

# Additive Schwarz Algorithms for Parabolic Convection-Diffusion Equations <sup>\*</sup>

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## Abstract

In this paper, we consider the solution of linear systems of algebraic equations that arise from parabolic finite element problems. We introduce three additive Schwarz type domain decomposition methods for general, not necessarily selfadjoint, linear, second order, parabolic partial differential equations and also study the convergence rates of these algorithms. The resulting preconditioned linear system of equations is solved by the generalized minimal residual method. Numerical results are also reported.

**Key words** Schwarz's alternating method, domain decomposition, parabolic convection-diffusion equation, finite elements.

**AMS(MOS) subject classifications.** 65N30, 65F10

## 1 Introduction

Domain decomposition techniques are powerful iterative methods for solving linear systems of equations that arise from discretizing partial differential equations. A systematic theory has been developed for elliptic finite element problems in the past few years, see [2, 3, 4, 5, 6, 7, 8, 9, 12, 13, 21] etc. In this paper, we are interested in solving the finite element parabolic convection-diffusion problems, obtained by using implicit schemes, such as the backward Euler scheme and the Crank-Nicolson scheme, in the time variable, and Galerkin approximation in the space variables. At a fixed time level, the resulting equation is equivalent to an elliptic problem which depends on a time step parameter. This suggests that we might apply the methods, originally proposed for elliptic equations, to the parabolic cases. The central mathematical question is then to estimate how the convergence rate depends on the space mesh and the time step parameters, especially in the case when the latter parameter is relatively large. Other domain decomposition methods for parabolic problems can be found in [11, 17, 20] and references therein.

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The outline of this paper is as follows. In the remainder of this section, we present the continuous and the discrete parabolic convection-diffusion equation, and also discuss some of their basic properties. In section 2, we outline an abstract additive Schwarz theory; cf. [8], without providing any detailed analysis. In section 3, we introduce three additive Schwarz type algorithms, which are instances of the abstract theory. The first one is an straightforward generalization of the additive Schwarz method for elliptic equation [9, 13], the second one is obtained by eliminating the coarse mesh space and using only local function spaces to take the advantage of the existence of a time step parameter. This algorithm is therefore more suitable for parallel computers than the first one. The third one is a non-overlapping domain decomposition method, but nevertheless fits into our general theory of additive Schwarz methods. Finally, in section 4, we report on some of our numerical experiments.

Throughout this paper,  $c$  and  $C$ , with or without subscripts, denote generic, strictly positive constants. Their values may be different at different occurrences, but are independent of the mesh parameters  $H$ ,  $h$  and  $\tau$ , which will be introduced later.

### 1.1 A parabolic convection-diffusion problem

We consider the following parabolic convection-diffusion problem: Find  $u(x, t)$ , such that

$$\begin{cases} \partial u / \partial t + Lu = f & \text{in } \Omega \times [0, T], \\ u(x, t) = 0 & \text{on } \partial\Omega \times [0, T], \\ u(x, 0) = u_0(x) & \text{in } \Omega, \end{cases} \quad (1)$$

where  $\Omega \subset \mathbb{R}^d$ ,  $d = 2$  or  $3$ , is a polygonal domain with boundary  $\partial\Omega$ .  $L$  is a strongly elliptic operator, which has the following form

$$Lu = - \sum_{i,j=1}^d \frac{\partial}{\partial x_i} (a_{ij}(x) \frac{\partial u}{\partial x_j}) + \sum_{i=1}^d b_i(x) \frac{\partial u}{\partial x_i} + c(x)u.$$

All coefficients are sufficiently smooth and  $a_{ij}(x) = a_{ji}(x)$  for all  $i, j$  and  $x$ .

A weak formulation of equation (1) is: Find  $u(x, t) \in H_0^1(\Omega)$ ,  $u(x, 0) = u_0(x)$  in  $\Omega$ , such that

$$\left( \frac{\partial u}{\partial t}, v \right) + B(u, v) = (f, v), \quad \forall v \in H_0^1(\Omega), \quad \forall t \in [0, T].$$

The bilinear form  $B(\cdot, \cdot)$  is defined as

$$B(u, v) = \sum_{i,j=1}^d \int_{\Omega} a_{ij} \frac{\partial u}{\partial x_i} \frac{\partial v}{\partial x_j} dx + \sum_{i=1}^d \int_{\Omega} b_i \frac{\partial u}{\partial x_i} v dx + \int_{\Omega} c u v dx$$

and the linear functional

$$(f, v) = \int_{\Omega} f v dx.$$

We assume that the bilinear form is

- (1) bounded, i.e.  $|B(u, v)| \leq C \|u\|_{H_0^1(\Omega)} \|v\|_{H_0^1(\Omega)}$ ,  $\forall u, v \in H_0^1(\Omega)$ , and
- (2) elliptic, i.e.  $|B(u, u)| \geq c \|u\|_{H_0^1(\Omega)}^2$ ,  $\forall u \in H_0^1(\Omega)$ .

The existence and uniqueness of the solution of the weak parabolic convection-diffusion equation present no problems, see [19]. We use two types of time discretization, namely, a backward Euler scheme and a Crank-Nicolson scheme. Let  $\Delta t_n$  be the  $n^{\text{th}}$  time step,  $M$  the number of steps and  $\sum_{n=1}^M \Delta t_n = T$ . For the first scheme, the time discrete problem is

$$\left( \frac{u^n - u^{n-1}}{\Delta t_n}, v \right) + B(u^n, v) = (f, v), \quad \forall v \in H_0^1(\Omega),$$

with  $u^0(x, t) = u_0(x)$  and  $n = 1, \dots, M$ . For the second scheme, we have

$$\left( \frac{u^n - u^{n-1}}{\Delta t_n}, v \right) + B\left(\frac{u^n + u^{n-1}}{2}, v\right) = (f, v), \quad \forall v \in H_0^1(\Omega),$$

with  $u^0(x, t) = u_0(x)$  and  $n = 1, \dots, M$ . Both schemes lead to the following problem: For a given function  $g_{n-1} \in H^{-1}(\Omega)$ , find  $w \in H_0^1(\Omega)$ , such that

$$D_\tau(w, v) \equiv (w, v) + \tau B(w, v) = (g_{n-1}, v), \quad \forall v \in H_0^1(\Omega), \quad (2)$$

where  $\tau$  is the time step parameter. For the backward Euler scheme

$$\begin{aligned} w &= u^n - u^{n-1}, \\ \tau &= \Delta t_n, \\ (g_{n-1}, v) &= \tau((f, v) - B(u^{n-1}, v)), \end{aligned}$$

and for the Crank-Nicolson scheme,

$$\begin{aligned} w &= u^n - u^{n-1}, \\ \tau &= \Delta t_n / 2, \\ (g_{n-1}, v) &= \tau(2(f, v) - B(u^{n-1}, v)). \end{aligned}$$

The stability of both schemes is well understood, see [19]. In this paper, we focus on the study of fast iterative algorithms for solving the linear systems at each time step. In general,  $D_\tau(\cdot, \cdot)$  is not symmetric. For technical reasons, it is convenient to separate the symmetric and the skewsymmetric parts of the bilinear form. We therefore introduce the bilinear forms  $A_\tau(u, v) = 1/2(D_\tau(u, v) + D_\tau(v, u))$ , which is symmetric, and  $N_\tau(u, v) = 1/2(D_\tau(u, v) - D_\tau(v, u))$ , which is skewsymmetric. It is easy to prove that  $N_\tau(\cdot, \cdot)$  is bounded in the following sense: There exists a constant  $C$ , such that

$$|N_\tau(u, v)| \leq C\tau \|u\|_{H_0^1(\Omega)} \|v\|_{L^2(\Omega)}, \quad \forall u, v \in H_0^1(\Omega).$$

Define the  $\tau$ -norm by  $\|\cdot\|_{\tau(\Omega)}^2 = \|\cdot\|_{L^2(\Omega)}^2 + \tau\|\cdot\|_{H_0^1(\Omega)}^2$ . The following lemma gives the boundness and positive definiteness of  $D_\tau(\cdot, \cdot)$ .

**Lemma 1** *There exist constants  $C$  and  $c > 0$ , which are independent of  $\tau$ , such that*

- (1)  $|D_\tau(u, v)| \leq C \|u\|_{\tau(\Omega)} \|v\|_{\tau(\Omega)}, \quad \forall u, v \in H_0^1(\Omega)$
- (2)  $D_\tau(u, u) \geq c \|u\|_{\tau(\Omega)}^2, \quad \forall u \in H_0^1(\Omega)$
- (3)  $|A_\tau(u, v)| \leq C \|u\|_{\tau(\Omega)} \|v\|_{\tau(\Omega)}, \quad \forall u, v \in H_0^1(\Omega)$
- (4)  $A_\tau(u, u) \geq c \|u\|_{\tau(\Omega)}^2, \quad \forall u \in H_0^1(\Omega).$

The proof follows immediately from the boundedness and ellipticity of  $B(\cdot, \cdot)$ .

If we denote  $\|u\|_{A_\tau} = \sqrt{A_\tau(u, u)}$ , then we know from this lemma that the  $A_\tau$ -norm is equivalent to the  $\tau$ -norm. We shall use the  $A_\tau$ -norm in our discussion. It can be shown, cf Grisvard [18] and Nečas [22], that the solution of  $B(u, v)$  is  $H^{1+\gamma}(\Omega)$ -regular in the following sense. For any  $g \in L^2(\Omega)$ , there exists a solution  $u \in H^{1+\gamma}(\Omega) \cap H_0^1(\Omega)$  of

$$B(v, u) = (g, v), \quad \forall v \in H_0^1(\Omega)$$

with

$$|u|_{H^{1+\gamma}(\Omega)} \leq C \|g\|_{L^2(\Omega)},$$

where  $\gamma \in [1/2, 1]$ . Based on this regularity result, we have

**Lemma 2** *For any  $g \in L^2(\Omega)$ , the equation*

$$D_\tau(v, u) = (g, v), \quad \forall v \in H_0^1(\Omega)$$

*has a solution  $u \in H^{1+\gamma}(\Omega) \cap H_0^1(\Omega)$  that satisfies  $|u|_{H^{1+\gamma}(\Omega)} \leq C/\tau \|g\|_{L^2(\Omega)}$ .*

**Proof:** Consider the equality

$$(u, u) + \tau B(u, u) = (g, u).$$

Since  $B(u, u) \geq 0$ , we have  $\|u\|_{L^2(\Omega)} \leq \|g\|_{L^2(\Omega)}$ , which implies that  $\|u - g\|_{L^2(\Omega)} \leq 2\|g\|_{L^2(\Omega)}$ . On the other hand, we have

$$B(v, u) = ((g - u)/\tau, v), \quad \forall v \in H_0^1(\Omega).$$

By using the regularity result of  $B(\cdot, \cdot)$ , we complete the proof.  $\square$

## 1.2 A finite element approximation

We solve equation (2) by a conformal Galerkin finite element method. For simplicity, we use piecewise linear triangular elements in  $\mathbb{R}^2$  and the corresponding tetrahedral elements in  $\mathbb{R}^3$ . In this subsection, we introduce a two level triangulation of  $\Omega$ , the corresponding finite element spaces and the Galerkin equation.

For a given polygonal region  $\Omega \in \mathbb{R}^d$ , let  $\{\Omega_i\}$  be a shape regular finite element triangulation of  $\Omega$  and  $H$  the maximal diameter of these  $\Omega_i$ 's. We sometimes refer  $\Omega_i$  as a substructure and  $\{\Omega_i\}$  as the coarse mesh or  $H$ -level triangulation of  $\Omega$ .

We further divide each substructure  $\Omega_i$  into smaller triangles, denoted as  $\tau_i^j$ ,  $j = 1, \dots$ . We assume that  $\{\tau_i^j\}$  form a shape regular finite element triangulation of  $\Omega$  and  $h$  is the maximal diameter of  $\tau_i^j$ . We call  $\{\tau_i^j\}$  the fine mesh or  $h$ -level triangulation of  $\Omega$ .

We next define the piecewise linear finite element function spaces over both  $H$ -level and  $h$ -level triangulations of  $\Omega$ .

$$V^H = \{v^H \mid \text{continuous on } \Omega, v^H|_{\Omega_i} \text{ linear on } \Omega_i, \forall i, v^H = 0 \text{ on } \partial\Omega\}$$

and

$$V^h = \{v^h \mid \text{continuous on } \Omega, v^h|_{\tau_i^j} \text{ linear on } \tau_i^j, \forall i, v^h = 0 \text{ on } \partial\Omega\}$$

It is obvious that  $V^H \subset V^h$ . We also use the notations

$$\Lambda^h = \{x \mid x \in \text{interior nodes of } h\text{-level subdivision}\}$$

and

$$\Lambda^H = \{x \mid x \in \text{interior nodes of } H\text{-level subdivision}\}.$$

To obtain a fully discrete problem, we discretize equation (2) in space by using a Galerkin finite element method and the subspace  $V^h \subset H_0^1(\Omega)$ . The Galerkin approximation of equation (2) reads as follows: Find  $\mathbf{u}^h \in V^h$ , such that

$$D_\tau(\mathbf{u}^h, v^h) = (g, v^h), \quad \forall v^h \in V^h. \quad (3)$$

Here  $g$  is different for different time discretizations. The existence and uniqueness of  $\mathbf{u}^h$  has been extensively studied in the literature, see [1, 19]. By using the nodal basis functions, equation (3) can be transformed into a linear system of equations, which is usually large and sparse, but not well conditioned. It is well known that the efficiency of any iterative methods to be used to solve this linear system of equations depends strongly on the conditioning of the stiffness matrix. The main focus of this paper is to provide some effective preconditioning techniques for solving this equation.

## 2 An abstract theory for the additive Schwarz method

In this section, we present an abstract theory developed in [8] for the additive Schwarz method. Interested reader should see [8] for detailed proof. Let  $V$  be a Hilbert space of real functions defined on  $\Omega \subset \mathbb{R}^d$  with an inner product  $(u, v)_V$  and the corresponding norm  $\|u\|_V$ . Let  $\omega$  be a subdomain of  $\Omega$ . We define  $(u, v)_{V(\omega)}$  to be the restriction of the inner product to  $\omega$ .

Let  $\mathcal{B}(\cdot, \cdot)$  be a bilinear form on  $V \times V$  and  $\mathcal{F}(\cdot)$  a linear functional on  $V$  such that

- $\mathcal{B}(\cdot, \cdot)$  is continuous; in particular  $|\mathcal{B}(u, v)| \leq C\|u\|_{V(\omega)}\|v\|_{V(\omega)}, \quad \forall u, v \in V$ ,

where  $\omega = \{\text{supp } u\} \cap \{\text{supp } v\}$ .

- $\mathcal{B}(\cdot, \cdot)$  is  $V$ -elliptic, i.e.  $\mathcal{B}(u, u) \geq c\|u\|_V^2, \quad \forall u \in V$ .

- $\mathcal{F}(\cdot)$  is continuous, i.e.  $|\mathcal{F}(u)| \leq C\|u\|_V, \quad \forall u \in V$ .

We define  $\mathcal{A}(u, v) = 1/2(\mathcal{B}(u, v) + \mathcal{B}(v, u))$ , a symmetric bilinear form, and  $\mathcal{N}(u, v) = 1/2(\mathcal{B}(u, v) - \mathcal{B}(v, u))$ , which is skewsymmetric.

It follows from the assumptions on  $\mathcal{B}(u, v)$  that  $\mathcal{A}(u, v)$  is elliptic and continuous in the same sense as  $\mathcal{B}$ . This implies that the norm corresponding to  $\mathcal{A}(\cdot, \cdot)$  is equivalent to the  $V$ -norm. In the following, we can therefore use  $(\cdot, \cdot)_{\mathcal{A}}$ , instead of  $(\cdot, \cdot)_V$ .

We are interested in solving the problem that reads as follows: Find  $\mathbf{u} \in V$ , such that

$$\mathcal{B}(\mathbf{u}, v) = \mathcal{F}(v), \quad \forall v \in V. \quad (4)$$

Let  $V_i$ ,  $i = 0, \dots, N$ , be subspaces of  $V$ , such that  $V = V_0 + \dots + V_N$ , i.e., for any  $v \in V$ , there exist  $v_i \in V_i$ ,  $i = 0, \dots, N$ , such that  $v = v_0 + \dots + v_N$ . Moreover, we assume that there exists a constant  $C_0^2$  such that

$$\sum_{i=0}^N \|v_i\|_{V(\omega_i)}^2 \leq C_0^2 \|v\|_V^2, \quad \forall v \in V, \quad (5)$$

where  $\omega_i$  is the support of  $V_i$ , i.e. the union of the supports of all functions in  $V_i$ .

This is usually called the bounded decomposition lemma. Note that the constant  $C_0^2$  may depend on the number of subregions  $N$  and also on some other parameters of  $V$ , which may be introduced in particular applications.

We also assume that there exists a constant  $C_\omega$ , such that

$$\sum_{i=0}^N \|v\|_{V(\omega_i)}^2 \leq C_\omega \|v\|_V^2, \quad \forall v \in V. \quad (6)$$

In fact this constant  $C_\omega$  is the maximal number of  $\omega_i$ 's to which any point  $x \in \Omega$  can belong.

For each subspace  $V_i$ ,  $0 \leq i \leq N$ , we define a projection  $P_i^{\mathcal{B}} = P_{V_i}^{\mathcal{B}} : V \rightarrow V_i$ , with respect to the bilinear form  $\mathcal{B}(\cdot, \cdot)$ , as the solution of

$$\mathcal{B}(P_i^{\mathcal{B}} u, v) = \mathcal{B}(u, v), \quad \forall v \in V_i.$$

Let us denote  $P^{\mathcal{B}} = P_0^{\mathcal{B}} + \dots + P_N^{\mathcal{B}} : V \rightarrow V$  and  $\mathbf{b} = P_0^{\mathcal{B}} \mathbf{u} + \dots + P_N^{\mathcal{B}} \mathbf{u}$ , where each component  $P_i^{\mathcal{B}} \mathbf{u}$  can be obtained, without *a priori* knowledge the solution  $\mathbf{u}$ , by solving

$$\mathcal{B}(P_i^{\mathcal{B}} \mathbf{u}, v) = \mathcal{F}(v), \quad \forall v \in V_i.$$

The additive Schwarz algorithm can be stated as: Find the solution  $\mathbf{u}$  of equation (4) by solving the following equation.

$$P^{\mathcal{B}} \mathbf{u} = \mathbf{b}. \quad (7)$$

We often refer to this equation as the derived equation with respect to the bilinear form  $\mathcal{B}(\cdot, \cdot)$  and the decomposition  $\{V_i\}$ , see [8, 13, 14]. The following theorem is easy to establish.

**Theorem 1** *If the operator  $P^{\mathcal{B}}$  is invertible, then the equations (4) and (7) have the same solution.*

In practice, the operator  $P^{\mathcal{B}}$  can be constructed in terms of the original stiffness matrix and the inverse of some small matrices. The inverse of the small matrices are imbedded, by zero, in a larger matrix. The explicit matrix for  $P^{\mathcal{B}}$  is normally not known. However, the matrix-vector-product  $P^{\mathcal{B}} u$ , which is all that is needed,

can be computed by solving a linear system of equations for each subregion. In the following theorem, we summarize the bounds of the operator  $P^{\mathcal{B}}$ , which determine the convergence rate of the iterative method used to solve (7). The result was established in [8].

**Theorem 2** (1) *There exists a constant  $C$ , such that*

$$\|P^{\mathcal{B}}u\|_{\mathcal{A}} \leq CC_{\omega}\|u\|_{\mathcal{A}}, \quad \forall u \in V.$$

(2) *There exists a constant  $c > 0$ , independent of  $C_0$ , such that*

$$\|P^{\mathcal{B}}u\|_{\mathcal{A}} \geq cC_0^{-2}\|u\|_{\mathcal{A}}, \quad \forall u \in V.$$

(3) *If there exists  $0 < \delta < 1$ , such that  $|\mathcal{N}(u, P^{\mathcal{B}}u)| \leq \delta \mathcal{B}(u, P^{\mathcal{B}}u)$ ,  $\forall u \in V$ , then*

$$(u, P^{\mathcal{B}}u)_{\mathcal{A}} \geq c(1 - \delta)C_0^{-2}(u, u)_{\mathcal{A}}, \quad \forall u \in V,$$

where  $c > 0$  is independent of  $C_0$  and  $\delta$ .

It is interesting to note that under similar assumptions for the bilinear forms and the skewsymmetric part, an abstract theory for multiplicative Schwarz methods can also be established, see [10].

### 3 Methods for parabolic convection-diffusion problems

In this section, we apply the abstract additive Schwarz method to parabolic convection-diffusion problems. Three algorithms of this type are discussed, namely, an additive Schwarz method ( $AS^2M$ ) for problems in  $\mathbb{R}^2$  and  $\mathbb{R}^3$ , a modified  $AS^2M$  for problems in  $\mathbb{R}^2$  and an iterative substructuring method ( $IS^2M$ ) for problems in  $\mathbb{R}^2$ . The convergence rates are given in Theorem 3 - 5.

#### 3.1 An additive Schwarz method for problems in $\mathbb{R}^2$ and $\mathbb{R}^3$

The additive Schwarz method was introduced by Dryja and Widlund [13] for elliptic problems. In this section, we generalize the method to parabolic convection-diffusion problems.

We first introduce our basic decomposition of  $\Omega$  and the corresponding projections. We first extend each subregion  $\Omega_i$ , introduced in section 1.2, to obtain  $\Omega'_i$ , such that  $\Omega_i \subset \Omega'_i$  and such that there exists a constant  $\alpha > 0$  such that

$$\text{distance}(\partial\Omega'_i, \partial\Omega_i \cap \Omega) \geq \alpha H, \quad \forall i.$$

We suppose that  $\partial\Omega'_i$  does not cut through any  $h$ -level elements. We make the same constructions for the subregions that meet the boundary except that we cut off the parts that are outside  $\Omega$ . To simplify the notations, we also denote  $\Omega'_0 = \Omega$ .

It is easy to see that the finite element space  $V^h$  can be decomposed into the sum of the coarse mesh function space  $V^H$  and a number of spaces which are supported

only in subregions  $\Omega'_i$ , i.e.  $V^h = V_0^h + V_1^h + \dots + V_N^h$ , where  $V_0^h = V^H$  and  $V_i^h = V^h \cap H_0^1(\Omega'_i)$ .

Let  $P_i^{D_\tau}$  be the projection from  $V^h$  to  $V_i^h$  with respect to the bilinear form  $D_\tau(\cdot, \cdot)$ , and  $P^{D_\tau} = P_0^{D_\tau} + \dots + P_N^{D_\tau}$ . We obtain the derived equation with respect to the bilinear form  $D_\tau(\cdot, \cdot)$  and the decomposition  $\{V_i^h\}$ ,

$$P^{D_\tau} u^h = g'_{n-1}. \quad (8)$$

Here  $g'_{n-1}$  can be computed without *a priori* knowledge of the solution  $u^h$  as described in the abstract theory. We will prove that this derived system is uniformly well conditioned under certain conditions, i.e., that the condition number of  $P^{D_\tau}$  does not change if (1) we refine the fine mesh size  $h$  to increase the accuracy; (2) we refine the coarse mesh size  $H$ , increasing the number of subproblems for parallel computing purpose; (3) we increase the time step  $\tau$ .

**Theorem 3** (1) *There exists a constant  $C_p$ , such that*

$$\|P^{D_\tau} u^h\|_{A_\tau} \leq C_p \|u^h\|_{A_\tau}, \quad \forall u^h \in V^h.$$

(2) *There exists a constant  $c > 0$ , such that*

$$\|P^{D_\tau} u^h\|_{A_\tau} \geq c \|u^h\|_{A_\tau}, \quad \forall u^h \in V^h.$$

(3) *If  $c_{H,\tau} = \max\{H, H^\gamma \sqrt{H^2/\tau + 1}\}$  is small enough, i.e.  $0 < c_{H,\tau} \leq \tilde{c}_0$ , then there exists a constant  $c_p(\tilde{c}_0) > 0$ , such that*

$$(u^h, P^{D_\tau} u^h)_{A_\tau} \geq c_p(\tilde{c}_0) (u^h, u^h)_{A_\tau}, \quad \forall u^h \in V^h.$$

**Remarks:**

(a) The problem is selfadjoint if the first order terms of  $L$  vanish and then it can be seen easily that the operator  $P^{D_\tau}$  is symmetric with respect to  $A_\tau$ -norm. The standard conjugate gradient method in  $A_\tau$ -norm can therefore be used. Theorem 3 shows that the derived linear system is optimal for the conjugate gradient method in the sense that the rate of convergence does not depend on the mesh parameters  $H$  and  $h$ , nor on the time step size  $\tau$ .

(b) Since, in general, (8) is a nonsymmetric but positive definite system, we can use the GMRES [16, 23] method to solve it. We use the  $(\cdot, \cdot)_{A_\tau}$  inner product.

(c) In general,  $\tilde{c}_0$  depends on the coefficients of the first order terms in  $L$ , the ellipticity constant of  $D_\tau$ , the bounds on  $D_\tau$  and also on the geometry of the domain  $\Omega$ , which is reflected in  $\gamma$ . We do not have an explicit relation between  $\tilde{c}_0$  and the skewsymmetric part of  $L$ . From the proof of the lower bound, we know that as the skewsymmetric coefficients in  $L$  increase,  $\tilde{c}_0$  decreases.

In order to prove the theorem, we need only to show that all the assumptions of Theorem 2 hold and, in addition, that all the constants that appear in that abstract result are independent of the mesh parameters  $H$ ,  $h$  and  $\tau$ . We begin by establishing some lemmas, which contain most of the basic results.

The following lemma is well known.



**Lemma 3** *There exist two constants  $c > 0$  and  $C$ , which depend only on the shape regularity of the finite element subdivision of  $\Omega$ , such that*

$$ch^d \sum_{x_i \in \Lambda^h} (u^h(x_i))^2 \leq \|u^h\|_{L^2(\Omega)}^2 \leq Ch^d \sum_{x_i \in \Lambda^h} (u^h(x_i))^2, \forall u^h \in V^h.$$

*The statement is also true if we replace  $V^h$  by  $V^H$  and  $h$  by  $H$ .*

**Lemma 4** *For all  $u^h \in V^h$ , there exist  $u_i^h \in V_i^h$ ,  $i = 0, 1, \dots, N$ , such that*

$$u^h = \sum_{i=0}^N u_i^h$$

*and there exists a constant  $C$ , which is independent of  $h$  and  $H$ , such that*

$$\sum_{i=0}^N \|u_i^h\|_{A_\tau(\Omega'_i)}^2 \leq C \|u^h\|_{A_\tau}^2.$$

**Proof:** A proof can be found in [14] for the  $H_0^1$  norm. Thus, our only task is to show that the same estimate holds for the  $L^2$  norm. The construction of  $u_i^h$ ,  $i = 0, 1, \dots, N$  is the same as that of [14]. Let  $Q_H$  be the  $L^2$  projection into the coarse mesh space defined as

$$(Q_H u^h, v^H) = (u^h, v^H), \quad \forall v^H \in V^H.$$

We take  $u_0^h = Q_H u^h$ ,  $w^h = u^h - Q_H u^h$ , and then set  $u_i^h = I_h(\theta_i(w^h))$ , where  $\{\theta_i\}$  is a partition of unity and  $\theta_i$  belongs to  $C_0^\infty(\Omega'_i)$ . It can be arranged so that  $\nabla \theta_i$  is bounded by  $const/H$ .  $I_h$  is a interpolation operator which uses the function values at the  $h$ -level nodes only. Because  $Q_H u^h$  is the  $L^2$  projection, we have

$$\|Q_H u^h\|_{L^2(\Omega)}^2 \leq \|u^h\|_{L^2(\Omega)}^2. \quad (9)$$

Next, we prove the  $L^2(\Omega)$  boundness for the other  $u_i^h$ . By using Lemma 3 for each  $\Omega'_i$ ,

$$\sum_{i=1}^N \|I^h(\theta_i w^h)\|_{L^2(\Omega'_i)}^2 \leq Ch^2 \sum_{i=1}^N \left( \sum_{x_j \in \Lambda^h \cap \Omega'_i} ((\theta_i w^h)(x_j))^2 \right).$$

Since  $|\theta_i| \leq 1$ ,  $i = 1, \dots, N$ , and  $\{\Omega'_i\}$  is a finite covering of  $\Omega$ , the right hand side can be bounded by  $h^2 \sum_{x_i \in \Lambda^h} (w^h(x_i))^2$ . Using Lemma 3 again, we obtain

$$\sum_{i=1}^N \|u_i^h\|_{L^2(\Omega'_i)}^2 \leq C \|w^h\|_{L^2(\Omega)}^2 \leq C \|u^h\|_{L^2(\Omega)}^2. \quad (10)$$

We conclude the proof by combining the estimates (9) and (10) and the estimate in the  $H_0^1$  norm.  $\square$

**Lemma 5** For any  $w \in H_0^1(\Omega) \cap H^{1+\gamma}(\Omega)$ , there exists a  $w^H \in V^H$ , such that

$$\|w - w^H\|_{\tau(\Omega)} \leq CH^\gamma \sqrt{H^2 + \tau} |w|_{H^{1+\gamma}(\Omega)}.$$

**Proof:** For any given  $w$ , let  $w^H$  be the solution of

$$a(w^H, v) = a(w, v), \quad \forall v \in V^H,$$

where  $a(u, v) = \int_\Omega \nabla u \nabla v d\Omega$ . By classical finite element theory, we know that

$$\|w - w^H\|_{H_0^1(\Omega)} \leq CH^\gamma |w|_{H^{1+\gamma}(\Omega)},$$

$$\|w - w^H\|_{L^2(\Omega)} \leq CH^{1+\gamma} |w|_{H^{1+\gamma}(\Omega)}.$$

Hence, we have  $\|w - w^H\|_{\tau(\Omega)} \leq CH^\gamma \sqrt{H^2 + \tau} |w|_{H^{1+\gamma}(\Omega)}$ .  $\square$

**Lemma 6** There exists a constant  $C$ , which is independent of  $H$  and  $\tau$ , such that

$$\|u^h - P_0^{D\tau} u^h\|_{L^2(\Omega)} \leq CH^\gamma (\sqrt{H^2 + \tau}/\tau) \|u^h\|_{\tau(\Omega)}, \quad \forall u^h \in V^h.$$

**Proof:** We first establish a bound for the  $\tau$ -norm. It is easy to verify that

$$D_\tau(u^h - P_0^{D\tau} u^h, u^h - P_0^{D\tau} u^h) = D_\tau(u^h - P_0^{D\tau} u^h, u^h).$$

Using Lemma 1, we get the estimate

$$\|u^h - P_0^{D\tau} u^h\|_{\tau(\Omega)} \leq C \|u^h\|_{\tau(\Omega)}.$$

The  $H_0^1(\Omega)$  (or  $\tau$ -norm) estimate and Nitsche's trick give the  $L^2(\Omega)$  estimate. We have

$$\|u^h - P_0^{D\tau} u^h\|_{L^2(\Omega)} = \sup_{\|v\|_{L^2(\Omega)} \neq 0} \frac{(u^h - P_0^{D\tau} u^h, v)}{\|v\|_{L^2(\Omega)}}.$$

For any fixed  $v \in L^2(\Omega)$ , we form an auxiliary problem: Find  $w \in H^{1+\gamma}(\Omega) \cap H_0^1(\Omega)$ , such that

$$D_\tau(\phi, w) = (\phi, v), \quad \forall \phi \in H_0^1(\Omega).$$

By Lemma 2, we have  $|w|_{H^{1+\gamma}(\Omega)} \leq C/\tau \|v\|_{L^2(\Omega)}$ . Let  $\phi = u^h - P_0^{D\tau} u^h \in H_0^1(\Omega)$ . Then

$$(u^h - P_0^{D\tau} u^h, v) = D_\tau(u^h - P_0^{D\tau} u^h, w).$$

Let  $w^H$  be the  $V^H$  approximation of  $w$  obtained in Lemma 5. Since  $w^H \in V^H$ , we have

$$\begin{aligned} |D_\tau(u^h - P_0^{D\tau} u^h, w)| &= |D_\tau(u^h - P_0^{D\tau} u^h, w - w^H)| \\ &\leq C \|u^h - P_0^{D\tau} u^h\|_{\tau(\Omega)} \|w - w^H\|_{\tau(\Omega)} \\ &\leq CH^\gamma \sqrt{H^2 + \tau} \|u^h\|_{\tau(\Omega)} |w|_{H^{1+\gamma}(\Omega)}. \end{aligned}$$

Combining the above results, we obtain

$$\|u^h - P_0^{D\tau} u^h\|_{L^2(\Omega)} \leq CH^\gamma (\sqrt{H^2 + \tau}/\tau) \|u^h\|_{\tau(\Omega)}. \square$$

In the next lemma, we estimate the contribution from the skewsymmetric part  $N_\tau(\cdot, \cdot)$ . We show that it is a lower order term compared with the symmetric part, and that it can be controlled if the coarse mesh size is fine enough.

**Lemma 7** *If  $\max\{H, H^\gamma \sqrt{H^2/\tau + 1}\}$  is small enough, there exists a constant  $0 < \delta < 1$ , such that*

$$|N_\tau(u^h, P^{D_\tau} u^h)| \leq \delta D_\tau(u^h, P^{D_\tau} u^h), \quad \forall u^h \in V^h.$$

**Proof:** Since  $N_\tau(P_i^{D_\tau} u^h, P_i^{D_\tau} u^h) = 0$ , it is easy to verify that

$$|N_\tau(u^h, P^{D_\tau} u^h)| \leq \sum_{i=0}^N |N_\tau(P_i^{D_\tau} u^h, u^h - P_i^{D_\tau} u^h)|.$$

Therefore, we need only to estimate the right hand side of the above inequality. Since the coarse mesh projection  $P_0^{D_\tau}$  is special, we consider it separately using Lemma 6.

(1)  $i = 0$ ; By Lemma 6

$$\begin{aligned} |N_\tau(P_0^{D_\tau} u^h, u^h - P_0^{D_\tau} u^h)| &\leq C\tau \|P_0^{D_\tau} u^h\|_{H_0^1(\Omega)} \|u^h - P_0^{D_\tau} u^h\|_{L^2(\Omega)} \\ &\leq CH^\gamma \sqrt{H^2 + \tau} \|P_0^{D_\tau} u^h\|_{H_0^1(\Omega)} \cdot \|u^h\|_{\tau(\Omega)}. \end{aligned}$$

It is easy to see that  $\sqrt{\tau} \|v\|_{H_0^1(\Omega)} \leq \|v\|_{\tau(\Omega)}$ . Hence,

$$|N_\tau(P_0^{D_\tau} u^h, u^h - P_0^{D_\tau} u^h)| \leq CH^\gamma \sqrt{H^2/\tau + 1} \|P_0^{D_\tau} u^h\|_{\tau(\Omega)} \|u^h\|_{\tau(\Omega)}.$$

From the definition of  $P_0^{D_\tau}$ , we have

$$D_\tau(P_0^{D_\tau} u^h, P_0^{D_\tau} u^h) = D_\tau(u^h, P_0^{D_\tau} u^h).$$

We thus have, by Lemma 1,  $\|P_0^{D_\tau} u^h\|_{\tau(\Omega)} \leq C \|u^h\|_{\tau(\Omega)}$ .

Therefore, the first term can be bounded as follows

$$|N_\tau(P_0^{D_\tau} u^h, u^h - P_0^{D_\tau} u^h)| \leq CH^\gamma \sqrt{H^2/\tau + 1} \|u^h\|_{\tau(\Omega)}^2.$$

(2)  $i \neq 0$ :

$$\begin{aligned} |N_\tau(P_i^{D_\tau} u^h, u^h - P_i^{D_\tau} u^h)| &\leq C\tau \|P_i^{D_\tau} u^h\|_{L^2(\Omega'_i)} (\|u^h\|_{H_0^1(\Omega'_i)} + \|P_i^{D_\tau} u^h\|_{H_0^1(\Omega'_i)}) \\ &\leq C\sqrt{\tau} \|P_i^{D_\tau} u^h\|_{L^2(\Omega'_i)} (\|u^h\|_{\tau(\Omega'_i)} + \|P_i^{D_\tau} u^h\|_{\tau(\Omega'_i)}). \end{aligned}$$

The first factor can be estimated by using Friedrich's inequality. Since  $P_i^{D_\tau} u^h \in H_0^1(\Omega'_i)$  and the diameter of  $\Omega'_i$  is of order  $H$ , we have

$$\|P_i^{D_\tau} u^h\|_{L^2(\Omega'_i)} \leq CH \|P_i^{D_\tau} u^h\|_{H_0^1(\Omega'_i)}.$$

By the definition of  $P_i^{D_\tau}$ ,

$$D_\tau(P_i^{D_\tau} u^h, P_i^{D_\tau} u^h) = D_\tau(u^h, P_i^{D_\tau} u^h).$$

We obtain, by Lemma 1,  $\|P_i^{D_\tau} u^h\|_{\tau(\Omega'_i)} \leq C \|u^h\|_{\tau(\Omega'_i)}$ .

Combining these inequalities, we obtain

$$|N_\tau(P_i^{D_\tau} u^h, u^h - P_i^{D_\tau} u^h)| \leq CH \|u^h\|_{\tau(\Omega'_i)}^2.$$

Putting the results in (1) and (2) together, we have

$$|N_\tau(P^{D_\tau} u^h, u^h)| \leq C \max\{H, H^\gamma \sqrt{H^2/\tau + 1}\} \sum_{i=0}^N \|u^h\|_{\tau(\Omega'_i)}^2, \quad (11)$$

which also holds in the  $A_\tau$ -norm. It is easy to see that

$$\sum_{i=0}^N \|P_i^{D_\tau} u^h\|_{A_\tau(\Omega'_i)}^2 = D_\tau(u^h, P^{D_\tau} u^h).$$

By using the decomposition lemma, we have

$$\begin{aligned} \|u^h\|_{A_\tau}^2 &= \sum_{i=0}^N D_\tau(P_i^{D_\tau} u^h, u_i^h) \\ &\leq C \sum_{i=0}^N \|P_i^{D_\tau} u^h\|_{A_\tau(\Omega'_i)} \|u_i^h\|_{A_\tau(\Omega'_i)} \\ &\leq C \sqrt{\sum_{i=0}^N \|u_i^h\|_{A_\tau(\Omega'_i)}^2} \sqrt{\sum_{i=0}^N \|P_i^{D_\tau} u^h\|_{A_\tau(\Omega'_i)}^2}. \end{aligned}$$

Therefore, we obtain

$$\|u^h\|_{A_\tau}^2 \leq C \sum_{i=0}^N \|P_i^{D_\tau} u^h\|_{A_\tau(\Omega'_i)}^2.$$

Hence

$$\|u^h\|_{A_\tau}^2 \leq CD_\tau(u^h, P^{D_\tau} u^h). \quad (12)$$

We complete the proof by combining estimates (11) and (12).  $\square$

The proof of Theorem 3 follows from Theorem 2 and the lemmas in this section.

### 3.2 The modified additive Schwarz method for problems in $\mathbb{R}^2$

In this subsection, we propose a modified version of  $AS_M$  obtained by dropping the coarse mesh space  $V_0^h$ , which provided the global transportation of information. We show that in some situations, with precise conditions given in the following theorem, the global space is not necessary. This is in contrast to the elliptic case, for fast convergence, this alternative algorithm is more suitable for parallel computers.

Let us define  $\tilde{P}^{D_\tau} = P_1^{D_\tau} + \cdots + P_N^{D_\tau}$ , where the  $P_i^{D_\tau}$  are the same as in the previous subsection. We have the following theorem.

**Theorem 4** (1) *There exists a constant  $C_{\tilde{p}}$ , such that*

$$\|\tilde{P}^{D_\tau} u^h\|_{A_\tau} \leq C_{\tilde{p}} \|u^h\|_{A_\tau}, \quad \forall u^h \in V^h.$$

(2) *There exists a constant  $c > 0$ , such that*

$$\|\tilde{P}^{D_\tau} u^h\|_{A_\tau} \geq c(1 + \tau/H^2)^{-1} \|u^h\|_{A_\tau}, \quad \forall u^h \in V^h.$$

(3) If  $c_{H,\tau} = H(1 + \tau/H^2)$  is small enough, i.e.  $0 < c_{H,\tau} \leq \tilde{c}_0$ , then there exists a constant  $c_{\tilde{p}}(\tilde{c}_0) > 0$ , such that

$$(u^h, \tilde{P}^{D\tau} u^h)_{A_\tau} \geq c_{\tilde{p}}(\tilde{c}_0)(u^h, u^h)_{A_\tau}, \quad \forall u^h \in V^h.$$

**Remarks:**

(a) For symmetric problems, parts (1) and (2) of this theorem show that if the factor  $\tau/H^2$  is small, the elimination of the coarse mesh space does not lead to slow convergence. This suggests that we can use the modified  $\mathcal{A}_\tau M$  when we have a relatively small time step or large substructures.

For nonsymmetric problems, we need that both  $H$  and  $\tau/H^2$  to be small in order to obtain the fast convergence. We cannot choose  $\tau$  and  $H$  independently.

(b) This theorem has been established only in  $\mathbb{R}^2$ . In higher dimension, the fine mesh size  $h$  enters our bounds.

Theorem 4 can be proved by using the results of Theorem 2, a new decomposition lemma given below and a new estimate of the skewsymmetric part. We begin with the decomposition lemma.

**Lemma 8** For all  $u^h \in V^h$ , there exist  $u_i^h \in V_i^h$ ,  $i = 1, \dots, N$ , such that

$$u^h = \sum_{i=1}^N u_i^h$$

and, there exists a constant  $C$ , which is independent of  $h$ ,  $H$  and  $\tau$ , such that

$$\sum_{i=1}^N \|u_i^h\|_{A_\tau(\Omega'_i)}^2 \leq C(1 + \tau/H^2) \|u^h\|_{A_\tau}^2.$$

**Proof:** We first construct the decomposition and then we do the estimates in both the  $H_0^1(\Omega)$  and the  $L^2(\Omega)$  norm. Let  $\{\theta_i, i = 1, \dots, N\}$  be the partition of unity, defined in the proof of Lemma 4, of  $\Omega$  and denote  $u_i^h = I_h(\theta_i u^h)$ . For each subregion  $\Omega'_i$ , we have, in the sense of equivalent norms, that

$$|u_i^h|_{H_0^1(\Omega'_i)}^2 = \sum ((\theta_i u^h)(x_l) - (\theta_i u^h)(x_m))^2,$$

where the sum is taken over all adjacent pairs of nodal points  $x_l$  and  $x_m$  in  $\Omega'_i$ . Let  $K \subset \Omega'_i$  be a single element and  $x_l, x_m \in \bar{K}$ . Let

$$\bar{\theta}_{ilm} = 1/2(\theta_i(x_l) + \theta_i(x_m)).$$

We then have

$$\begin{aligned} & (\theta_i u^h)(x_l) - (\theta_i u^h)(x_m) = \\ & (\theta_i(x_l) - \bar{\theta}_{ilm})u^h(x_l) - (\theta_i(x_m) - \bar{\theta}_{ilm})u^h(x_m) + \bar{\theta}_{ilm}(u^h(x_l) - u^h(x_m)), \end{aligned}$$

which can be bounded from above by

$$C(h/H \max_{\bar{K}} \{ |u^h(x)| \} + |u^h(x_l) - u^h(x_m)|).$$

By squaring this estimate, using the triangle inequality and summing over all  $K \subset \Omega'_i$ , we obtain

$$\begin{aligned} |u_i^h|_{H_0^1(\Omega'_i)}^2 &\leq C(H^{-2} \sum_{K \subset \Omega'_i} \max_{\bar{K}} \{|u^h(x)|\}^2 \cdot h^2 + |u^h|_{H_0^1(\Omega'_i)}^2) \\ &\leq C(H^{-2} \|u^h\|_{L^2(\Omega'_i)}^2 + |u^h|_{H_0^1(\Omega'_i)}^2). \end{aligned}$$

Using Friedrich's inequality, it is easy to establish that

$$\sum_{i=1}^N \|u_i^h\|_{H_0^1(\Omega'_i)}^2 \leq (1 + CH^2) \sum_{i=1}^N |u_i^h|_{H_0^1(\Omega'_i)}^2.$$

Therefore,

$$\sum_{i=1}^N \|u_i^h\|_{H_0^1(\Omega'_i)}^2 \leq C(H^{-2} \|u^h\|_{L^2(\Omega)}^2 + |u^h|_{H_0^1(\Omega)}^2) \leq CH^{-2} \|u^h\|_{H_0^1(\Omega)}^2.$$

By combining the estimate given for the  $L^2$  norm in the proof of Lemma 4 and the definition of the  $\tau$ -norm, we complete the proof.  $\square$

In the next lemma, we estimate the contribution from the skewsymmetric part  $N_\tau(\cdot, \cdot)$ .

**Lemma 9** *There exists a constant  $C$ , which is independent of  $h, H$  and  $\tau$ , such that*

$$|N_\tau(u^h, \tilde{P}^{D_\tau} u^h)| \leq CH(1 + \tau/H^2) D_\tau(u^h, \tilde{P}^{D_\tau} u^h), \quad \forall u^h \in V^h.$$

The proof is very similar, replacing the decomposition Lemma 4 by Lemma 8, to that of Lemma 7 part (2).

### 3.3 An iterative substructuring method for problems in $\mathbb{R}^2$

We present a domain decomposition method which appears to be of non-overlapping type, but nevertheless fits into the additive Schwarz framework. In the first two algorithms, we extended each  $H$ -level substructure  $\Omega_i$  to a larger region  $\Omega'_i$ , which defined our subspace and the corresponding projection.  $\{\Omega'_i\}$  is thus the basic decomposition of the domain. Now, instead of extending each region, we combine each pair of adjacent substructures. Denote by  $\Gamma_{ij}$  the common edge of two adjacent substructures  $\Omega_i$  and  $\Omega_j$ , and let  $\Omega_{ij} = \Omega_i \cup \Gamma_{ij} \cup \Omega_j$ . We are now going to use  $\{\Omega_{ij}\}$  as the basic decomposition of  $\Omega$ , on which we will define our projections. This will eventually lead to an iterative substructuring algorithm. Let  $\Omega_{00} = \Omega$  and denote

$$\Lambda^E = \{(i, j) \mid x_i, x_j \in \Lambda^H, x_i, x_j \text{ adjacent, or } i = j = 0\}.$$

For each  $(0, 0) \neq (i, j) \in \Lambda^E$ , we define  $V_{ij}^h \equiv H_0^1(\Omega_{ij}) \cap V^h$  and also  $V_{00}^h = V^H$ . Denote by  $P_{ij}^{D_\tau}$  the projection associated with  $V_{ij}^h$  using the bilinear form  $D_\tau(\cdot, \cdot)$ . It is easy to verify that

$$V^h = \sum_{(i,j) \in \Lambda^E} V_{ij}^h.$$

Let us also denote

$$P^{D\tau} = \sum_{(i,j) \in \Lambda^E} P_{ij}^{D\tau} : V^h \longrightarrow V^h$$

and the derived equation can thus be defined as

$$P^{D\tau} \mathbf{u} = \mathbf{b}^{D\tau}.$$

In the next theorem, we present the bounds for the operator  $P^{D\tau}$ . We use the notation

$$\nu(h, H, \tau) = (1 + H^2/\tau + \log H/h)(1 + \log H/h).$$

**Theorem 5** (1) *There exists a constant  $C_p$ , such that*

$$\|P^{D\tau} u^h\|_{A_\tau} \leq C_p \|u^h\|_{A_\tau}, \quad \forall u^h \in V^h.$$

(2) *There exists a constant  $c > 0$ , such that*

$$\|P^{D\tau} u^h\|_{A_\tau} \geq c/\nu(h, H, \tau) \|u^h\|_{A_\tau}, \quad \forall u^h \in V^h.$$

(3) *There exists a constant  $c_p > 0$ , such that for  $H\sqrt{\nu}$  small enough,*

$$(u^h, P^{D\tau} u^h)_{A_\tau} \geq c_p/\nu(h, H, \tau) (u^h, u^h)_{A_\tau}, \quad \forall u^h \in V^h.$$

As before, we need only to prove that the assumptions of Theorem 2 hold and study how the constants depend on the mesh parameters  $h, H$  and  $\tau$ . The proofs look more complicated because of the lack of overlap between the substructures, but the idea behind the proof is similar. We need to establish a bounded decomposition lemma for  $V^h$  functions. The bounds depend mildly on the parameters  $h$  and  $H$ .

We first give a lemma, which plays an important role in the traditional theory for iterative substructuring algorithms. Variations of this result, which dates back at least to 1966, are given in Bramble [3], Bramble, Pasciak and Schatz [4] and Yserentant [25].

**Lemma 10** *Let  $\Omega_i$  be a substructure,  $u^h \in V_i^h$  and let  $\overline{u_i^h} = \frac{1}{\text{area}(\Omega_i)} \int_{\Omega_i} u^h d\Omega$ . Then*

$$\|u^h - \overline{u_i^h}\|_{L^\infty(\Omega_i)}^2 \leq C(1 + \log H/h) |u^h|_{H^1(\Omega_i)}^2.$$

**Lemma 11** *For any  $u^h \in V^h$ , there exist  $u_{ij}^h \in V_{ij}^h$ ,  $(i, j) \in \Lambda^E$ , such that  $u^h = \sum_{(i,j) \in \Lambda^E} u_{ij}^h$  and the decomposition is bounded as follows*

$$\begin{aligned} (a) \quad & \sum_{(i,j) \in \Lambda^E} \|u_{ij}^h\|_{H_0^1(\Omega_{ij})}^2 \leq C(1 + \log H/h)^2 \|u^h\|_{H_0^1(\Omega)}^2 \\ (b) \quad & \sum_{(i,j) \in \Lambda^E} \|u_{ij}^h\|_{L^2(\Omega_{ij})}^2 \leq C(\|u^h\|_{L^2(\Omega)}^2 + H^2(1 + \log H/h) \|u^h\|_{H_0^1(\Omega)}^2). \end{aligned}$$

**Proof:** The construction of the decomposition and the proof of (a) were given by Dryja and Widlund in [14]. We shall focus on (b) only. Let  $I_H$  be the interpolation operator using the function values at the  $H$ -level nodes only. For a given  $u^h \in V^h$ , we estimate the  $L^2(\Omega)$  norm of  $u_{00}^h = I_H u^h$ . Let us consider one substructure at a time.

Assume that  $\Omega_i$  has vertices  $T_1, T_2, T_3$  and denote  $\alpha_i = \frac{1}{\text{area}(\Omega_i)} \int_{\Omega_i} u^h d\Omega$ . Then, we have

$$\|I_H u^h\|_{L^2(\Omega_i)}^2 \leq C(\|I_H u^h - \alpha_i\|_{L^2(\Omega_i)}^2 + \|\alpha_i\|_{L^2(\Omega_i)}^2).$$

Since the function  $I_H u^h - \alpha_i$  is linear in the region  $\Omega_i$ , a straightforward calculation shows that the  $L^2$  norm in  $\Omega_i$  can be bounded by

$$CH^2((I_H u^h - \alpha_i)^2|_{T_1} + (I_H u^h - \alpha_i)^2|_{T_2} + (I_H u^h - \alpha_i)^2|_{T_3}).$$

This expression is, using Lemma 10, bounded by  $CH^2(1 + \log H/h) |u^h|_{H^1(\Omega_i)}^2$ . We bound the  $\|\alpha_i\|_{L^2(\Omega_i)}$  term, by using Schwarz's inequality

$$\|\alpha_i\|_{L^2(\Omega_i)}^2 = \int_{\Omega_i} \left(\frac{1}{\text{area}(\Omega_i)} \int_{\Omega_i} u^h d\Omega\right)^2 d\Omega \leq \|u^h\|_{L^2(\Omega_i)}^2.$$

Therefore, we obtain

$$\|I_H u^h\|_{L^2(\Omega_i)}^2 \leq C(\|u^h\|_{L^2(\Omega_i)}^2 + H^2(1 + \log H/h) |u^h|_{H^1(\Omega_i)}^2).$$

By summing these inequalities, using Friedrich's inequality and replacing the  $H^1$  semi-norm by the  $H^1$  norm, we obtain

$$\|I_H u^h\|_{L^2(\Omega)}^2 \leq C(\|u^h\|_{L^2(\Omega)}^2 + H^2(1 + \log H/h) \|u^h\|_{H_0^1(\Omega)}^2).$$

Now consider the case  $(i, j) \neq (0, 0)$ . Recalling the construction of [14],  $u_{ij}^h = I_h(\theta_{ij}(u^h - I_H u^h))$ , where  $\{\theta_{ij}\}$  is a partition of unity of  $\Omega \setminus \{H\text{-level nodes}\}$ . More precise, as described in [14],  $\theta_{ij}$  must equal to zero at all  $H$ -level nodes and equal to one at all nodal points in the interior of  $\Gamma_{ij}$ . Let us denote  $\Lambda_{ij}^h = \Omega_{ij} \cap \Lambda^h$ . By Lemma 3, we have

$$\|u_{ij}^h\|_{L^2(\Omega_{ij})}^2 \leq Ch^2 \sum_{x_k \in \Lambda_{ij}^h} ((I_h(\theta_{ij}(u^h - I_H u^h)))(x_k))^2.$$

Since  $|\theta_{ij}| \leq 1$ , and  $x_k$  is a nodal point, the interpolation operator  $I_h$  can be removed. Therefore the right hand side can be bounded by

$$Ch^2 \sum_{x_k \in \Lambda_{ij}^h} ((u^h - I_H u^h)(x_k))^2.$$

Taking into account that the value at a nodal point contributes at most three times, we obtain

$$\sum_{(i,j) \in \Lambda^E \setminus (0,0)} \|u_{ij}^h\|_{L^2(\Omega_{ij})}^2 \leq Ch^2 \sum_{x_k \in \Lambda^h} ((u^h - I_H u^h)(x_k))^2.$$

We can bound the right hand side, using Lemma 3, by  $C\|u^h - I_H u^h\|_{L^2(\Omega)}^2$ , which can then be bounded by  $C(\|u^h\|_{L^2(\Omega)}^2 + H^2(1 + \log H/h) \|u^h\|_{H_0^1(\Omega)}^2)$ .  $\square$

An immediate consequence of this lemma is the following:



**Lemma 12** *The decomposition in Lemma 11 is bounded in the  $A_\tau$ -norm, i.e. there exists a constant  $C$ , independent of  $H, h$  and  $\tau$ , such that*

$$\sum_{(i,j) \in \Lambda^E} \|u_{ij}^h\|_{\tau(\Omega_{ij})}^2 \leq C\nu(h, H, \tau) \|u^h\|_{\tau(\Omega)}^2.$$

**Lemma 13** *There exists a constant  $C$ , independent of  $h, H$  and  $\tau$ , such that*

$$|N_\tau(u^h, P^{D_\tau} u^h)| \leq C \max\{H, H^\gamma \sqrt{H^2/\tau + 1}\} \nu(h, H, \tau) D_\tau(u^h, P^{D_\tau} u^h), \quad \forall u^h \in V^h.$$

**Proof:** It is easy to see that the coarse mesh projection  $P_{00}^{D_\tau}$  is identical to the coarse mesh projection  $P_0^{D_\tau}$  defined previously. It follows from the proof of Lemma 7 part (1) and (2), that we have

$$|N_\tau(P^{D_\tau} u^h, u^h)| \leq C \max\{H, H^\gamma \sqrt{H^2/\tau + 1}\} \sum_{(i,j) \in \Lambda^E} \|u^h\|_{\tau(\Omega_{ij})}^2. \quad (13)$$

By using the same argument as when equation (12) was proved and the decomposition lemma 12, we have

$$\|u^h\|_{A_\tau}^2 \leq C\nu(h, H, \tau) D_\tau(u^h, P^{D_\tau} u^h). \square$$

The proof of Theorem 5 is completed by using Theorem 2, the lemmas in this section and the assumption that  $\max\{H, H^\gamma \sqrt{H^2/\tau + 1}\} \nu(h, H, \tau)$  is small enough.

## 4 Numerical results

In this section, we present some numerical results with model problems. To specify our model problems, we need only to give the elliptic parts of the parabolic operators. We consider the following linear second order elliptic operator defined on  $\Omega = [0, 1] \times [0, 1] \subset \mathbb{R}^2$ ,

$$Lu = -\frac{\partial}{\partial x} \left( \xi \frac{\partial u}{\partial x} \right) - \frac{\partial}{\partial y} \left( \eta \frac{\partial u}{\partial y} \right) + \alpha \frac{\partial u}{\partial x} + \beta \frac{\partial u}{\partial y} + \gamma u = f,$$

with the homogenous Dirichlet boundary condition. The coefficients are specified as follows.

**Example 0.**  $\xi = 1$ ,  $\eta = 1$  and  $\alpha = \beta = \gamma = 0$ . This is a selfadjoint problem, used to test the iterative substructuring algorithm.  $f$  is chosen so that the solution is  $u = xe^{xy} \sin(\pi x) \sin(\pi y)$ .

**Example 1.**  $\xi = 1 + x^2 + y^2$ ,  $\eta = e^{xy}$ ,  $\alpha = 5(x + y)$ ,  $\beta = 1/(1 + x + y)$  and  $\gamma = 0$ .  $u$  is the same as in Example 0.

**Example 2.** The coefficients are chosen as  $\xi = \sigma$ ,  $\eta = \sigma$ ,  $\alpha = 1$ ,  $\beta = 1$  and  $\gamma = 1$ .  $\sigma$  will be specified later.  $u$  is the same as in Example 0.

The  $H$ -level subdivision of  $\Omega$  is as shown in Figure 1. We further divide each subregion in a similar fashion into  $h$ -level triangles, which are not shown in the figure.

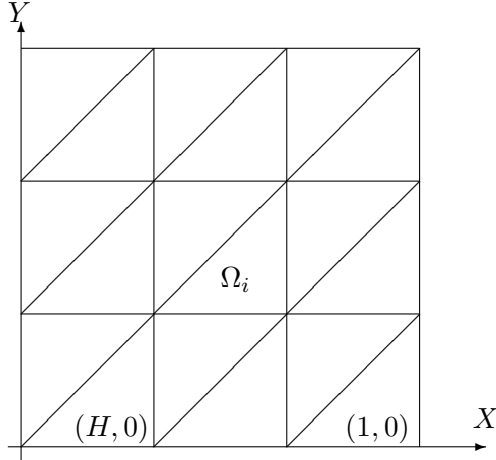


Figure 1. H-level subdivision

We use **ovlp** to denote the size of the overlap, i.e.  $\mathbf{ovlp} = \text{distance}(\partial\Omega'_i, \partial\Omega_i \cap \Omega)$ . In our Fortran program, all the subproblems are solved exactly by the band solver from LINPACK. The stopping criterion for the GMRES method is  $\|r_i\|_{A_\tau} / \|r_0\|_{A_\tau} \leq 10^{-4}$ , where  $r_i$  is the residual after the  $i^{\text{th}}$  step. The programs were run in single precision on a CONVEX C-1 computer at New York University.

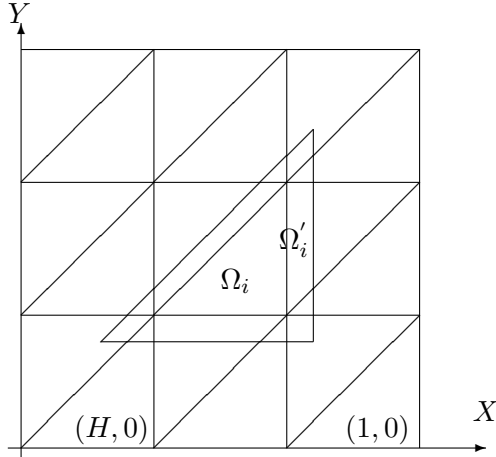


Figure 2. extended subregions

We assume that the time step has the form

$$\tau = h^\epsilon,$$

where  $\epsilon > 0$ . The main issue here is to determine how the rate of convergence ( number of iterations ) depends on  $\epsilon$ . In the case of  $\epsilon = 2.0$ , the stiffness matrix is well conditioned with a condition number independent of the mesh parameters  $h$  and  $\tau$ . In

general, the smaller  $\epsilon$  is ( the larger the time step is ) the more ill-conditioned is the problem. If  $\epsilon$  is very small, the accuracy of the solution is lost. We include some test results by using unusually large time step  $\tau$  to show the robustness of the algorithm. In our experiments, we have used  $\epsilon \in [0.25, 1.5]$ .

#### 4.1 Tests of the additive Schwarz method

	$\epsilon$	1.5	1.0	0.5	0.25
Example 0	$h=1/45, H=1/3, \text{ovlp}=5h$	10	12	12	13
	$h=1/60, H=1/3, \text{ovlp}=6h$	10	12	13	13
	$h=1/60, H=1/10, \text{ovlp}=2h$	10	11	12	12
	$h=1/45, H=1/3, \text{ovlp}=h$	15	18	18	19
Example 1	$h=1/15, H=1/3, \text{ovlp}=2h$	12	13	14	14
	$h=1/30, H=1/3, \text{ovlp}=4h$	12	14	14	14
	$h=1/45, H=1/3, \text{ovlp}=6h$	12	14	15	15
	$h=1/60, H=1/3, \text{ovlp}=8h$	11	14	15	15
	$h=1/15, H=1/5, \text{ovlp}=1h$	11	12	12	12
	$h=1/30, H=1/5, \text{ovlp}=2h$	12	14	14	14
	$h=1/45, H=1/5, \text{ovlp}=3h$	12	13	14	15
	$h=1/60, H=1/5, \text{ovlp}=4h$	11	13	15	15
Example 2 $\sigma = \sqrt{2}/30$	$h=1/45, H=1/3, \text{ovlp}=5h$	11	17	28	35
	$h=1/45, H=1/9, \text{ovlp}=2h$	9	12	14	16
	$h=1/45, H=1/15, \text{ovlp}=1h$	7	8	10	10
	$h=1/60, H=1/10, \text{ovlp}=2h$	9	11	14	15
	$h=1/60, H=1/15, \text{ovlp}=1h$	7	9	10	11
	$h=1/60, H=1/20, \text{ovlp}=1h$	6	8	10	10
Example 2 $\sigma = \sqrt{2}/100$	$h=1/60, H=1/15, \text{ovlp}=1h$	8	11	18	20
	$h=1/60, H=1/20, \text{ovlp}=1h$	7	9	13	16

This table shows that the number of iterations is almost independent on the mesh parameters as well as the time step size, especially for symmetric problems. For non-symmetric, see Theorem 3. For nonsymmetric problems, a relatively small coarse mesh size is needed for reducing the total number of iterations, which confirms Theorem 3 part(3). The larger the Reynold's number, which is roughly the ratio of the first order term coefficients and the second order term coefficients, is, the smaller the coarse mesh size is needed. As in the stationary case, using less overlap can increase the total number of iterations. This can be seen from the last row of Example 0.

## 4.2 Tests of the modified additive Schwarz method

	$\epsilon$	1.5	1.0	0.5	0.25
Example 0	$h=1/45, H=1/3, \text{ovlp}=5h$	10	12	14	15
	$h=1/60, H=1/3, \text{ovlp}=6h$	9	12	15	16
	$h=1/60, H=1/10, \text{ovlp}=2h$	12	24	34	35
	$h=1/45, H=1/3, \text{ovlp}=h$	12	21	27	28
Example 1	$h=1/15, H=1/3, \text{ovlp}=2h$	13	15	15	15
	$h=1/30, H=1/3, \text{ovlp}=4h$	13	14	16	16
	$h=1/45, H=1/3, \text{ovlp}=6h$	13	15	16	17
	$h=1/60, H=1/3, \text{ovlp}=8h$	12	15	17	17
	$h=1/15, H=1/5, \text{ovlp}=1h$	14	18	19	20
	$h=1/30, H=1/5, \text{ovlp}=2h$	15	20	23	24
	$h=1/45, H=1/5, \text{ovlp}=3h$	13	19	24	25
	$h=1/60, H=1/5, \text{ovlp}=4h$	13	18	23	25

	$\epsilon$	1.5	1.0	0.75	0.5
Example 2 $\sigma = \sqrt{2}/30$	$h=1/45, H=1/3, \text{ovlp}=5h$	13	22	31	42
	$h=1/45, H=1/5, \text{ovlp}=3h$	12	18	23	32
	$h=1/45, H=1/9, \text{ovlp}=2h$	10	14	19	30
	$h=1/45, H=1/15, \text{ovlp}=1h$	7	12	22	36
	$h=1/60, H=1/10, \text{ovlp}=2h$	10	13	20	33
	$h=1/60, H=1/15, \text{ovlp}=1h$	8	13	24	41
Example 2 $\sigma = \sqrt{2}/100$	$h=1/45, H=1/9, \text{ovlp}=2h$	10	19	31	+50
	$h=1/45, H=1/15, \text{ovlp}=1h$	7	13	25	48
	$h=1/60, H=1/15, \text{ovlp}=1h$	7	13	25	+50
	$h=1/60, H=1/20, \text{ovlp}=1h$	6	12	24	+50

This table shows that if the factor  $\tau/H^2$  is small, the results are quite satisfactory compared with the  $AS_M$  which also uses a coarse mesh. This is true especially for symmetric and mildly nonsymmetric problems; cf. Example 0 and Example 1. However, if  $\tau/H^2$  is large, i.e. the time step or the number of substructures large, the number of iterations can be large; cf. the third row of Example 0. For the problems with high Reynold's number, the algorithm becomes very sensitive to the parameter  $\epsilon$ ; cf. the last column of the second table. See also Theorem 4.

We also note that using less overlap can increase the total number of iterations. This can be seen by comparing the first and last rows in Example 0.

By comparing with the  $AS_M$  which uses coarse mesh, we see that this simpler algorithm is more sensitive to the mesh parameters and the overlapping factor.

### 4.3 Tests of the iterative substructuring method

Example 0	$\epsilon$	1.5	1.0	0.5	0.25
	$h=1/45, H=1/3$	14	14	14	14
	$h=1/60, H=1/3$	15	15	14	14
	$h=1/60, H=1/10$	11	12	12	13
Example 1	$h=1/45, H=1/5$	18	17	17	17
	$h=1/45, H=1/9$	14	14	14	14
	$h=1/60, H=1/5$	18	17	17	18
	$h=1/60, H=1/10$	13	14	14	14
Example 2 $\sigma = \sqrt{2}/30$	$h=1/45, H=1/5$	15	25	35	38
	$h=1/45, H=1/9$	13	18	20	22
	$h=1/45, H=1/15$	10	12	13	14
	$h=1/60, H=1/5$	15	22	36	41
	$h=1/60, H=1/10$	13	17	20	22
	$h=1/60, H=1/15$	11	14	15	16
Example 2 $\sigma = \sqrt{2}/100$	$h=1/45, H=1/15$	11	15	20	22
	$h=1/60, H=1/15$	12	18	24	27
	$h=1/60, H=1/20$	10	14	17	18

For the symmetric problems, we see that the number of iterations is insensitive to the fine mesh size  $h$ , but depends on the coarse mesh size  $H$ . The number of iterations is reduced when the coarse mesh size is reduced. This is not true for elliptic problems. The algorithm is not very sensitive to the time step parameter.

For nonsymmetric problems, the algorithm becomes more sensitive to the time step parameter and the coarse mesh size is more important for controlling the number of iterations.

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