\|_2$$

*with  $k_\varepsilon = \mathcal{O}(|\log \varepsilon|^n)$ .*

Gathering all previous results, we obtain the existence of  $\mathcal{H}$ -matrix approximants to the inverse stiffness matrix.

**Theorem 3.13.** *Let  $\varepsilon_h > 0$  be the finite element error from (30) and  $p$  the depth of the cluster tree  $T_I$  defined in subsection 2.1. Then there is a constant  $\tilde{c} > 0$  defining  $k := \tilde{c}p^2 \log^{n+1} \frac{p}{\varepsilon_h}$  and there is  $C_{\mathcal{H}} \in \mathcal{H}(P, k)$  such that*

$$(31) \quad \|S^{-1} - C_{\mathcal{H}}\|_2 \leq c_n \varepsilon_h,$$

*where  $c_n = c_n(\kappa, \lambda, \|L^{-1}\|_{H^1 \leftarrow H^{-1}}, \mathbf{b}, \mathbf{c}, d, \eta, \text{diam } \Omega)$ . If  $\varepsilon_h = \mathcal{O}(h^\beta)$  for some  $\beta > 0$ ,  $k = \mathcal{O}(\log^{n+3} N)$  holds.*

Theorem 3.13 states that  $S^{-1}$  can be approximated to an accuracy determined by the FE error, which is sufficient since the accuracy of the solution cannot be improved by a better approximation of  $S^{-1}$ . In the following section, however, an inversion algorithm is devised, which, as we will see in the numerical results, can reach any prescribed accuracy.

#### 4. ALGORITHMS

Since  $\mathcal{H}(P, k)$  is not a linear space, we have to replace the usual matrix operations by truncated ones. Starting from the  $\mathcal{H}$ -matrix addition, we define an  $\mathcal{H}$ -matrix multiplication. Using these modified operations it is possible to define an  $\mathcal{H}$ -matrix inversion based on the Frobenius formulas. These ideas already appeared in the early papers on  $\mathcal{H}$ -matrices (cf. [15, 16]).

**4.1. Truncated addition.** In order to make the sum of two  $\mathcal{H}(P, k)$ -matrices be in  $\mathcal{H}(P, k)$ , we have to add them blockwise and truncate each sum  $UV^T$ ,  $U = (U_1, U_2) \in \mathbb{R}^{s \times 2k}$ ,  $V = (V_1, V_2) \in \mathbb{R}^{t \times 2k}$ , of two rank  $k$  blocks  $U_1V_1^T$  and  $U_2V_2^T$  to a matrix of rank at most  $k$ . For this purpose we have to assume that for a given precision  $\varepsilon > 0$ , a matrix  $R$  of rank  $\ell \leq k$  exists such that  $\|UV^T - R\|_2 < \varepsilon$ . The matrix  $R$  can be found by the following algorithm, which was also used in [2] for finding the approximant of lowest rank in an  $\varepsilon$ -neighborhood of a low-rank matrix.

```

procedure truncate( $U, V, k, \text{var } \tilde{U}, \text{var } \tilde{V}$ )
begin
  Compute the  $QR$ -decompositions  $U = Q_U R_U$  and  $V = Q_V R_V$ .
  Compute  $M := R_U R_V^T \in \mathbb{R}^{2k \times 2k}$ .
  Compute the singular value decomposition  $M = XSY^T$ .
  Find the smallest  $\ell$  such that  $s_{\ell+1} \leq \varepsilon s_1$ , where  $s_1 \geq \dots \geq s_{2k}$  are the
    diagonal entries of  $S$ .
  Let  $S_\ell$  and  $Y_\ell$  be the first  $\ell$  columns of  $S$  and  $Y$ , respectively.
  Compute  $\tilde{U} := Q_U X S_\ell$  and  $\tilde{V} := Q_V Y_\ell$ .
end

```

Obviously,  $\tilde{U}\tilde{V}^T$  has rank  $\ell$  and for the error in spectral norm it holds that

$$\|UV^T - \tilde{U}\tilde{V}^T\|_2 = \frac{s_{\ell+1}}{s_1} \|UV^T\|_2 \leq \varepsilon \|UV^T\|_2.$$

The actual rank of a block within an  $\mathcal{H}$ -matrix may therefore be less than  $k$ .

The truncated addition will be denoted by  $\oplus_\varepsilon$  and we define the addition of two submatrices  $A, B$  in the entries  $\hat{b} \in T_{I \times I}$  by

$$A \oplus B = \{A_b \oplus_\varepsilon B_b \text{ for all } b \in P, b \text{ is a descendant of } \hat{b} \text{ in } T_{I \times I}\}.$$

The previous truncation algorithm needs  $\mathcal{O}(k^2(|s| + |t|))$  operations if  $b = s \times t$ . Hence, exploiting the block hierarchy the complexity for the  $\mathcal{H}$ -matrix addition of two matrices from  $\mathcal{H}(P, k)$  can be shown to be of order  $\eta^{-n} k^2 N \log N$ .

**4.2. Truncated matrix-matrix multiplication.** Since the partition  $P$  consists of the leaves of the block cluster tree  $T_{I \times I}$ , we are able to recursively define a modified matrix-matrix multiplication  $C \stackrel{\oplus}{=} A \odot B$ ,  $A \in \mathcal{H}(P, k)$ ,  $B \in \mathcal{H}(P, k)$ ,

making use of the partitioned matrix-matrix multiplication. Let  $r \times s, s \times t, r \times t \in T_{I \times I}$  be block clusters. In order to define what is meant with  $C_{r \times t} \stackrel{\oplus}{=} A_{r \times s} \odot B_{s \times t}$  we have to distinguish three cases.

- (1) All three blocks  $r \times s, s \times t$ , and  $r \times t$  have sons in the tree  $T_{I \times I}$ .

$$\begin{bmatrix} C_{r_1 \times t_1} & C_{r_1 \times t_2} \\ C_{r_2 \times t_1} & C_{r_2 \times t_2} \end{bmatrix} \stackrel{\oplus}{=} \begin{bmatrix} A_{r_1 \times s_1} & A_{r_1 \times s_2} \\ A_{r_2 \times s_1} & A_{r_2 \times s_2} \end{bmatrix} \odot \begin{bmatrix} B_{s_1 \times t_1} & B_{s_1 \times t_2} \\ B_{s_2 \times t_1} & B_{s_2 \times t_2} \end{bmatrix}$$

is recursively defined by

$$\begin{aligned} C_{r_1 \times t_1} &\stackrel{\oplus}{=} A_{r_1 \times s_1} \odot B_{s_1 \times t_1}, & C_{r_1 \times t_1} &\stackrel{\oplus}{=} A_{r_1 \times s_2} \odot B_{s_2 \times t_1}, \\ C_{r_1 \times t_2} &\stackrel{\oplus}{=} A_{r_1 \times s_1} \odot B_{s_1 \times t_2}, & C_{r_1 \times t_2} &\stackrel{\oplus}{=} A_{r_1 \times s_2} \odot B_{s_2 \times t_2}, \\ C_{r_2 \times t_1} &\stackrel{\oplus}{=} A_{r_2 \times s_1} \odot B_{s_1 \times t_1}, & C_{r_2 \times t_1} &\stackrel{\oplus}{=} A_{r_2 \times s_2} \odot B_{s_2 \times t_1}, \\ C_{r_2 \times t_2} &\stackrel{\oplus}{=} A_{r_2 \times s_1} \odot B_{s_1 \times t_2}, & C_{r_2 \times t_2} &\stackrel{\oplus}{=} A_{r_2 \times s_2} \odot B_{s_2 \times t_2}. \end{aligned}$$

- (2) One of the blocks  $r \times s$  and  $s \times t$  is a leaf in  $T_{I \times I}$ .

Assume that  $s \times t$  is a leaf, then  $B_{s \times t}$  has a representation  $B_{s \times t} = U_B V_B^T$ ,  $U_B \in \mathbb{R}^{s \times k}$ ,  $V_B \in \mathbb{R}^{t \times k}$ .

$$C_{r \times t} \stackrel{\oplus}{=} A_{r \times s} \odot B_{s \times t}$$

in this case is defined as

$$C_{r \times t} := C_{r \times t} \oplus A_{r \times s} U_B V_B^T,$$

where  $A_{r \times s} U_B$  are  $k$   $\mathcal{H}$ -matrix-vector products.

- (3)  $r \times t$  has no sons in  $T_{I \times I}$ , and  $r \times s, s \times t$  have sons in the tree  $T_{I \times I}$ .  
For the definition of

$$C_{r \times t} \stackrel{\oplus}{=} \begin{bmatrix} A_{r_1 \times s_1} & A_{r_1 \times s_2} \\ A_{r_2 \times s_1} & A_{r_2 \times s_2} \end{bmatrix} \odot \begin{bmatrix} B_{s_1 \times t_1} & B_{s_1 \times t_2} \\ B_{s_2 \times t_1} & B_{s_2 \times t_2} \end{bmatrix},$$

we introduce matrices  $R_1, R_2, R_3$ , and  $R_4$  by

$$R_1 = R_2 = R_3 = R_4 = 0,$$

and

$$\begin{aligned} R_1 &\stackrel{\oplus}{=} A_{r_1 \times s_1} \odot B_{s_1 \times t_1}, & R_1 &\stackrel{\oplus}{=} A_{r_1 \times s_2} \odot B_{s_2 \times t_1}, \\ R_2 &\stackrel{\oplus}{=} A_{r_1 \times s_1} \odot B_{s_1 \times t_2}, & R_2 &\stackrel{\oplus}{=} A_{r_1 \times s_2} \odot B_{s_2 \times t_2}, \\ R_3 &\stackrel{\oplus}{=} A_{r_2 \times s_1} \odot B_{s_1 \times t_1}, & R_3 &\stackrel{\oplus}{=} A_{r_2 \times s_2} \odot B_{s_2 \times t_1}, \\ R_4 &\stackrel{\oplus}{=} A_{r_2 \times s_1} \odot B_{s_1 \times t_2}, & R_4 &\stackrel{\oplus}{=} A_{r_2 \times s_2} \odot B_{s_2 \times t_2}, \end{aligned}$$

and set

$$C_{r \times t} := C_{r \times t} \oplus_k \left[ \left( \begin{bmatrix} R_1 & 0 \\ 0 & 0 \end{bmatrix} \oplus_k \begin{bmatrix} 0 & R_2 \\ 0 & 0 \end{bmatrix} \right) \oplus_k \left( \begin{bmatrix} 0 & 0 \\ R_3 & 0 \end{bmatrix} \oplus_k \begin{bmatrix} 0 & 0 \\ 0 & R_4 \end{bmatrix} \right) \right].$$

If the truncation accuracy  $\varepsilon$  was chosen to be the machine precision, then  $C \stackrel{\oplus}{=} A \odot B$  would coincide with  $C := C + AB$ . The complexity of the truncated matrix-matrix multiplication can be estimated as  $\mathcal{O}(\eta^{-n} k^2 N \log^2 N)$  (cf. [11]).

**4.3. Inversion.** We assume that each block  $A_{s \times s}$ ,  $s \in T_I$ , of  $A \in \mathcal{H}(P, k)$  is invertible. This is, for example, the case if  $A$  is positive definite. The matrix block  $A_{s \times s}$  corresponding to  $s \in T_I \setminus \mathcal{L}(T_I)$  is subdivided into the sons of  $s \times s$ :

$$A_{s \times s} = \begin{bmatrix} A_{s_1 \times s_1} & A_{s_1 \times s_2} \\ A_{s_2 \times s_1} & A_{s_2 \times s_2} \end{bmatrix}.$$

According to the Frobenius formulas for the inverse of  $A$ , it holds that

$$A_{s \times s}^{-1} = \begin{bmatrix} A_{s_1 \times s_1}^{-1} + A_{s_1 \times s_1}^{-1} A_{s_1 \times s_2} S^{-1} A_{s_2 \times s_1} A_{s_1 \times s_1}^{-1} & -A_{s_1 \times s_1}^{-1} A_{s_1 \times s_2} S^{-1} \\ -S^{-1} A_{s_2 \times s_1} A_{s_1 \times s_1}^{-1} & S^{-1} \end{bmatrix},$$

where  $S$  is the Schur complement  $S = A_{s_2 \times s_2} - A_{s_2 \times s_1} A_{s_1 \times s_1}^{-1} A_{s_1 \times s_2}$ . The  $\mathcal{H}$ -matrix inverse  $C_{s \times s}$  of  $A_{s \times s}$  is defined by replacing the matrix-matrix multiplication and the addition by the  $\mathcal{H}$ -versions. We need a temporary matrix  $T \in \mathcal{H}(P, k)$ , which together with  $C$  is initialized to zero.

```

procedure invertH( $s, A, \text{var } C$ )
begin
  if  $s \in \mathcal{L}(T_I)$  then  $C_{s \times s} := A_{s \times s}^{-1}$  is the usual inverse.
  else begin
    invertH( $s_1, A, C$ ).
     $T_{s_1 \times s_2} \stackrel{\oplus}{=} C_{s_1 \times s_1} \odot A_{s_1 \times s_2}$ .
     $T_{s_2 \times s_1} \stackrel{\oplus}{=} A_{s_2 \times s_1} \odot C_{s_1 \times s_1}$ .
     $A_{s_2 \times s_2} \stackrel{\oplus}{=} A_{s_2 \times s_1} \odot T_{s_1 \times s_2}$ .
    invertH( $s_2, A, C$ ).
     $C_{s_1 \times s_2} \stackrel{\oplus}{=} T_{s_1 \times s_2} \odot C_{s_2 \times s_2}$ .
     $C_{s_2 \times s_1} \stackrel{\oplus}{=} C_{s_2 \times s_2} \odot T_{s_2 \times s_1}$ .
     $C_{s_1 \times s_1} \stackrel{\oplus}{=} T_{s_1 \times s_2} \odot C_{s_2 \times s_1}$ .
  end
end

```

The matrix  $A$  is destroyed during the previous algorithm, and  $C \in \mathcal{H}(P, k)$  contains an approximant of  $A^{-1}$ . The cost for the computation of the  $\mathcal{H}$ -inverse is mainly determined by the cost for the  $\mathcal{H}$ -multiplication. Therefore, an approximation to the inverse of  $A$  can be obtained with complexity  $\mathcal{O}(\eta^{-n} k^2 N \log^2 N)$ .

## 5. NUMERICAL EXPERIMENTS

In this section the practical influence of the various terms of the differential operator (6) on the efficiency and accuracy of the  $\mathcal{H}$ -inverse is investigated. For simplicity all tests are performed on a uniform triangulation of the unit square  $\Omega := (0, 1)^2$  in  $\mathbb{R}^2$ . In each case the stiffness matrix  $S$  is built in the  $\mathcal{H}$ -matrix format; see Example 2.4. Then the inversion algorithm from subsection 4.3 is applied to it with a relative truncation accuracy  $\varepsilon$ . Hence, rank  $k$  is adaptively chosen and is therefore expected to vary among the blocks. All tests were carried out on a single processor of a SunFire 6800 – 900MHz.<sup>1</sup>

<sup>1</sup>The  $\mathcal{H}$ -matrix library which was used for the tests is available under <http://www.mathematik.uni-leipzig.de/~bebendorf/AHMED.html>.

Let  $u_h \in \mathbb{R}^N$  be the finite element solution; i.e., the solution of  $Su_h = b$ , where  $b$  is the vector with the components

$$b_i = \int_{\Omega} f \varphi_i \, dx, \quad i = 1, \dots, N,$$

and  $\tilde{u}_h = Cb$ , where  $C$  is the computed  $\mathcal{H}$ -matrix approximant of  $S^{-1}$ . Since

$$(32) \quad \|u_h - \tilde{u}_h\|_2 = \|u_h - CSu_h\|_2 \leq \|I_N - CS\|_2 \|u_h\|_2,$$

the expression  $\|I_N - CS\|_2$  is an upper bound on the relative accuracy of  $\tilde{u}_h$  compared with the finite element solution  $u_h$ . Note that  $\tilde{u}_h$  cannot be a better approximation of  $u$  than  $u_h$  is, since the proposed method is built on top of the finite element method. Hence, in the following computations we will rely on the expression  $\|I_N - CS\|_2$  as a measure of accuracy.

**5.1. Principal parts.** In the first example we consider operators  $L = -\operatorname{div} A(x) \nabla$ . The coefficients  $A$  of  $L$  are chosen to be of the form

$$A(x) = \begin{bmatrix} 1 & 0 \\ 0 & \alpha(x) \end{bmatrix}, \quad x \in \Omega,$$

where  $\alpha(x) = 1$  in the lower region of Figure 1 and a random number from the interval  $[0, a]$  in the remaining part of the unit square. In order to avoid averaging effects, the coefficient  $\alpha$  possesses a two-level random structure: the randomly chosen coefficient on each triangle is multiplied by a coefficient chosen randomly on a scale of length  $\sqrt{h}$ , where the grid size  $h$  is defined through  $h(\sqrt{N/2} + 1) = 1$ .

In Table 1 the accuracy  $\|I_N - CS\|_2$  of the  $\mathcal{H}$ -matrix  $C$  and the CPU time consumption are compared for different  $a$  and different problem sizes  $N$ . The truncation accuracy  $\varepsilon$  is chosen such that  $\|I_N - CS\|_2$  is of order  $h$ .

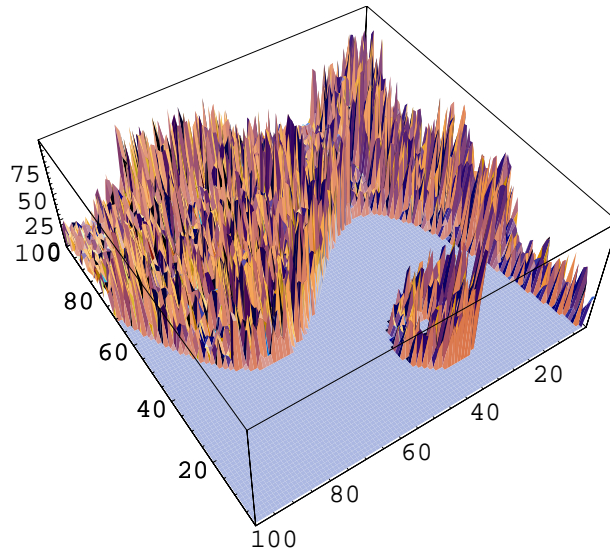


FIGURE 1. The coefficient  $\alpha(x)$



TABLE 1.

$N$	$h$	$\varepsilon$	$a = 1.0$		$a = 10.0$		$a = 100.0$	
			time [s]	accuracy	time [s]	accuracy	time [s]	accuracy
14400	$8.3e-3$	$1e-5$	29.4	$8.5e-3$	28.4	$1.1e-2$	29.0	$1.1e-1$
38025	$5.1e-3$	$2e-6$	146.7	$6.2e-3$	138.6	$6.9e-3$	143.2	$1.5e-2$
65025	$3.9e-3$	$5e-7$	341.3	$3.4e-3$	327.6	$3.2e-3$	338.1	$6.0e-3$
129600	$2.8e-3$	$2e-7$	1074.9	$2.6e-3$	1008.0	$2.8e-3$	1032.3	$4.7e-3$
278784	$1.9e-3$	$5e-8$	3452.1	$1.6e-3$	3201.7	$2.0e-3$	3320.1	$4.4e-3$
529984	$1.3e-3$	$2e-8$	9210.0	$1.3e-3$	8529.0	$1.5e-3$	8868.7	$3.1e-3$

TABLE 2.

	$a = 1.0$		$a = 10.0$		$a = 100.0$	
	time [s]	accuracy	time [s]	accuracy	time [s]	accuracy
$b = 1$	333.4	$3.2e-3$	334.4	$5.0e-3$	385.4	$8.0e-3$
$b = 10$	327.6	$2.8e-3$	326.8	$3.0e-3$	373.7	$2.7e-3$
$b = 100$	329.7	$3.1e-3$	330.2	$2.6e-3$	394.3	$2.0e-3$
$b = 1000$	331.9	$2.7e-3$	333.1	$3.5e-3$	388.2	$1.8e-3$

In the next set of tests the same quantities for a smooth but oscillating coefficient in the principal part are computed; i.e.,  $\alpha$  is chosen to be the function

$$\alpha(x) = a(1 + \cos(b2\pi x) \sin(b2\pi y)).$$

By changing the coefficient  $a$  we are able to prescribe the amplitude of the oscillation and by  $b$  the number of oscillations in  $x$ - and  $y$ -direction of  $\Omega = (0, 1)^2$ . Table 2 contains the results for  $N = 65025$  and  $\varepsilon = 5e-7$ .

The experiments show that in the absence of lower-order terms neither the accuracy nor the CPU times needed to compute the approximant do depend much on the coefficients.

**5.2. Convection-diffusion.** In this section operators of the type

$$L = -\Delta + \mathbf{c} \cdot \nabla$$

will be considered. In the first example the convection coefficient  $\mathbf{c}$  is randomly chosen; i.e.,  $\mathbf{c}(x) \in [-a, a]^2$  for  $x \in \Omega$ . Table 3 shows  $\|I_N - CS\|_2$  for different parameters  $a$ .

In the next example we investigate operators

$$Lu = -\epsilon \Delta u + u_x + u_y$$

for different parameters  $\epsilon > 0$ . We are particularly interested in a convection dominated setting. Table 4 shows the results.

TABLE 3.

$N$	$h$	$\varepsilon$	$a = 1.0$		$a = 10.0$		$a = 100.0$	
			time [s]	accuracy	time [s]	accuracy	time [s]	accuracy
14641	$8.2e-3$	$1e-5$	31.4	$7.1e-3$	31.5	$7.5e-3$	31.2	$9.5e-3$
38416	$5.1e-3$	$2e-6$	153.6	$6.5e-3$	154.1	$6.6e-3$	155.7	$8.0e-3$
65025	$3.9e-3$	$5e-7$	355.9	$2.9e-3$	356.9	$3.0e-3$	356.2	$3.5e-3$
129600	$2.8e-3$	$2e-7$	1100.3	$2.5e-3$	1104.3	$2.5e-3$	1106.4	$2.6e-3$
278784	$1.9e-3$	$5e-8$	3503.2	$1.5e-3$	3505.8	$1.6e-3$	3518.7	$1.5e-3$
597529	$1.3e-3$	$2e-8$	11072.6	$1.4e-3$	11111.1	$1.5e-3$	11105.7	$2.0e-3$

TABLE 4.

$N$	$\epsilon = 0.1$		$\epsilon = 0.01$		$\epsilon = 0.001$	
	time [s]	accuracy	time [s]	accuracy	time [s]	accuracy
14641	30.0	$5.7e-3$	29.4	$9.4e-4$	60.9	$2.7e-4$
38416	145.5	$4.2e-3$	154.9	$4.6e-4$	259.3	$1.2e-4$
65025	336.8	$1.9e-3$	362.3	$2.6e-4$	481.0	$2.5e-5$
129600	1046.5	$1.6e-3$	1133.0	$2.0e-4$	1233.0	$1.8e-5$
278784	3311.4	$9.4e-4$	3665.6	$1.1e-4$	2998.5	$1.3e-5$
597529	10452.4	$1.0e-3$	11701.8	$1.2e-4$	11370.6	$9.5e-6$

Since the tables above show that it is possible to find an  $\mathcal{H}$ -matrix  $C$  that approximates  $S^{-1}$ , from (32) it is obvious that  $\tilde{u}_h$  approximates the finite element solution  $u_h$ . This is especially true in the presence of boundary layers, as illustrated in the following example. The solution of

$$-\epsilon \Delta u + u_x + u_y = f,$$

where

$$f(x, y) = (x + y)(1 - e^{(x-1)/\epsilon} e^{(y-1)/\epsilon}) + (x - y)(e^{(y-1)/\epsilon} - e^{(x-1)/\epsilon})$$

with zero boundary conditions is known to be

$$u(x, y) = xy(1 - e^{(x-1)/\epsilon})(1 - e^{(y-1)/\epsilon}).$$

Figure 2 compares the restrictions of  $u$  and  $\tilde{u}_h$  to the set  $\{(x, x), x \in (0, 1)\}$  for  $\epsilon = 0.01$  and  $N = 14641$ . Obviously, the proposed inversion procedure is able to handle boundary layers as long as the underlying finite element method is stable.

In the next example we consider convection in a direction that is aligned with the grid; i.e., operators

$$Lu = -\epsilon \Delta u + u_x$$

for different parameters  $\epsilon > 0$  are investigated. The CPU times for computing the approximants and their accuracies can be found in Table 5.

It is well known that if the ratio  $\epsilon/h$  gets small, the finite element discretization suffers from the loss of stability. As a consequence the stiffness matrix becomes ill conditioned. It may even happen that  $S$  is not invertible. This behavior is observable for small  $N$  if  $\epsilon$  tends to zero: when changing  $\epsilon = 0.01$  to  $\epsilon = 0.001$

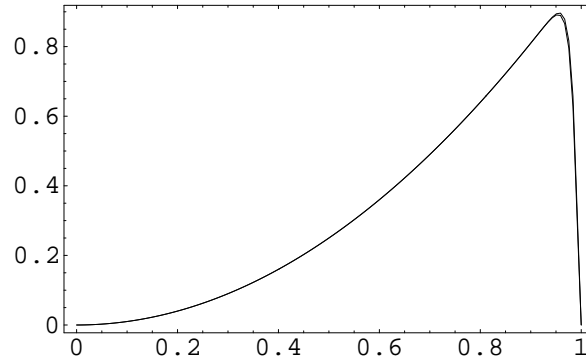
FIGURE 2. The solutions  $u$  and  $\tilde{u}_h$ .

TABLE 5.

$N$	$\epsilon = 0.1$		$\epsilon = 0.01$		$\epsilon = 0.001$	
	time [s]	accuracy	time [s]	accuracy	time [s]	accuracy
14641	30.3	$6.3e-3$	31.4	$6.6e-4$	47.4	$1.3e-4$
38416	145.7	$4.9e-3$	161.4	$5.7e-4$	226.1	$6.3e-5$
65025	339.0	$2.2e-3$	365.7	$1.8e-4$	463.9	$1.8e-5$
129600	1053.9	$1.9e-3$	1152.1	$1.8e-4$	1235.8	$1.4e-5$
278784	3323.3	$1.1e-3$	3657.8	$1.5e-4$	4135.5	$6.3e-6$
597529	10457.6	$1.0e-3$	11646.8	$1.3e-4$	12208.7	$7.9e-6$

the CPU time for  $N = 14641$  increases by a factor of 1.5, while it is almost not influenced in the case  $N = 597529$ , where the discretization is stable.

The proposed inversion procedure can also be applied to Shishkin meshes, which should be used for a better convergence of the finite element method.

**5.3. Diffusion-reaction.** As a third kind of example we consider the operator  $Lu = -\Delta u + du$  with a randomly chosen reaction term  $d(x) \in [0, a]$  for  $x \in \Omega$ . Since by adding a positive  $d$  to the operator  $-\Delta$ , the distance of the spectrum to zero is increased. Hence, the larger the  $d$ , the better the approximation works. The respective results can be found in Table 6.

Negative  $d$  (i.e., the Helmholtz equation) can also be handled as long as the inverse of  $L$  is guaranteed to exist. If an integral formulation of the Helmholtz equation is discretized, it is known that, due to the oscillatory kernel, the rank of the blocks in the  $\mathcal{H}$ -matrix approximant depends on the wave number. Therefore, the method can still be applied, but it is only efficient if the wave number is not too large. The same effect can be observed here.

Each column of Table 7 shows the approximation results for the respective  $d$  in the case  $N = 129600$ . In order to be able to guarantee  $\|I_N - CS\|_2 \sim h$ , we have to choose a higher truncation accuracy  $\varepsilon$  if the modulus of  $d$  is increased. Note that  $d$  is now a constant. Hence, for large wave numbers the inversion procedure can be applied, but gets less efficient.

As a last example we consider values  $d$  that are close to the eigenvalues of the operator  $\Delta$ . In the case of resonance (i.e.,  $-d$  is an eigenvalue of  $S$ ), the stiffness matrix  $S$  is not invertible. Close to resonance the stiffness matrix is ill conditioned. The functions

$$u_k(x, y) := \sin(2\pi kx) \sin(2\pi ky), \quad k = 0, 1, \dots,$$

TABLE 6.

$N$	$a = 10.0$		$a = 100.0$		$a = 1000.0$	
	time [s]	accuracy	time [s]	accuracy	time [s]	accuracy
14641	31.2	$5.5e-3$	32.3	$1.8e-3$	34.7	$2.1e-4$
38416	154.4	$4.9e-3$	161.6	$1.3e-3$	172.2	$7.2e-5$
65025	357.1	$2.1e-3$	369.1	$6.1e-3$	394.2	$4.6e-5$
129600	1106.5	$1.9e-3$	1133.9	$5.8e-4$	1222.8	$5.4e-5$
278784	3502.8	$1.2e-3$	3610.1	$3.7e-4$	3876.2	$3.6e-5$
597529	11100.7	$1.1e-3$	11500.8	$3.6e-4$	12242.4	$3.2e-5$

TABLE 7.

	$d = -1.0$	$d = -10.0$	$d = -100.0$	$d = -1000.0$	$d = -10000.0$
$\varepsilon$	$2.0e-7$	$1.0e-7$	$2.0e-9$	$5.0e-10$	$2.0e-11$
accuracy	$2.5e-3$	$2.8e-3$	$2.5e-3$	$3.8e-3$	$2.9e-3$
time [s]	1029.1	1076.7	1386.6	1678.3	3144.2

TABLE 8.

$d$	-39.0	-39.5	-40.0	-40.15	-40.2	-40.25	-41.0
$\varepsilon$	$2.0e-8$	$2.0e-8$	$1.0e-8$	$2.0e-9$	$5.0e-10$	$2.0e-9$	$2.0e-8$
accuracy	$1.9e-3$	$1.7e-3$	$2.0e-3$	$1.6e-3$	$2.3e-3$	$1.3e-3$	$1.5e-3$
time [s]	1109.4	1114.6	1168.2	1299.4	1414.0	1292.6	1109.1

solve the eigenproblem

$$\begin{aligned} -\Delta u &= \lambda u \quad \text{in } \Omega = (0, 1)^2, \\ u &= 0 \quad \text{on } \partial\Omega, \end{aligned}$$

with corresponding eigenvalues  $\lambda_k := (2\pi k)^2$ . In Table 8  $N = 129600$  and  $-d$  is chosen in a neighborhood of the second eigenvalue  $\lambda_1 = 4\pi^2$ .

From the numerical experiments above we conclude that the proposed method is robust with respect to nonsmooth and anisotropic coefficients. Even convection-diffusion problems with dominating convection can be solved efficiently without special adaptation of the algorithm to this class of problems. Hence, the proposed inversion procedure can be applied whenever a stable discretization of the operator  $L$  is available. Helmholtz' equation can be treated but the algorithms are only efficient for relatively small wave numbers.

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#### REFERENCES

- [1] M. Bebendorf: *Approximation of boundary element matrices*. Numer. Math. **86**, 565–589, 2000. MR2001j:65022
- [2] M. Bebendorf: *Effiziente numerische Lösung von Randintegralgleichungen unter Verwendung von Niedrigrang-Matrizen*. dissertation.de, Verlag im Internet, 2001. ISBN 3-89825-183-7.
- [3] M. Bebendorf and S. Rjasanow: *Adaptive Low-Rank Approximation of Collocation Matrices*. Computing **70**(1), 1–24, 2003. MR2004a:65177
- [4] M. Bebendorf and W. Hackbusch: *Existence of  $\mathcal{H}$ -Matrix Approximants to the Inverse FE-Matrix of Elliptic Operators with  $L^\infty$ -Coefficients*. Numer. Math. **95**, 1–28, 2003. MR2004e:65128
- [5] M. Bebendorf: *A Note on the Poincaré-Inequality for Convex Domains*. J. Anal. Appl. **22**, 751–756, 2003. MR2004e:65128
- [6] G. Beylkin, R. Coifman, and V. Rokhlin: *Fast wavelet transforms and numerical algorithms*. I. Comm. Pure Appl. Math. **44**(2), 141–183, 1991.
- [7] W. Dahmen, S. Prössdorf, and R. Schneider: *Wavelet approximation methods for pseudodifferential equations. II. Matrix compression and fast solution*. Adv. Comput. Math. **1**(3-4), 259–335, 1993. MR95g:65149
- [8] W. Dahmen, S. Prössdorf, and R. Schneider: *Wavelet approximation methods for pseudodifferential equations. I. Stability and convergence*. Math. Z. **215**(4), 583–620, 1994. MR95g:65148

- [9] G. Dolzmann and S. Müller: *Estimates for Green's matrices of elliptic systems by  $L^p$  theory*. Manuscripta Math. **88**, 261–273, 1995. MRMR96g:35054
- [10] D. Gilbarg and N. S. Trudinger: *Elliptic partial differential equations of second order*, Reprint of the 1998 edition, Springer, Berlin, 2001. MR2001k:35004
- [11] L. Grasedyck: *Theorie und Anwendungen Hierarchischer Matrizen*. Dissertation, Universität Kiel, 2001. MR
- [12] M. Grüter and K.-O. Widman: *The Green function for uniformly elliptic equations*. Manuscripta Math. **37**, 303–342, 1982. MR83h:35033
- [13] L. Greengard and V. Rokhlin: *A new version of the fast multipole method for the Laplace equation in three dimensions*. Acta Numerica, 1997, pages 229–269. Cambridge Univ. Press, Cambridge, 1997. MR99c:65012
- [14] W. Hackbusch: *Theorie und Numerik elliptischer Differentialgleichungen*. B. G. Teubner, Stuttgart, 1996 - English translation: *Elliptic differential equations. Theory and numerical treatment*. Springer-Verlag, Berlin, 1992. MR94b:35001
- [15] W. Hackbusch: *A sparse matrix arithmetic based on  $\mathcal{H}$ -matrices. I. Introduction to  $\mathcal{H}$ -matrices*. Computing **62**, 89–108, 1999. MR2000c:65039
- [16] W. Hackbusch and B. N. Khoromskij: *A sparse  $\mathcal{H}$ -matrix arithmetic. II. Application to multi-dimensional problems*. Computing **64**, 21–47, 2000. MR2001i:65053
- [17] W. Hackbusch and Z. P. Nowak: *On the fast matrix multiplication in the boundary element method by panel clustering*. Numer. Math. **54**, 463–491, 1989. MR89k:65162
- [18] S. Prössdorf and B. Silbermann: *Numerical analysis for integral and related operator equations*, Akademie Verlag, Berlin, 1991. MR94f:65126a
- [19] E. Tyrtysnikov: *Mosaic-skeleton approximations*. Calcolo **33**(1-2), 47–57 (1998), 1996. MR99f:15005

FAKULTÄT FÜR MATHEMATIK UND INFORMATIK, UNIVERSITÄT LEIPZIG, AUGUSTUSPLATZ 10/11,  
D-04109 LEIPZIG, GERMANY

*E-mail address:* bebendorf@math.uni-leipzig.de