A FETI-DP TYPE DOMAIN DECOMPOSITION ALGORITHM FOR THREE-DIMENSIONAL INCOMPRESSIBLE STOKES EQUATIONS

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Abstract. The FETI-DP algorithms, proposed by the authors in [SIAM J. Numer. Anal., 51 (2013), pp. 1235–1253] and [Internat. J. Numer. Methods Engrg., 94 (2013), pp. 128–149] for solving incompressible Stokes equations, are extended to three-dimensional problems. A new analysis of the condition number bound for using the Dirichlet preconditioner is given. The algorithm and analysis are valid for mixed finite elements with both continuous and discontinuous pressures. An advantage of this new analysis is that the numerous coarse level velocity components, required in the previous analysis to enforce the divergence free subdomain boundary velocity conditions, are no longer needed. This greatly reduces the size of the coarse level problem in the algorithm, especially for three-dimensional problems. The coarse level velocity space can be chosen as simple as for solving scalar elliptic problems corresponding to each velocity component. Both the Dirichlet and lumped preconditioners are analyzed using the same framework in this new analysis. Their condition number bounds are proved to be independent of the number of subdomains for fixed subdomain problem size. Numerical experiments in both two and three dimensions, using mixed finite elements with both continuous and discontinuous pressures, demonstrate the convergence rate of the algorithms.

Key words. domain decomposition, incompressible Stokes, FETI-DP, BDDC, divergence free

AMS subject classifications. 65F10, 65N30, 65N55

1. Introduction. Mixed finite elements are often used to solve incompressible Stokes and Navier-Stokes equations. Continuous pressures have been used in many mixed finite elements, e.g., the well known Taylor-Hood finite elements [27]. However, most domain decomposition methods require that the pressure be discontinuous, when they are used to solve the indefinite linear systems arising from such mixed finite element discretizations; see, e.g., [4, 5, 6, 10, 11, 16, 19, 21, 23, 24, 28, 29]. Several domain decomposition algorithms allow to use continuous pressures, e.g., Klawonn and Pavarino [16], Goldfeld [9], Šístek et. al. [25], Benhassine and Bendali [1], and Kim and Lee [15], even though no convergence rate analysis of those algorithm is known for the continuous pressure case.

Recently, the authors [20, 30] proposed and analyzed a FETI-DP (Dual-Primal Finite Element Tearing and Interconnecting) type domain decomposition algorithm for solving the incompressible Stokes equation in two dimensions. Both discontinuous and continuous pressures can be used in the mixed finite element discretization. In both cases, the indefinite system of linear equations can be reduced to a symmetric positive semi-definite system. Therefore, the preconditioned conjugate gradient method can be applied and a scalable convergence rate of the algorithm has been proved.

The lumped and Dirichlet preconditioners have been studied in [20] and [30], respectively. For the lumped preconditioner it was shown both experimentally and analytically in [20], that the coarse level space can be chosen the same as for solving scalar elliptic problems corresponding to each velocity component to achieve a scalable convergence rate. Similar observations for the lumped preconditioner have also been

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pointed out earlier by Kim and Lee [13, 14, 12, with Park], even though their studies are only for using discontinuous pressures.

For the Dirichlet preconditioner studied in [30], a distinctive feature is the application of subdomain discrete harmonic extensions in the preconditioner. In other existing FETI-DP and BDDC (Balancing Domain Decomposition by Constraints) algorithms, cf. [19, 21], subdomain discrete Stokes extensions have been used and the coarse level velocity space has to contain sufficient components to enforce divergence free subdomain boundary velocity conditions. Those complicated and numerous coarse level velocity components, especially for three-dimensional problems as discussed in [21], are not needed for the implementation of the Dirichlet preconditioner in [30]. But they are still required in [30] just for the analysis, where subdomain Stokes extensions were used, to obtain a scalable condition number bound.

In this paper, we provide a new analysis for the algorithms in [20, 30], which can analyze both the lumped and Dirichlet preconditioners in the same framework. It does not use any subdomain Stokes extensions and those additional coarse level components to enforce divergence free subdomain boundary velocity conditions are no longer needed. For both the lumped and Dirichlet preconditioners, the coarse level space can be chosen as simple as for solving scalar elliptic problems corresponding to each velocity component. This greatly simplifies the requirements on the coarse level space for the case of the Dirichlet preconditioner, especially in three dimensions. This paper is presented in the context of solving three-dimensional problems; the same approach can be applied to two-dimensional problems as well.

The remainder of this paper is organized as follows. The finite element discretization of the incompressible Stokes equation is introduced in Section 2. A domain decomposition approach is described in Section 3, and the system is reduced to a symmetric positive semi-definite problem in Section 4. A few preliminary results used in the condition number bound estimates are given in Section 5. The lumped and Dirichlet preconditioners are introduced in Section 6, and the condition number bounds of the preconditioned systems are established in Section 7. At the end, numerical results of solving the incompressible Stokes equation in both two and three dimensions are given in Section 8 to demonstrate the convergence rate of the algorithm.

2. Finite element discretization. We consider solving the following incompressible Stokes problem on a bounded, three-dimensional polyhedral domain $\Omega$ with a Dirichlet boundary condition,

$$
\begin{aligned}
-\Delta u^* + \nabla p^* &= f, & \text{in } \Omega, \\
-\nabla \cdot u^* &= 0, & \text{in } \Omega, \\
u^* &= u_{\partial\Omega}^*, & \text{on } \partial\Omega,
\end{aligned}
$$

(2.1)

where the boundary velocity $u_{\partial\Omega}^*$ satisfies the compatibility condition $\int_{\partial\Omega} u_{\partial\Omega}^* \cdot n = 0$. For simplicity, we assume that $u_{\partial\Omega}^* = 0$ without losing any generality.

The weak solution of (2.1) is given by: find $u^* \in (H^1_0(\Omega))^3 = \{ v \in (H^1(\Omega))^3 \mid v = 0 \text{ on } \partial\Omega \}$ and $p^* \in L^2(\Omega)$, such that

$$
\begin{aligned}
 a(u^*, v) + b(v, p^*) &= (f, v), & \forall v \in (H^1_0(\Omega))^3, \\
b(u^*, q) &= 0, & \forall q \in L^2(\Omega),
\end{aligned}
$$

(2.2)

where $a(u^*, v) = \int_\Omega \nabla u^* \cdot \nabla v$, $b(u^*, q) = -\int_\Omega (\nabla \cdot u^*) q$, $(f, v) = \int_\Omega f \cdot v$. We note that the solution of (2.2) is not unique, with the pressure $p^*$ different up to an additive constant.
A mixed finite element is used to solve (2.2). In this paper we apply a mixed finite element with continuous pressures, e.g., the Taylor-Hood type mixed finite elements. The same algorithm and analysis can be applied to mixed finite elements with discontinuous pressures as well; see [30]. Denote the velocity finite element space by $W \subset (H^1_0(\Omega))^3$, and the pressure finite element space by $Q \subset L^2(\Omega)$. The finite element solution $(u, p) \in W \oplus Q$ of (2.2) satisfies

\[(2.3) \quad \begin{bmatrix} A & B^T \\ B & 0 \end{bmatrix} \begin{bmatrix} u \\ p \end{bmatrix} = \begin{bmatrix} f \\ 0 \end{bmatrix},\]

where $A$, $B$, and $f$ represent respectively the restrictions of $a(\cdot, \cdot)$, $b(\cdot, \cdot)$ and $(f, \cdot)$ to the finite-dimensional spaces $W$ and $Q$. We use the same notation in this paper to represent both a finite element function and the vector of its nodal values.

The coefficient matrix in (2.3) is rank deficient even though $A$ is symmetric positive definite. $\text{Ker}(B^T)$, the kernel of $B^T$, contains all constant pressures in $Q$. $\text{Im}(B)$, the range of $B$, is orthogonal to $\text{Ker}(B^T)$ and consists of all vectors in $Q$ with zero average. For a general right-hand side vector $(f, g)$ in (2.3), the existence of solution requires that $g \in \text{Im}(B)$, i.e., $g$ has zero average; for the right-hand side given in (2.3), $g = 0$ and the solution always exists. When the pressure is considered in the quotient space $Q/\text{Ker}(B^T)$, the solution is unique. In this paper, when $q \in Q/\text{Ker}(B^T)$, we always assume that $q$ has zero average.

Let $h$ represent the characteristic diameter of the mixed elements. We assume that the mixed finite element space $W \times Q$, is inf-sup stable in the sense that there exists a positive constant $\beta$, independent of $h$, such that

\[(2.4) \quad \sup_{w \in W} \frac{\langle q, Bw \rangle^2}{\langle w, Aw \rangle} \geq \beta^2 \langle q, Zq \rangle, \quad \forall q \in Q/\text{Ker}(B^T),\]

cf. [3, Chapter III, §7]. Here, as always used in this paper, $\langle \cdot, \cdot \rangle$ represents the inner (or semi-inner) product of two vectors. The matrix $Z$ represents the mass matrix defined on the pressure finite element space $Q$, i.e., for any $q \in Q$, $\|q\|_{Z,2}^2 = \langle q, Zq \rangle$. It is easy to see, cf. [31, Lemma B.31], that $Z$ is spectrally equivalent to $h^3I$ for three-dimensional problems, i.e., there exist positive constants $c$ and $C$, such that

\[(2.5) \quad ch^3I \leq Z \leq Ch^3I,\]

where $I$ represents the identity matrix. Here, as in other places of this paper, $c$ and $C$ represent generic positive constants which are independent of $h$ and the subdomain diameter $H$ (described in the following section).

3. A non-overlapping domain decomposition approach. The domain $\Omega$ is decomposed into $N$ non-overlapping polyhedral subdomains $\Omega_i$, $i = 1, 2, \ldots, N$. Each subdomain is the union of a bounded number of elements, with the diameter of the subdomain in the order of $H$. We use $\Gamma$ to represent the subdomain interface which contains all the subdomain boundary nodes shared by neighboring subdomains; we assume that the subdomain meshes have matching nodes across $\Gamma$. $\Gamma$ is composed of subdomain faces, which are regarded as open subsets of $\Gamma$ shared by two subdomains, subdomain edges, which are regarded as open subsets of $\Gamma$ shared by more than two subdomains, and of the subdomain vertices, which are end points of edges.

The velocity and pressure finite element spaces $W$ and $Q$ are decomposed into

\[W = W_I \oplus W_\Gamma, \quad Q = Q_I \oplus Q_\Gamma,\]
where $W_I$ and $Q_I$ are direct sums of independent subdomain interior velocity spaces $W_I^{(i)}$, and interior pressure spaces $Q_I^{(i)}$, respectively, i.e.,

$$W_I = \bigoplus_{i=1}^{N} W_I^{(i)}, \quad Q_I = \bigoplus_{i=1}^{N} Q_I^{(i)}.$$ 

$W_\Gamma$ and $Q_\Gamma$ are subdomain interface velocity and pressure spaces, respectively. All functions in $W_\Gamma$ and $Q_\Gamma$ are continuous across $\Gamma$; their degrees of freedom are shared by neighboring subdomains.

To formulate the domain decomposition algorithm, we introduce a partially sub-assembled subdomain interface velocity space $\tilde{W}_\Gamma$,

$$\tilde{W}_\Gamma = W_\Delta \bigoplus W_{II} = \left( \bigoplus_{i=1}^{N} W_{\Delta}^{(i)} \right) \bigoplus W_{II}.$$ 

$W_{II}$ is the continuous, coarse level, primal velocity space which is typically spanned by subdomain vertex nodal basis functions, and/or by interface edge/face-cutoff functions with constant nodal values on each edge/face, or with values of positive weights on these edges/faces. The primal, coarse level velocity degrees of freedom are shared by neighboring subdomains. The complimentary space $W_\Delta$ is the direct sum of independent subdomain dual interface velocity spaces $W_\Delta^{(i)}$, which correspond to the remaining subdomain interface velocity degrees of freedom and are spanned by basis functions which vanish at the primal degrees of freedom. Thus, an element in $\tilde{W}_\Gamma$ typically has a continuous primal velocity component and a discontinuous dual velocity component.

It is well known that, for domain decomposition algorithms, the coarse space $W_{II}$ should be sufficiently rich to achieve a scalable convergence rate. On the other hand, a large coarse level problem will certainly degrade the parallel performance of the algorithm. Therefore it is important to keep the size of the coarse level problem as small as possible. When the Dirichlet preconditioner was used in the FETI-DP algorithm for solving the incompressible Stokes equation [19] and similarly in the BDDC algorithm [21], subdomain discrete Stokes extensions were used and $W_{II}$ has to contain sufficient subdomain interface components such that functions in $W_\Delta$ have zero flux across the subdomain boundaries. Such requirements lead to a large coarse level velocity space, especially for three-dimensional problems, cf. [21].

In [30], a FETI-DP type algorithm is proposed for solving two-dimensional incompressible Stokes problems. A distinctive feature of the Dirichlet preconditioner used in that algorithm is the application of subdomain discrete harmonic extensions, instead of subdomain discrete Stokes extensions. As a result, the divergence free subdomain boundary velocity conditions are not needed in that algorithm. However, the analysis, given in [30] for the Dirichlet preconditioner, still uses subdomain Stokes extensions and requires the same type coarse level velocity space as discussed in [21] to establish a scalable condition number bound estimate. In this paper, a new analysis is offered and it is sufficient for $W_{II}$ to be spanned just by the subdomain vertex nodal basis functions and subdomain edge-cutoff functions corresponding to each velocity component, as for solving three-dimensional scalar elliptic problems, cf. [31, Section 6.4.2].

The functions $w_\Delta$ in $W_\Delta$ are in general not continuous across $\Gamma$. To enforce their
continuity, we define a Boolean matrix $B_\Delta$ of the form

$$B_\Delta = \begin{bmatrix} B_\Delta^{(1)} & B_\Delta^{(2)} & \cdots & B_\Delta^{(N)} \end{bmatrix},$$

constructed from $\{0, 1, -1\}$. On each row of $B_\Delta$, there are only two nonzero entries, 1 and $-1$, corresponding to one velocity degree of freedom shared by two neighboring subdomains, such that for any $w_\Delta$ in $W_\Delta$, each row of $B_\Delta w_\Delta = 0$ implies that these two degrees of freedom from the two neighboring subdomains be the same. We note that, in three dimensions, a velocity degree of freedom on a subdomain edge is shared by more than two subdomains, e.g., by four subdomains. In this case, a minimum of three continuity constraints can be applied to enforce the continuity of this velocity degree of freedom among the four subdomains, which corresponds to the use of non-redundant Lagrange multipliers. In this paper, the fully redundant Lagrange multipliers are used, which means, e.g., for a subdomain edge velocity degree of freedom shared by four subdomains, six Lagrange multipliers are used to enforce all the six possible continuity constraints among them, cf. [31, Section 6.3.1].

We denote the range of $B_\Delta$ applied on $W_\Delta$ by $\Lambda$, the vector space of the Lagrange multipliers. Solving the original fully assembled linear system (2.3) is then equivalent to: find $(u_I, p_I, u_\Delta, u_{II}, p_{II}, \lambda) \in W_I \oplus Q_I \oplus W_\Delta \oplus W_{II} \oplus Q_{II} \oplus \Lambda$, such that

$$\begin{bmatrix} \begin{bmatrix} A_{II} & B_{II}^T & A_{I\Delta} & A_{II} & B_{I\Delta}^T & 0 \\ B_{II} & 0 & B_{I\Delta} & B_{II} & 0 & 0 \\ A_{I\Delta} & B_{I\Delta}^T & A_{\Delta\Delta} & A_{II} & B_{\Delta\Delta} & B_{I\Delta}^T \\ A_{II} & B_{II}^T & A_{I\Delta} & A_{II} & B_{II} & 0 \\ B_{II} & 0 & B_{I\Delta} & B_{II} & 0 & 0 \\ 0 & 0 & B_\Delta & 0 & 0 & 0 \end{bmatrix} \end{bmatrix} \begin{bmatrix} u_I \\ p_I \\ u_\Delta \\ u_{II} \\ p_{II} \\ \lambda \end{bmatrix} = \begin{bmatrix} f_I \\ 0 \\ f_\Delta \\ f_{II} \\ 0 \\ 0 \end{bmatrix},$$

where the sub-blocks in the coefficient matrix represent the restrictions of $A$ and $B$ in (2.3) to appropriate subspaces. The leading three-by-three block can be made block diagonal with each diagonal block corresponding to one subdomain.

The coefficient matrix in (3.1) is singular. The trivial null space vectors are those with $\lambda$ in the null space of $B_\Delta^T$ and other components zero. Such singularity, due to the rank deficiency of $B_\Delta$, needs not to be worried, since the Lagrange multiplier vector $\lambda$ will be confined in $\Lambda$, the range of $B_\Delta$. The only meaningful basis vector in the null space of (3.1) corresponds to the one-dimensional null space of the original incompressible Stokes system (2.3), and is specified in the following lemmas.

We first need to introduce a positive scaling factor $\delta^I(x)$ for each node $x$ on $\Gamma$. Let $N_x$ be the number of subdomains sharing $x$, and we define $\delta^I(x) = 1/N_x$. Given such scaling factors at the subdomain interface nodes, we can define a scaled operator $B_{\Delta,D}$. We note that each row of $B_\Delta$ has only two nonzero entries, 1 and $-1$, connecting two neighboring subdomains sharing a node $x$ on $\Gamma$. Multiplying each entry by the scaling factor $\delta^I(x)$ gives us $B_{\Delta,D}$. Namely

$$B_{\Delta,D} = \begin{bmatrix} D_\Delta B_{\Delta}^{(1)} & D_\Delta B_{\Delta}^{(2)} & \cdots & D_\Delta B_{\Delta}^{(N)} \end{bmatrix},$$

where $D_\Delta$ is a diagonal matrix and contains $\delta^I(x)$ on its diagonal. We also see from the definition of $B_{\Delta,D}$ that the scalings on all the Lagrange multipliers related to
on the same subdomain interface node are the same, from which we have the following lemma.

**Lemma 3.1.** The null space of $B^T_{\Delta}$ is the same as the null space of $B^T_{\Delta,D}$; the range of $B_\Delta$ is the same as the range of $B_{\Delta,D}$.

**Lemma 3.2.** For any $\lambda \in \Lambda$, $B_{\Delta} B^T_{\Delta,D} \lambda = B_{\Delta,D} B^T_{\Delta} \lambda = \lambda$.

**Proof:** The equality $B_{\Delta} B^T_{\Delta,D} \lambda = \lambda$ can be found directly at [31, Page 175]. A similar approach is used here to prove $B_{\Delta,D} B^T_{\Delta} \lambda = \lambda$. Since $\Lambda$ is the range of $B_\Delta$ applied on $W_\Delta$, let us denote $\lambda = B_\Delta w_\Delta$ for a certain $w_\Delta \in W_\Delta$. Then just look at one entry of $\lambda$, e.g., the Lagrange multiplier connecting two neighboring subdomain velocity degrees of freedom, represented by $w^{(i)}(x)$ and $w^{(j)}(x)$, of $w_\Delta$ at a subdomain interface boundary node $x$ between $\Omega_i$ and $\Omega_j$. The value of that Lagrange multiplier equals $w^{(i)}(x) - w^{(j)}(x)$ (or $w^{(j)}(x) - w^{(i)}(x)$ depending on the choice of signs in $B_\Delta$).

On the other hand the corresponding entry of $B_{\Delta,D} B^T_{\Delta} B_\Delta w_\Delta$ equals

$$\delta^i(x) \left( B^T_{\Delta} B_\Delta w_\Delta \right)^{(i)}(x) - \delta^j(x) \left( B^T_{\Delta} B_\Delta w_\Delta \right)^{(j)}(x)$$

$$= \delta^i(x) \sum_{k \in N_x} \left( w^{(i)}(x) - w^{(k)}(x) \right) - \delta^j(x) \sum_{k \in N_x} \left( w^{(j)}(x) - w^{(k)}(x) \right)$$

$$= \sum_{k \in N_x} \delta^i(x) \left( w^{(i)}(x) - w^{(j)}(x) \right) = w^{(i)}(x) - w^{(j)}(x),$$

where in the last step, we used the fact that $\sum_{k \in N_x} \delta^i(x) = 1$. \hfill \qed

**Lemma 3.3.** Let $1_{p_l} \in Q_1, 1_{p_r} \in Q_1$ represent vectors with value 1 on each entry. Then

$$[B^T_{\Delta} \ B^T_{\Gamma\Delta}] \begin{bmatrix} 1_{p_l} \\ 1_{p_r} \end{bmatrix} = B^T_{\Delta} \lambda,$$

where

$$\lambda = B_{\Delta,D} [B^T_{\Gamma\Delta} \ B^T_{\Gamma\Delta}] \begin{bmatrix} 1_{p_l} \\ 1_{p_r} \end{bmatrix} \in \Lambda.$$

**Proof:** The left side of (3.2) contains face integrals of the normal component of the dual subdomain interface velocity finite element basis functions across the subdomain interface. For a face velocity degree of freedom, which is shared by two neighboring subdomains, the face integrals of their normal components on the two neighboring subdomains are negative of each other, since their normal directions are opposite. This pair of opposite values can then be represented by the product of $B^T_{\Delta}$ and a Lagrange multiplier with value equal to the face integral of the corresponding basis function.

Now we consider a subdomain edge velocity degree of freedom, which is shared by more than two subdomains, e.g., by four subdomains $\Omega_i, \Omega_j, \Omega_k,$ and $\Omega_l$. A two-dimensional projection of such an edge node is shown in Figure 3.1, where the third direction points outward directly and the dashed lines represent the element boundaries. $F_{ij}, F_{jk}, F_{kl}$ and $F_{li}$ represent the common element faces connected to this edge node on the subdomain interfaces, e.g., $F_{ij}$ represents the element faces sharing this node between $\Omega_i$ and $\Omega_j$, while the elements on $\Omega_i$ and $\Omega_k$ have no common face. Denote the integration of the normal component of this velocity basis function on $F_{ij}, F_{jk}, F_{kl}$ and $F_{li}$ by $I_{ij}, I_{jk}, I_{kl}$ and $I_{li}$, respectively, with a certain
chosen normal direction for each element face. Then the entries of the left side vector in (3.2) corresponding to this edge velocity degree of freedom on the four subdomains $\Omega_i, \Omega_j, \Omega_k, \Omega_l$, can be represented by $I_{ij} + I_{li}, -I_{ij} + I_{jk}, -I_{jk} - I_{kl}, \text{ and } I_{kl} - I_{li}$, respectively, where the choice of positive and negative signs are due to the fact that the normal direction to the same element faces shared by two neighboring subdomains are opposite of each other. Take $I_{ij}, I_{jk}, I_{kl}, I_{li}$ as the four Lagrange multiplier values as illustrated in Figure 3.1. Then the four subdomain face integral values $I_{ij} + I_{li}, -I_{ij} + I_{jk}, -I_{jk} - I_{kl}, \text{ and } I_{kl} - I_{li}$, can be represented as the product of corresponding $B_T^\Delta$ with a Lagrange multiplier vector containing these four Lagrange multiplier values and zero elsewhere.

The above has just shown that the left side of $(3.2)$ can be represented by the product of $B_T^\Delta$ with a Lagrange multiplier vector $\lambda$ if $\lambda$ is not in $\Lambda$, i.e., not in the range of $B_D$, it can always be written as the sum of its components in $\Lambda$ and in the null space of $B_T^\Delta$. Then we just take its component in $\Lambda$ as $\lambda$, which does not change the product $B_T^\Delta \lambda$. By multiplying $B_D, D$ to both sides of (3.2) and using Lemma 3.2, we have (3.3).

**Lemma 3.4.** The basis vector in the null space of (3.1), corresponding to the one-dimensional null space of the original incompressible Stokes system (2.3), is

\begin{equation}
(3.4) \quad \begin{pmatrix} 0, & 1_{p_l}, & 0, & 0, & 1_{p_r}, & -B_D, D \left[B_T^\Delta \right] \left[ B_D^T \right] B_T^\Delta \left[ 1_{p_l} \right] \\
\end{pmatrix}.
\end{equation}

**Proof:** Since the null space of (2.3) consists of all constant pressures, substituting the vector (3.4) into (3.1) gives zero blocks on the right-hand side, except at the third block where

\begin{equation}
(3.5) \quad f_\Delta = \begin{bmatrix} B_T^\Delta & B_T^\Gamma \end{bmatrix} \begin{bmatrix} 1_{p_l} \\
1_{p_r} \end{bmatrix} - B_D, D \left[B_T^\Delta \right] \left[ B_D^T \right] B_T^\Delta \left[ 1_{p_l} \right],
\end{equation}

which also equals zero from (3.2) and (3.3) in Lemma 3.3.

**4. A reduced symmetric positive semi-definite system.** The system (3.1) can be reduced to a Schur complement problem for the variables $(p_r, \lambda)$. Since the leading four-by-four block of the coefficient matrix in (3.1) is invertible, the variables
can be eliminated and we obtain

\begin{equation}
G \begin{bmatrix} p_T \\ \lambda \end{bmatrix} = g,
\end{equation}

where

\begin{equation}
G = B_C \tilde{A}^{-1} B_C^T, \quad g = B_C \tilde{A}^{-1} \begin{bmatrix} f_I \\ 0 \\ f_\Delta \\ f_{II} \end{bmatrix},
\end{equation}

with

\begin{equation}
\tilde{A} = \begin{bmatrix} A_{II} & B_{II}^T & A_{I\Delta} & A_{I\Pi} \\ B_{II} & 0 & B_{I\Delta} & B_{I\Pi} \\ A_{I\Delta} & B_{I\Delta}^T & A_{\Delta\Delta} & A_{\Delta\Pi} \\ A_{I\Pi} & B_{I\Pi}^T & A_{\Pi\Delta} & A_{\Pi\Pi} \end{bmatrix}
\end{equation}

and $B_C = \begin{bmatrix} B_{II} & 0 & B_{I\Delta} & B_{I\Pi} \\ 0 & 0 & B_{\Delta} & 0 \end{bmatrix}$.

We can see that $-G$ is the Schur complement of the coefficient matrix of (3.1) with respect to the last two row blocks, i.e.,

\[
\begin{bmatrix} I & 0 \\ -B_C \tilde{A}^{-1} & I \end{bmatrix} \begin{bmatrix} \tilde{A} & B_C^T \\ B_C & 0 \end{bmatrix} \begin{bmatrix} I & -\tilde{A}^{-1} B_C^T \\ 0 & I \end{bmatrix} = \begin{bmatrix} \tilde{A} & 0 \\ 0 & -G \end{bmatrix}.
\]

From the Sylvester law of inertia, namely, the number of positive, negative, and zero eigenvalues of a symmetric matrix is invariant under a change of coordinates, we can see that the number of zero eigenvalues of $G$ is the same as the number of zero eigenvalues (with multiplicity counted) of the original coefficient matrix of (3.1), and all other eigenvalues of $G$ are positive. Therefore $G$ is symmetric positive semi-definite.

The basis vectors of the null space of $G$ also inherit those from the null space of (3.1), and the only interesting basis vector is

\begin{equation}
\begin{bmatrix} 1_{pr} \\ -B_{\Delta,D} \end{bmatrix} \begin{bmatrix} B_{I\Delta}^T & B_{I\Delta}^T \end{bmatrix} \begin{bmatrix} 1_{pr} \\ 1_{pr} \end{bmatrix},
\end{equation}

which is derived from Lemma 3.4. The other null space vectors of $G$ are all vectors with $\lambda$ in the null space of $B_\Delta^T$ and $p_T = 0$. The range of $G$ contains all vectors orthogonal to those null vectors. Denote $X = Q_I \oplus \Lambda$, where, as defined earlier, $\Lambda$ is the range of $B_\Delta$. Then the range of $G$, denoted by $R_G$, is the subspace of $X$ orthogonal to (4.4), i.e.,

\begin{equation}
R_G = \left\{ \begin{bmatrix} g_{pr} \\ g_\lambda \end{bmatrix} \in X \mid g_{pr}^T 1_{pr} - g_\lambda^T \begin{bmatrix} B_{\Delta,D} & B_{I\Delta}^T \end{bmatrix} \begin{bmatrix} 1_{pr} \\ 1_{pr} \end{bmatrix} = 0 \right\}.
\end{equation}

The restriction of $G$ to its range $R_G$ is positive definite. The fact that the solution of (3.1) always exists for any given $(f_I, f_\Delta, f_{II})$ on the right-hand side implies that the solution of (4.1) exists for any $g$ defined by (4.2). Therefore $g \in R_G$. When the conjugate gradient method (CG) is applied to solve (4.1) with zero initial guess, all the iterates are in the Krylov subspace generated by $G$ and $g$, which is also a
subspace of $R_G$, and where the CG cannot break down. After obtaining $(p_{\Gamma}, \lambda)$ from solving (4.1), the other components $(u_I, p_I, u_{\Delta}, u_{\Pi})$ in (3.1) are obtained by back substitution.

In the rest of this section, we discuss the implementation of multiplying $G$ by a vector. The main operation is the product of $\tilde{A}^{-1}$ with a vector, cf. (4.2). We denote

$$A_{rr} = \begin{bmatrix} A_{II} & B_{II}^T & A_{I\Delta} \\ B_{II} & 0 & B_{I\Delta} \\ A_{\Delta I} & B_{I\Delta}^T & A_{\Delta\Delta} \end{bmatrix}, \quad A_{rr}^{-1} = \begin{bmatrix} f_I \\ 0 \\ f_{\Delta} \end{bmatrix},$$

and define the Schur complement

$$S_{\Pi} = A_{\Pi\Pi} - A_{\Pi r} A_{rr}^{-1} A_{r\Pi},$$

which is symmetric positive definite from the Sylvester law of inertia. $S_{\Pi}$ defines the coarse level problem in the algorithm. The product

$$\begin{bmatrix} A_{II} & B_{II}^T & A_{I\Pi} \\ B_{II} & 0 & B_{I\Pi} \\ A_{\Pi I} & B_{I\Pi}^T & A_{\Pi\Pi} \end{bmatrix}^{-1} \begin{bmatrix} f_I \\ 0 \\ f_{\Pi} \end{bmatrix}$$

can then be represented by

$$\begin{bmatrix} A_{rr}^{-1} f_r \\ 0 \end{bmatrix} + \begin{bmatrix} -A_{rr}^{-1} A_{r\Pi} \\ I_{\Pi} \end{bmatrix} S_{\Pi}^{-1} \begin{bmatrix} f_{\Pi} - A_{\Pi r} A_{rr}^{-1} f_r \end{bmatrix},$$

which requires solving the coarse level problem once and independent subdomain Stokes problems with Neumann type boundary conditions twice.

5. Preliminary results. Denote

$$(5.1) \quad \tilde{W} = W_I \bigoplus \tilde{W}_I = W_I \bigoplus W_\Delta \bigoplus W_{\Pi}.$$ 

For any $w$ in $\tilde{W}$, denote its restriction to subdomain $\Omega_i$ by $w^{(i)}$. A subdomain-wise $H^1$-seminorm can be defined for functions in $\tilde{W}$ by

$$|w|_{H^1}^2 = \sum_{i=1}^N |w^{(i)}|_{H^1(\Omega_i)}^2.$$ 

We also define

$$\tilde{V} = W_I \bigoplus Q_I \bigoplus W_\Delta \bigoplus W_{\Pi},$$

and its subspace

$$(5.2) \quad \tilde{V}_0 = \left\{ v = (w_I, p_I, w_\Delta, w_{\Pi}) \in \tilde{V} \mid B_{II} w_I + B_{I\Delta} w_\Delta + B_{I\Pi} w_{\Pi} = 0 \right\}.$$
For any \( v = (w_I, p_I, w_\Delta, w_{II}) \in \tilde{V}_0 \), we have \( w = (w_I, w_\Delta, w_{II}) \in \tilde{W} \). Then

\[
\langle v, v \rangle_\tilde{A} = \begin{bmatrix} w_I \\ w_\Delta \\ w_{II} \end{bmatrix}^T \begin{bmatrix} A_{II} & A_{I\Delta} & A_{I\Pi} \\ A_{\Delta I} & A_{\Delta\Delta} & A_{\Delta\Pi} \\ A_{\Pi I} & A_{\Pi\Delta} & A_{\Pi\Pi} \end{bmatrix} \begin{bmatrix} w_I \\ w_\Delta \\ w_{II} \end{bmatrix}
\]

\[(5.3) = \sum_{i=1}^{N} \begin{bmatrix} w^{(i)}_I \\ w^{(i)}_\Delta \\ w^{(i)}_{II} \end{bmatrix}^T \begin{bmatrix} A_{II}^{(i)} & A_{I\Delta}^{(i)} & A_{I\Pi}^{(i)} \\ A_{\Delta I}^{(i)} & A_{\Delta\Delta}^{(i)} & A_{\Delta\Pi}^{(i)} \\ A_{\Pi I}^{(i)} & A_{\Pi\Delta}^{(i)} & A_{\Pi\Pi}^{(i)} \end{bmatrix} \begin{bmatrix} w^{(i)}_I \\ w^{(i)}_\Delta \\ w^{(i)}_{II} \end{bmatrix} = \sum_{i=1}^{N} \left| w^{(i)}_I \right|^2 \]

where the superscript \( (i) \) is used to represent the restrictions of corresponding vectors and matrices to subdomain \( \Omega_i \). We can see from (5.3) that for any \( v \in \tilde{V}_0 \), the value \( \langle v, v \rangle_\tilde{A} \) is independent of its pressure component \( p_I \). \( \langle \cdot, \cdot \rangle_\tilde{A} \) defines a semi-inner product on \( \tilde{V}_0 \); \( \langle v, v \rangle_\tilde{A} = 0 \) if and only if the velocity component of \( v \) is constant on \( \Omega \) and is in fact zero due to the zero boundary condition on \( \partial \Omega \), while its pressure component can be arbitrary.

Denote

\[(5.4) \quad \tilde{B} = \begin{bmatrix} B_{II} & B_{I\Delta} & B_{I\Pi} \\ B_{\Gamma I} & B_{\Gamma\Delta} & B_{\Gamma\Pi} \end{bmatrix}, \]

cf. (3.1). The following lemma on the stability of \( \tilde{B} \) can be found at [20, Lemma 5.1].

**Lemma 5.1.** For any \( w \in \tilde{W} \) and \( q \in Q \), \( \langle \tilde{B}w, q \rangle \leq |w|_{H^1} \| q \|_{L^2} \).

The following lemma will also be used and can be found at [10, Lemma 2.3].

**Lemma 5.2.** Let \( (u, p) \in W \oplus Q \) satisfy

\[(5.5) \quad \begin{bmatrix} A & B^T \\ B & 0 \end{bmatrix} \begin{bmatrix} u \\ p \end{bmatrix} = \begin{bmatrix} f \\ g \end{bmatrix}, \]

where \( A \) and \( B \) are as in (2.3), \( f \in \tilde{W} \), and \( g \in \text{Im}(B) \subset Q \). Let \( \beta \) be the inf-sup constant specified in (2.4). Then

\[\|u\|_A \leq \|f\|_{A^{-1}} + \frac{1}{\beta} g \|_{Z^{-1}},\]

where \( Z \) is the mass matrix defined in Section 2.

**6. Jump operators and preconditioners.** We first define certain jump operators across the subdomain interface \( \Gamma \), which will be used for the analysis of the preconditioners.

Denote the restriction operator from \( \tilde{V} \) onto \( W_\Delta \) by \( \tilde{R}_\Delta \), i.e., for any \( v = (w_I, p_I, w_\Delta, w_{II}) \in \tilde{V} \), \( \tilde{R}_\Delta v = w_\Delta \). Define \( P_{D,L} : \tilde{V} \to \tilde{V} \), by

\[P_{D,L} = \tilde{R}_\Delta^T B_{\Delta\Delta}^T B_\Delta \tilde{R}_\Delta.\]

Following this definition, given any \( v = (w_I, p_I, w_\Delta, w_{II}) \in \tilde{V} \), the dual velocity component of \( P_{D,L} v \), on any subdomain interface node \( x \) in subdomain \( \Omega_i \), is given
by, cf. [31, Equation (6.70)],
\[
\left(R_\Delta (P_{D,L}v)\right)^{(i)}(x) = \sum_{j \in \mathcal{N}_x} \delta'(x) \left(w_\Delta^{(j)}(x) - w_\Delta^{(j)}(x)\right),
\]
which represents the so-called jump of the dual velocity component \(w_\Delta\) across the subdomain interface \(\Gamma\). All other components of \(P_{D,L}v\) equal zero. We also have
\[
\langle P_{D,L}v, P_{D,L}v \rangle_{\tilde{A}} = (H_{\Delta,D} B_{\Delta,D} \tilde{R}_\Delta v)^T \tilde{A} (H_{\Delta,D} B_{\Delta,D} \tilde{R}_\Delta v) = \langle B_{\Delta,D}^T B_{\Delta,D} w_\Delta, B_{\Delta,D}^T B_{\Delta,D} w_\Delta \rangle_{A_{\Delta \Delta}}.
\]
(6.1)
Together with (5.3), we have the following lemma, which can be found at [22, Section 6.1].

**Lemma 6.1.** There exists a constant \(C\) and a function \(\Phi_L(H/h)\), such that for all \(v \in \tilde{V}_0\), \(\langle P_{D,L}v, P_{D,L}v \rangle_{\tilde{A}} \leq C \Phi_L(H/h) \langle v, v \rangle_{\tilde{A}}\). Here, \(\Phi_L(H/h) = (H/h)(1 + \log(H/h))\), when the coarse level space is spanned by the subdomain vertex nodal basis functions and subdomain edge-cutoff functions corresponding to each velocity component.

When applying \(P_{D,L}\) to a vector, the jump of the dual subdomain interface velocities is extended by zero to the interior of subdomains. To improve the stability of the jump operator, the jump can be extended to the interior of subdomains by subdomain discrete harmonic extension. We define a Schur complement operator \(H_\Delta^{(i)} : w_\Delta^{(i)} \to w_\Delta^{(i)}\) by, for any \(u_\Delta^{(i)} \in w_\Delta^{(i)}\),
\[
A^{(i)}_{II} A^{(i)}_{\Delta \Delta} \begin{bmatrix} A^{(i)}_{II} \ A^{(i)}_{\Delta \Delta} \ A^{(i)}_{II} \ A^{(i)}_{\Delta \Delta} \end{bmatrix} \begin{bmatrix} u^{(i)}_{II} \ u^{(i)}_{\Delta \Delta} \ u^{(i)}_{II} \ u^{(i)}_{\Delta \Delta} \end{bmatrix} = \begin{bmatrix} 0 \ H_\Delta^{(i)} u^{(i)}_{\Delta \Delta} \end{bmatrix}.
\]
(6.2)
To multiply \(H_\Delta^{(i)}\) by a vector \(u_\Delta^{(i)}\), a subdomain elliptic problem on \(\Omega_i\) with given boundary velocity \(u_\Delta^{(i)}\) and \(u^{(i)}_{II} = 0\) needs to be solved.

Using \(H_\Delta^{(i)}\), we define the second jump operator \(P_{D,D} : \tilde{V} \to \tilde{V}\), as follows: for any given \(v = (w_I, p_I, w_\Delta, w_{II}) \in \tilde{V}\), the subdomain interior velocity part of \(P_{D,D}v\) on each subdomain \(\Omega_i\) is taken as \(u^{(i)}_{II}\) in the solution of (6.2), with given subdomain boundary velocity \(u^{(i)}_{\Delta} = B_{\Delta,D}^T B_{\Delta,D} w_\Delta\). Here \(B_{\Delta,D}^T\) represents restriction of \(B_{\Delta,D}\) on subdomain \(\Omega_i\) and is a map from \(\Lambda\) to \(w_\Delta^{(i)}\). The other components of \(P_{D,D}v\) are kept zero. Therefore
\[
\langle P_{D,D}v, P_{D,D}v \rangle_{\tilde{A}} = \sum_{i=1}^N \begin{bmatrix} u^{(i)}_{II} \ u^{(i)}_{\Delta \Delta} \end{bmatrix}^T \begin{bmatrix} A^{(i)}_{II} & A^{(i)}_{\Delta \Delta} & A^{(i)}_{II} & A^{(i)}_{\Delta \Delta} \end{bmatrix} \begin{bmatrix} u^{(i)}_{II} \ u^{(i)}_{\Delta \Delta} \ u^{(i)}_{II} \ u^{(i)}_{\Delta \Delta} \end{bmatrix} = \sum_{i=1}^N u^{(i)}_{II} A^{(i)}_{\Delta \Delta} u^{(i)}_{\Delta \Delta} \sum_{i=1}^N u^{(i)}_{II} A^{(i)}_{\Delta \Delta} u^{(i)}_{\Delta \Delta}
\]
(6.3)
\[
\leq C \Phi_D(H/h) \sum_{i=1}^N \begin{bmatrix} w^{(i)}_I \ w^{(i)}_{II} \end{bmatrix}^2_{H^{1/2}(\partial \Omega_i)} \leq C \Phi_D(H/h) \sum_{i=1}^N \begin{bmatrix} w^{(i)}_I \ w^{(i)}_{II} \end{bmatrix}^2_{H^{1/2}(\partial \Omega_i)} = C \Phi_D(H/h) \sum_{i=1}^N \begin{bmatrix} w^{(i)}_I \ w^{(i)}_{II} \end{bmatrix}^2_{H^1(\Omega_i)} = C \Phi_D(H/h) \sum_{i=1}^N \begin{bmatrix} w^{(i)}_I \ w^{(i)}_{II} \end{bmatrix}^2_{H^1(\Omega_i)},
\]
The first inequality in (6.3) is a well established result, cf., [31, Lemma 6.36]. Since for any \( v \in \tilde{V}_0 \), \( \langle v, v \rangle_{\tilde{A}} = |w|_{H^1(\Omega)}^2 \), cf. (5.3), we have the following lemma.

**Lemma 6.2.** There exists a constant \( C \) and a function \( \Phi_D(H/h) \), such that for all \( v \in \tilde{V}_0 \), \( \langle P_D,D v, P_D,D v \rangle_{\tilde{A}} \leq C \Phi_D(H/h) \langle v, v \rangle_{\tilde{A}} \). Here \( \Phi_D(H/h) = (1 + \log (H/h))^2 \), when the coarse level space is spanned by the subdomain vertex nodal basis functions and subdomain edge-cutoff functions corresponding to each velocity component.

**Remark 6.3.** We note that the coarse spaces in Lemmas 6.1 and 6.2 are not necessarily the coarse space of minimal size to achieve a scalable bound \( \Phi_L(H/h) \) and \( \Phi_D(H/h) \), respectively; several other choices can be found at [31, Section 6.4.2]. As for the edge-cutoff function in Lemma 6.2, the proof given above requires the edge-cut-off function of each velocity component. Meanwhile, enriching the coarse space may improve those bounds. For the case of \( \Phi_L(H/h) \) in Lemma 6.1, it has been proved that the bound can be improved to \( \Phi_L \), for two-dimensional problem with additional edge-cut-off functions in the coarse space [22] and for three dimensions with additional face-cut-off functions in the coarse space [13, 14]. But, to the best of our knowledge, no such improvements have been proved in literature for the case of \( \Phi_D(H/h) \) in Lemma 6.2 in either two or three dimensions.

To introduce the preconditioners, we write \( G \), defined in (4.2) and (4.3), in a two-by-two block structure. Denote the first row of \( B_C \) by

\[
\tilde{B}_T = \begin{bmatrix} B_{T} & 0 & B_{T\Delta} & B_{\Gamma \Gamma \Gamma} \end{bmatrix},
\]

and note that \( \tilde{R}_{\Delta} \) is the restriction operator from \( \tilde{V} \) onto \( W_{\Delta} \). Then \( G \) can be written as

\[
G = \begin{bmatrix} G_{\lambda p \lambda p} & G_{\lambda p \lambda} \\ G_{\lambda p} & G_{\lambda \lambda} \end{bmatrix},
\]

where

\[
G_{\lambda p \lambda p} = \tilde{B}_T \tilde{A}^{-1} \tilde{B}_T^T, \quad G_{\lambda p \lambda} = \tilde{B}_T \tilde{A}^{-1} \tilde{R}_{\Delta}^T B_{\Delta},
\]

\[
G_{\lambda p} = B_{\Delta} \tilde{R}_{\Delta} \tilde{A}^{-1} \tilde{B}_T, \quad G_{\lambda \lambda} = B_{\Delta} \tilde{R}_{\Delta} \tilde{A}^{-1} \tilde{R}_{\Delta}^T B_{\Delta}.
\]

We consider a block diagonal preconditioner for (4.1). As for two-dimensional problems, the first diagonal block \( G_{\lambda p \lambda p} \) of \( G \) can be shown spectrally equivalent to \( h^3 I_{p\lambda} \), where \( I_{p\lambda} \) is the identity matrix of the same dimension as \( G_{\lambda p \lambda p} \); see [20, 30]. Therefore, in the following block diagonal preconditioners, the inverse of \( G_{\lambda p \lambda p} \) is approximated by \( \alpha h^{-3} I_{p\lambda} \). Here \( \alpha \) is a given constant. We will show in the next section that \( \alpha \) has only a minor effect on the condition number bound of the preconditioned operator and its value is typically taken as 1, cf. Remark 7.7. We introduce \( \alpha \) in the preconditioner just for the convenience in the numerical experiments to demonstrate the convergence rates of the proposed algorithm.

The inverse of the second diagonal block \( B_{\Delta} \tilde{R}_{\Delta} \tilde{A}^{-1} \tilde{R}_{\Delta}^T B_{\Delta} \), can be approximated by the lumped block

\[
M_{\lambda \lambda}^{-1} = B_{\Delta,D} \tilde{R}_{\Delta} \tilde{A} \tilde{R}_{\Delta}^T B_{\Delta,D}^T.
\]

This leads to the following lumped preconditioner for solving (4.1)

\[
M_{\lambda \lambda}^{-1} = \begin{bmatrix} \alpha h^{-3} I_{p\lambda} & 0 \\ 0 & M_{\lambda \lambda}^{-1} \end{bmatrix}.
\]
Applying subdomain discrete harmonic extensions in the preconditioning step, we have the following Dirichlet preconditioner

\begin{equation}
M_D^{-1} = \begin{bmatrix}
\alpha h^{-3} I_{p_F} & M_{\lambda,D}^{-1}
\end{bmatrix},
\end{equation}

where

\begin{equation}
M_{\lambda,D}^{-1} = B_{\Delta,D} H_{\Delta} B_{\Delta,D}^T.
\end{equation}

Here $H_{\Delta} : W_{\Delta} \rightarrow W_{\Delta}$ represents the direct sum of $H_{\Delta}^{(i)}$, $i = 1, \ldots, N$, defined in (6.2).

Remark 6.4. The lumped and Dirichlet preconditioners (6.6) and (6.7) are of the same feature as those in the original FETI-DP algorithms [7, 8, 17] developed for solving elliptic problems, e.g., the Laplace equation. The only difference here is the additional diagonal block in the two preconditioners corresponding to the pressure variables; removing all the pressure variables and their related blocks in the algorithms studied in this paper will just lead to the original FETI-DP algorithms.

We can see from Lemma 3.1 that both $M_{\lambda,L}^{-1}$ and $M_{\lambda,D}^{-1}$ are symmetric positive definite when restricted on $\Lambda$. Therefore both the lumped and Dirichlet preconditioners $M_L^{-1}$ and $M_D^{-1}$ are symmetric positive definite in the range of $G$.

7. Condition number bounds. In the following, we use the same framework to establish the condition number bounds for both the lumped and Dirichlet preconditioned operators $M_L^{-1}G$ and $M_D^{-1}G$. Let $M^{-1}$, $M_{\lambda}^{-1}$, $P_D$, and $\Phi$ represent both $M_L^{-1}$, $M_{\lambda,L}^{-1}$, $P_{D,L}$, $\Phi_L$, for the lumped preconditioner case, and $M_D^{-1}$, $M_{\lambda,D}^{-1}$, $P_{D,D}$, $\Phi_D$, for the Dirichlet preconditioner case, respectively, when they apply in the proofs.

When the conjugate gradient method is applied to solving the preconditioned system

\begin{equation}
M^{-1}Gx = M^{-1}g,
\end{equation}

with zero initial guess, all iterates belong to the Krylov subspace generated by the operator $M^{-1}G$ and the vector $M^{-1}g$, which is a subspace of the range of $M^{-1}G$. We denote the range of $M^{-1}G$ by $R_{M^{-1}G}$ and note that both preconditioners are symmetric positive definite in the range of $G$. We have the following lemma, cf. [30, Lemma 6].

Lemma 7.1. The conjugate gradient method applied to solving (7.1) with zero initial guess cannot break down.

Proof: We just need to show that for any $0 \neq x \in R_{M^{-1}G}$, $\langle x, Gx \rangle \neq 0$, i.e., to show $Gx \neq 0$. Let $0 \neq x = M^{-1}Gy$, for a certain $y \in X$ and $y \neq 0$. Then $Gx = GM^{-1}Gy$, which cannot be zero since $Gy \neq 0$ and $y^T GM^{-1}Gy \neq 0$. □

The following lemma will be used to provide the upper eigenvalue bound of the preconditioned operator. It is similar to [20, Lemma 6.4] and [30, Lemmas 8 and 11].

Lemma 7.2. There exist positive constants $C_1$ and $C_2$, such that for all $v \in V_0$,

$$\langle M^{-1}B_Cv, B_Cv \rangle \leq (C_1 \alpha + C_2 \Phi(H/h)) \langle \tilde{A}v, v \rangle,$$

where $\Phi(H/h)$ is defined in Lemmas 6.1 and 6.2, respectively.
Proof: Given \( v = (w_I, q, w_\Delta, w_\Pi) \in \tilde{V}_0 \), let \( g_{pr} = B_{II}w_I + B_{I\Delta}w_\Delta + B_{I\Pi}w_\Pi \). From (4.3), (6.5)–(6.8), (6.1), and (6.3), we have

\[
\langle M^{-1}B_Cv, B_Cv \rangle = \alpha h^{-3} \langle g_{pr}, g_{pr} \rangle + \left( B_\Delta \tilde{R}_\Delta v \right)^T M^{-1}_\lambda B_\Delta \tilde{R}_\Delta v
\]

\[
= \alpha h^{-3} \langle g_{pr}, g_{pr} \rangle + (P_Dv, P_Dv)_\tilde{A} \\
\leq \alpha h^{-3} \langle g_{pr}, g_{pr} \rangle + C_2 \Phi(H/h) \langle v, v \rangle_{\tilde{A}},
\]

(7.2)

where we used Lemmas 6.1 and 6.2 for the last inequality. It is sufficient to bound the first term of the right-hand side in the above inequality.

Since \( v \in \tilde{V}_0 \), we have \( B_{II}w_I + B_{I\Delta}w_\Delta + B_{I\Pi}w_\Pi = 0 \), cf. (5.2). Then

\[
\langle g_{pr}, g_{pr} \rangle = \left[ \begin{array}{c}
B_{II}w_I + B_{I\Delta}w_\Delta + B_{I\Pi}w_\Pi \\
B_{II}w_I + B_{I\Delta}w_\Delta + B_{I\Pi}w_\Pi
\end{array} \right]^T \left[ \begin{array}{c}
B_{II}w_I + B_{I\Delta}w_\Delta + B_{I\Pi}w_\Pi \\
B_{II}w_I + B_{I\Delta}w_\Delta + B_{I\Pi}w_\Pi
\end{array} \right] = \langle \tilde{B}w, \tilde{B}w \rangle,
\]

where \( \tilde{B} \) is defined in (5.4) and \( w = (w_I, w_\Delta, w_\Pi) \in \tilde{W} \). From (2.5) and the stability of \( \tilde{B} \), cf. Lemma 5.1, we have

\[
(7.3) \quad h^{-3} \langle g_{pr}, g_{pr} \rangle = h^{-3} \langle \tilde{B}w, \tilde{B}w \rangle \leq C \left( \tilde{B}w, \tilde{B}w \right)_{\tilde{Z}^{-1}} = C \max_{q \in \tilde{Q}} \frac{\left( \tilde{B}w, q \right)^2}{\left( q, q \right)_{\tilde{Z}}} \\
\leq C_1 \max_{q \in \tilde{Q}} \frac{|w_I|^2_{H^1} ||q||^2_{L^2}}{||q||^2_{L^2}} = C_1 |w_I|^2_{H^1} = C_1 \langle v, v \rangle_{\tilde{A}},
\]

where for the last equality, we used (5.3). \( \square \)

The following lemma will be used to provide the lower eigenvalue bound of the preconditioned operator. In [30, Lemmas 9 and 12], the lower eigenvalue bounds for the lumped and Dirichlet preconditioners were analyzed differently. In the analysis of the Dirichlet preconditioner, subdomain discrete Stokes extensions were used. Such extensions require enforcing the same type divergence free subdomain boundary velocity conditions as discussed in [21], even though they are not necessary for implementing the algorithm. The new proof given in the next lemma works for both the lumped and Dirichlet preconditioners. It does not use the subdomain Stokes extensions and those additional subdomain divergence free boundary conditions are no longer needed. For both type of preconditioners, the coarse level velocity space can be chosen as simple as for solving scalar elliptic problems corresponding to each velocity component.

LEMA 7.3. There exists a constant \( C \), such that for any nonzero \( y = (g_{pr}, g_\lambda) \in R_G \), there exists \( v \in \tilde{V}_0 \), which satisfies \( B_Cv = y, \langle v, v \rangle_{\tilde{A}} \neq 0 \), and

\[
\langle \tilde{A}v, v \rangle \leq C \max \left\{ 1, \left( 1 + \frac{1}{h^2} \right) M^{-1}y, y \right\}.
\]

Proof: Given \( y = (g_{pr}, g_\lambda) \in R_G \), take \( u^{(I)}_{\Delta} = B^T_{\Delta,D}g_\lambda, u^{(I)}_{\Pi} = 0 \), and \( p^{(I)} = 0 \). On each subdomain \( \Omega_i \), let \( u^{(I,0)}_{j} \) be zero for the lumped preconditioner, and be obtained for the Dirichlet preconditioner through the solution of (6.2) with given subdomain boundary values \( u^{(j)}_{\Delta} = u^{(j)}_{\Delta} \). Let \( v^{(I,0)} = \left( u^{(I,0)}_{I}, p^{(I,0)}_{I}, u^{(I,0)}_{\Delta}, u^{(I,0)}_{\Pi} \right) \), the corresponding global vectors \( v^{(I)} = \left( u^{(I)}_{I}, p^{(I)}_{I}, u^{(I)}_{\Delta}, u^{(I)}_{\Pi} \right) \), and \( u^{(I)} = \left( u^{(I)}_{I}, u^{(I)}_{\Delta}, u^{(I)}_{\Pi} \right) \).
Then we have
\begin{equation}
B_{Cv}^{(I)} = \begin{bmatrix}
B_{\Gamma I} & 0 & B_{\Gamma \Delta}
0 & B_{\Delta} & 0
\end{bmatrix}
\begin{bmatrix}
\mathbf{u}_I^{(I)} \\
p_I^{(I)} \\
\mathbf{u}_\Delta^{(I)} \\
\mathbf{u}_\Pi^{(I)}
\end{bmatrix} = \begin{bmatrix}
B_{\Gamma I} \mathbf{u}_I^{(I)} + B_{\Gamma \Delta} \mathbf{u}_\Delta^{(I)} + B_{\Gamma \Pi} \mathbf{u}_\Pi^{(I)}
\end{bmatrix}
\end{equation}

where we have used Lemma 3.2. Also
\begin{equation}
|\mathbf{u}^{(I)}|_{H^1}^2 = \begin{bmatrix}
\mathbf{u}_I^{(I)} \\
\mathbf{u}_\Delta^{(I)} \\
\mathbf{u}_\Pi^{(I)}
\end{bmatrix}^T
\begin{bmatrix}
A_{II} & A_{I\Delta} & A_{I\Pi} \\
A_{I\Delta} & A_{\Delta\Delta} & A_{\Delta\Pi} \\
A_{I\Pi} & A_{\Delta\Pi} & A_{\Pi\Pi}
\end{bmatrix}
\begin{bmatrix}
\mathbf{u}_I^{(I)} \\
\mathbf{u}_\Delta^{(I)} \\
\mathbf{u}_\Pi^{(I)}
\end{bmatrix} = \begin{cases}
|\mathbf{u}_\Delta^{(I)}|_{A\Delta\Delta}^2, & \text{for the lumped preconditioner,} \\
|\mathbf{u}_\Delta^{(I)}|_{B\Delta\Delta}^2, & \text{for the Dirichlet preconditioner.}
\end{cases}
\end{equation}

We consider a solution to the following fully assembled system of linear equations of the form (2.3): find \( (\mathbf{u}_I^{(II)}, p_I^{(II)}, \mathbf{u}_\Gamma^{(II)}, p_\Gamma^{(II)}) \) \( \in \mathbf{W}_I \oplus Q_I \oplus \mathbf{W}_\Gamma \oplus Q_\Gamma \), such that
\begin{equation}
\begin{bmatrix}
A_{II} & B_{II}^T & A_{I\Gamma} & B_{I\Gamma}^T \\
B_{II} & 0 & B_{I\Gamma} & 0 \\
A_{I\Gamma} & B_{I\Gamma} & A_{\Gamma\Gamma} & B_{\Gamma\Gamma}^T \\
B_{I\Gamma} & 0 & B_{\Gamma\Gamma} & 0
\end{bmatrix}
\begin{bmatrix}
\mathbf{u}_I^{(II)} \\
p_I^{(II)} \\
\mathbf{u}_\Gamma^{(II)} \\
p_\Gamma^{(II)}
\end{bmatrix} = \begin{bmatrix}
0 \\
-B_{II} \mathbf{u}_I^{(I)} - B_{I\Delta} \mathbf{u}_\Delta^{(I)} - B_{I\Pi} \mathbf{u}_\Pi^{(I)} \\
0 \\
g_{pr} - B_{\Gamma I} \mathbf{u}_I^{(I)} - B_{\Gamma \Delta} \mathbf{u}_\Delta^{(I)} - B_{\Gamma \Pi} \mathbf{u}_\Pi^{(I)}
\end{bmatrix},
\end{equation}

where we know that the particularly chosen right-hand side is essentially
\begin{equation}
\begin{bmatrix}
0 \\
-B_{II} \mathbf{u}_I^{(I)} - B_{I\Delta} \mathbf{u}_\Delta^{(I)} \\
0 \\
g_{pr} - B_{\Gamma I} \mathbf{u}_I^{(I)} - B_{\Gamma \Delta} \mathbf{u}_\Delta^{(I)}
\end{bmatrix}.
\end{equation}

Since \( (g_{pr}, g_\lambda) \in R_G \), we have, cf. (4.5),
\((-B_{I\Delta} \mathbf{u}_\Delta^{(I)})^T 1_{pr} + (g_{pr} - B_{\Gamma \Delta} \mathbf{u}_\Delta^{(I)})^T 1_{pr} = g_{pr}^T 1_{pr} - g_\lambda^T B_{\Delta,D} (B_{I\Delta}^T 1_{pr} + B_{I\Delta}^T 1_{pr}) = 0.\)

Meanwhile,
\((-B_{II} \mathbf{u}_I^{(I)})^T 1_{pr} + (-B_{\Gamma I} \mathbf{u}_I^{(I)})^T 1_{pr} = -\int_\Omega (\nabla \cdot \mathbf{u}_I^{(I)}) 1 = 0.\)

We have that the right-hand side vector (7.7) has zero average, which implies existence of the solution to (7.6).
Denote \( u^{(I)} = (u_I^{(I)}, u_I^{(I)}) \). Then from Lemma 5.2 and (2.5), we have

\[
|u^{(I)}|_{H^1}^2 \leq \frac{1}{\beta^2} \left\| \begin{bmatrix} -B_{II} u_I^{(I)} & B_{II} u_I^{(I)} \\ B_{II} u_I^{(I)} & B_{II} u_I^{(I)} \end{bmatrix} \right\|_{\mathbb{Z}^{-1}}^2 \\
\leq \frac{1}{\beta^2} \left\| \begin{bmatrix} B_{II} u_I^{(I)} + B_{II} u_I^{(I)} + B_{II} u_I^{(I)} \\ B_{II} u_I^{(I)} + B_{II} u_I^{(I)} + B_{II} u_I^{(I)} \end{bmatrix} \right\|_{\mathbb{Z}^{-1}}^2 + \frac{1}{\beta^2} \left\| \begin{bmatrix} 0 \\ g_{pr} \end{bmatrix} \right\|_{\mathbb{Z}^{-1}}^2
\]

(7.8) \[
\leq \frac{1}{\beta^2} |u_I^{(I)}|_{H^1}^2 + \frac{C}{\beta^2 h^3} \langle g_{pr}, g_{pr} \rangle,
\]

where the bound on the first term is obtained in the same way as in (7.3).

Split the continuous subdomain interface velocity \( u_I^{(I)} \) into the dual part \( u_{\Delta}^{(I)} \) and the primal part \( u_{II}^{(I)} \), and denote \( v^{(I)} = (u_I^{(I)}, p_I^{(I)}, u_{\Delta}^{(I)}, u_{II}^{(I)}) \). Let \( v = v^{(I)} + v^{(II)} \). Then we have from (7.6) that \( v \in V_0 \), and

\[
B_C v^{(II)} = \begin{bmatrix} B_{II} & 0 & B_{II} & 0 \\ 0 & 0 & B_{\Delta} & 0 \\ B_{II} & 0 & B_{II} & 0 \\ 0 & 0 & B_{\Delta} & 0 \end{bmatrix} \begin{bmatrix} u_I^{(I)} \\ p_I^{(I)} \\ u_{\Delta}^{(I)} \\ u_{II}^{(I)} \end{bmatrix} = \begin{bmatrix} g_{pr} - B_{II} u_I^{(I)} - B_{II} u_{II}^{(I)} \\ 0 \end{bmatrix}.
\]

Together with (7.4), we have \( B_C v = y \). From (5.3) and (7.8), we have

\[
|v|_{A}^2 = |u^{(I)} + u^{(II)}|_{H^1}^2 \leq |u^{(I)}|_{H^1}^2 + |u^{(II)}|_{H^1}^2 = \left( 1 + \frac{1}{\beta^2} \right) |u^{(I)}|_{H^1}^2 + \frac{C}{\beta^2 h^3} \langle g_{pr}, g_{pr} \rangle
\]

\[
= \begin{cases} 
(1 + \frac{1}{\beta^2}) |u_{\Delta}^{(II)}|_{A_{\Delta \Delta}}^2 + \frac{C}{\beta^2 h^3} \langle g_{pr}, g_{pr} \rangle, & \text{for the lumped preconditioner,} \\
(1 + \frac{1}{\beta^2}) |u_{\Delta}^{(II)}|_{H_{\Delta}}^2 + \frac{C}{\beta^2 h^3} \langle g_{pr}, g_{pr} \rangle, & \text{for the Dirichlet preconditioner,}
\end{cases}
\]

where we used (7.5) in the last equality.

On the other hand, we have from (6.5)–(6.8)

\[
(M^{-1} y, y) = \frac{\alpha}{h^3} \langle g_{pr}, g_{pr} \rangle + g_{pr}^T M_{\lambda}^{-1} g_{\lambda}
\]

\[
= \begin{cases} 
\frac{\alpha}{h^3} \langle g_{pr}, g_{pr} \rangle + g_{pr}^T B_{\Delta, D} A_{\Delta \Delta} B_{\Delta, D}^T g_{\lambda}, & \text{for the lumped preconditioner,} \\
\frac{\alpha}{h^3} \langle g_{pr}, g_{pr} \rangle + g_{pr}^T B_{\Delta, D} H_{\Delta} B_{\Delta, D}^T g_{\lambda}, & \text{for the Dirichlet preconditioner},
\end{cases}
\]

\[
= \begin{cases} 
\frac{\alpha}{h^3} \langle g_{pr}, g_{pr} \rangle + |u_{\Delta}^{(II)}|_{A_{\Delta \Delta}}^2, & \text{for the lumped preconditioner,} \\
\frac{\alpha}{h^3} \langle g_{pr}, g_{pr} \rangle + |u_{\Delta}^{(II)}|_{H_{\Delta}}^2, & \text{for the Dirichlet preconditioner.}
\end{cases}
\]
It is not difficult to see that \( \langle v, v \rangle_A \neq 0 \). Otherwise, all the velocity components of \( v \) would be zero, cf. \((5.3)\), and then \( B_C v \) would be zero, which conflicts with \( B_C v = y \) and \( y \) being nonzero.

The proofs of the following two lemmas can be found at \([20, \text{Lemmas} 6.6 \text{and} 6.3]\).

**Lemma 7.4.** For any \( v = (w_I, p_I, w_\Delta, w_\Pi) \in \bar{V}_0, B_C v \in R_G \).

**Lemma 7.5.** For any \( x \in R_{M^{-1}G} \),

\[
\langle Mx, x \rangle = \max_{y \in R_G, y \neq 0} \frac{\langle y, x \rangle^2}{\langle M^{-1}y, y \rangle}.
\]

The condition number bound of the preconditioned operator \( M^{-1}G \) is given in the following theorem.

**Theorem 7.6.** There exist positive constants \( c, C_1 \) and \( C_2 \), such that for all \( x \in R_{M^{-1}G} \),

\[
\min \left\{ 1, \alpha \right\} \frac{c \beta^2}{(1 + \beta^2)} \langle Mx, x \rangle \leq \langle Gx, x \rangle \leq (C_1 \alpha + C_2 \Phi(H/h)) \langle Mx, x \rangle.
\]

**Proof:** We only need to prove the above inequalities for any nonzero \( x \in R_{M^{-1}G} \).

We know from Lemma 7.1 that

\[
0 \neq \langle Gx, x \rangle = x^T B_C \tilde{A}^{-1} B_C^T x = x^T B_C \tilde{A}^{-1} \tilde{A} \tilde{A}^{-1} B_C^T x = \left( \tilde{A}^{-1} B_C^T x, \tilde{A}^{-1} B_C^T x \right)_A.
\]

Therefore \( \tilde{A}^{-1} B_C^T x \neq 0 \). Also note that \( \tilde{A}^{-1} B_C^T x \in \bar{V}_0 \) and \( \langle \cdot, \cdot \rangle_A \) defines a semi-inner product on \( \bar{V}_0 \), cf \((5.3)\), and then we have

\[
\langle Gx, x \rangle = \max_{v \in \bar{V}_0, \langle v, v \rangle_A \neq 0} \frac{\langle v, \tilde{A}^{-1} B_C^T x \rangle_A^2}{\langle v, v \rangle_A} = \max_{v \in \bar{V}_0, \langle v, v \rangle_A \neq 0} \frac{\langle B_C v, x \rangle^2}{\langle \tilde{A} v, v \rangle}.
\]

**Lower bound:** From Lemma 7.3, we know that for any nonzero \( y \in R_G \), there exists \( w \in \bar{V}_0 \), such that \( B_C w = y \), \( \langle w, w \rangle_A \neq 0 \), \( \langle \tilde{A} w, w \rangle \leq \max \left\{ 1, \frac{1}{\alpha} \right\} \frac{C(1 + \beta^2)}{\beta^2} \langle M^{-1}y, y \rangle \).

Then from \((7.9)\), we have

\[
\langle Gx, x \rangle \geq \frac{\langle B_C w, x \rangle^2}{\langle \tilde{A} w, w \rangle} \geq c \max \left\{ 1, \frac{1}{\alpha} \right\} \frac{\beta^2}{(1 + \beta^2)} \langle M^{-1}y, y \rangle \langle y, x \rangle^2.
\]

Since \( y \) is arbitrary, using Lemma 7.5, we have

\[
\langle Gx, x \rangle \geq c \max \left\{ 1, \frac{1}{\alpha} \right\} \frac{\beta^2}{(1 + \beta^2)} \frac{\langle y, x \rangle^2}{\langle M^{-1}y, y \rangle} = \min \left\{ 1, \alpha \right\} \frac{c \beta^2}{(1 + \beta^2)} \langle Mx, x \rangle.
\]

**Upper bound:** From \((7.9)\) and the fact that \( \langle Gx, x \rangle \neq 0 \), we have

\[
\langle Gx, x \rangle = \max_{v \in \bar{V}_0, \langle v, v \rangle_A \neq 0} \frac{\langle B_C v, x \rangle^2}{\langle \tilde{A} v, v \rangle} = \max_{v \in \bar{V}_0, \langle v, v \rangle_A \neq 0, B_C v \neq 0} \frac{\langle B_C v, x \rangle^2}{\langle \tilde{A} v, v \rangle}.
\]
where the maximum only needs to be considered among $v$ also satisfying $B_C v \neq 0$. Then using Lemmas 7.2, 7.4, and 7.5, we have

$$
\langle Gx, x \rangle \leq (C_1 \alpha + C_2 \Phi(H, h)) \max_{v \in \tilde{V}_0, y \neq 0} \frac{(B_C v, x)^2}{(M^{-1}B_C v, B_C v)} \leq (C_1 \alpha + C_2 \Phi(H, h)) \max_{y \in \mathbb{R}^G, y \neq 0} \frac{(M^{-1}y, y)}{(M^{-1}y, y)} = (C_1 \alpha + C_2 \Phi(H, h)) \langle Mx, x \rangle.
$$

\[ \square \]

Remark 7.7. We can see from Theorem 7.6 that, for $\alpha \geq 1$, the condition number bound of $M^{-1}G$ is proportional to $\alpha + C\Phi(H, h)$, and we should take smaller $\alpha$ to achieve faster convergence. When $\alpha \leq 1$, the condition number bound is proportional to $1 + C\Phi(H, h)\alpha$ and we should take larger $\alpha$. This explains why the value of $\alpha$ in (6.6) and (6.7) is typically taken as 1. We introduce $\alpha$ in the preconditioner just for the convenience to demonstrate the convergence rates of the proposed algorithm in the following section.

8. Numerical experiments. We illustrate the convergence rate of the proposed algorithm by solving the incompressible Stokes problem (2.1) in both two and three dimensions, on $\Omega = [0, 1]^2$ and $\Omega = [0, 1]^3$, respectively. The right-hand side $f$ is chosen such that the exact solution is

$$
u = \begin{bmatrix} \sin^3(\pi x) \sin^2(\pi y) \cos(\pi y) \\ -\sin^2(\pi x) \sin^3(\pi y) \cos(\pi x) \end{bmatrix}, \quad p = x^2 - y^2,$$

for two dimensions, and for three dimensions

$$
u = \begin{bmatrix} \sin^2(\pi x) (\sin(2\pi y) \sin(\pi z) - \sin(\pi y) \sin(2\pi z)) \\ \sin^2(\pi y) (\sin(2\pi z) \sin(\pi x) - \sin(\pi z) \sin(2\pi x)) \\ \sin^2(\pi z) (\sin(2\pi x) \sin(\pi y) - \sin(\pi x) \sin(2\pi y)) \end{bmatrix}, \quad p = xyz - \frac{1}{8},$$

with zero Dirichlet boundary condition. We also test a mixed boundary condition with the same right-hand side $f$.

Three types of mixed finite elements are used in the experiments. They are two- and three-dimensional $Q_2-Q_1$ Taylor-Hood mixed finite elements with continuous pressures [2, 26], and a three-dimensional mix finite element with discontinuous pressures, which is used in [18]. For the continuous pressure elements, the velocity space contains piecewise bi-quadratic functions and the pressure space contains piecewise bilinear functions in two dimensions; piecewise tri-quadratic functions for the velocity and piecewise trilinear functions for the pressure in three dimensions. The discontinuous pressure element is illustrated in Figure 8.1, where the velocity is spanned by 1, $x, y, z, zx, zy$ on each prism, and the pressure is a constant on each eight prisms.

![Fig. 8.1. A three-dimensional mixed finite element with discontinuous pressures.](image-url)
following tables list the minimum and maximum eigenvalues (\(\lambda_{\text{min}}\) and \(\lambda_{\text{max}}\)) of the iteration matrix \(M^{-1}G\) and the iteration count (\(\text{iter}\)), for using both the lumped and Dirichlet preconditioners and in different cases. Here \(\lambda_{\text{min}}\) and \(\lambda_{\text{max}}\) are estimated by using the tridiagonal Lanczos matrix generated in the iteration.

We note that the lower eigenvalue bound of the preconditioned operator is proportional to the inf-sup constant \(\beta^2\), defined in (2.4), as shown in Theorem 7.6. Estimate of \(\beta^2\) can be obtained as in [9, Lemma 2.3] and their values for these three mixed elements are

\[
(8.1) \quad \beta^2_{2D} = 0.0719; \quad \beta^2_{3D,\text{continuous}} = 0.0189; \quad \beta^2_{3D,\text{discontinuous}} = 0.1053.
\]

From these values, we can expect and have indeed observed in the tables that \(\lambda_{\text{min}}\), for three-dimensional finite element with continuous pressure (Table 8.3), is quite small compared with those for the two-dimensional case (Table 8.1) and for the three-dimensional discontinuous pressure case (Table 8.2).

Table 8.1 shows the performance for solving the two-dimensional problem. The coarse level velocity space is spanned by the subdomain corner nodal basis functions corresponding to each velocity component. Table 8.2 shows the performance for solving the three-dimensional problem using the discontinuous pressures mixed finite elements illustrated in Figure 8.1. The coarse level velocity space is spanned by the subdomain vertex nodal basis functions and subdomain edge-cutoff functions corresponding to each velocity component. We take \(\alpha = 1\) in both the lumped and Dirichlet preconditioners (6.6) and (6.7). These results are consistent with our theory. The error of the finite element solution is also listed in Tables 8.1 and 8.2.

### Table 8.1

<table>
<thead>
<tr>
<th>(H^k)</th>
<th>#sub</th>
<th>(|\mathbf{u} - \mathbf{u}_h|_2)</th>
<th>(|p - p_h|_2)</th>
<th>(\lambda_{\text{min}})</th>
<th>(\lambda_{\text{max}})</th>
<th>(\text{iter})</th>
<th>(\lambda_{\text{min}})</th>
<th>(\lambda_{\text{max}})</th>
<th>(\text{iter})</th>
</tr>
</thead>
<tbody>
<tr>
<td>4</td>
<td>(9.50e-7)</td>
<td>1.15e-5</td>
<td>0.3066</td>
<td>32.28</td>
<td>31</td>
<td>0.2983</td>
<td>4.40</td>
<td>18</td>
<td></td>
</tr>
<tr>
<td>8</td>
<td>(5.98e-8)</td>
<td>7.00e-7</td>
<td>0.3068</td>
<td>38.42</td>
<td>51</td>
<td>0.2859</td>
<td>5.03</td>
<td>24</td>
<td></td>
</tr>
<tr>
<td>16</td>
<td>(3.74e-9)</td>
<td>4.33e-8</td>
<td>0.3067</td>
<td>85.32</td>
<td>62</td>
<td>0.2966</td>
<td>6.04</td>
<td>25</td>
<td></td>
</tr>
<tr>
<td>32</td>
<td>(2.36e-10)</td>
<td>3.04e-9</td>
<td>0.3070</td>
<td>192.32</td>
<td>83</td>
<td>0.3070</td>
<td>7.19</td>
<td>27</td>
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</tbody>
</table>

### Table 8.2

<table>
<thead>
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<th>#sub</th>
<th>(|\mathbf{u} - \mathbf{u}_h|_2)</th>
<th>(|p - p_h|_2)</th>
<th>(\lambda_{\text{min}})</th>
<th>(\lambda_{\text{max}})</th>
<th>(\text{iter})</th>
<th>(\lambda_{\text{min}})</th>
<th>(\lambda_{\text{max}})</th>
<th>(\text{iter})</th>
</tr>
</thead>
<tbody>
<tr>
<td>4</td>
<td>(9.50e-7)</td>
<td>1.15e-5</td>
<td>0.3024</td>
<td>15.91</td>
<td>34</td>
<td>0.2706</td>
<td>4.15</td>
<td>21</td>
<td></td>
</tr>
<tr>
<td>8</td>
<td>(5.98e-8)</td>
<td>7.00e-7</td>
<td>0.3067</td>
<td>37.25</td>
<td>46</td>
<td>0.2859</td>
<td>5.03</td>
<td>24</td>
<td></td>
</tr>
<tr>
<td>16</td>
<td>(3.74e-9)</td>
<td>4.33e-8</td>
<td>0.3068</td>
<td>85.32</td>
<td>62</td>
<td>0.2966</td>
<td>6.04</td>
<td>25</td>
<td></td>
</tr>
<tr>
<td>32</td>
<td>(2.36e-10)</td>
<td>3.04e-9</td>
<td>0.3070</td>
<td>192.32</td>
<td>83</td>
<td>0.3070</td>
<td>7.19</td>
<td>27</td>
<td></td>
</tr>
</tbody>
</table>

Tables 8.3 and 8.4 are for solving the three-dimensional problem using the mixed finite element with continuous pressures. The coarse velocity space is spanned by the subdomain vertex nodal basis functions and subdomain edge-cutoff functions corresponding to each velocity component. In Table 8.3, \(\alpha = 1\); in Table 8.4, \(\alpha = 1/2\). As expected, \(\lambda_{\text{min}}\) is smaller than those in Tables 8.1 and 8.2, cf. (8.1). However, the error of the finite element solution using continuous pressures shown in Table 8.3 is much better than that shown in Table 8.2 for using discontinuous pressures.
Table 8.2
Solving 3D problem using discontinuous pressure discretization, $\alpha = 1$ in (6.6) and (6.7).

<table>
<thead>
<tr>
<th>$\frac{H}{h}$</th>
<th>#sub</th>
<th>$|u - u_h|_2$</th>
<th>$|p - p_h|_2$</th>
<th>$\lambda_{\min}$</th>
<th>$\lambda_{\max}$</th>
<th>iter</th>
<th>$\lambda_{\min}$</th>
<th>$\lambda_{\max}$</th>
<th>iter</th>
</tr>
</thead>
<tbody>
<tr>
<td>$3^3$</td>
<td>3</td>
<td>1.24e-2</td>
<td>6.15e-2</td>
<td>0.2509</td>
<td>3.72</td>
<td>21</td>
<td>0.2510</td>
<td>3.15</td>
<td>18</td>
</tr>
<tr>
<td>4</td>
<td>4</td>
<td>6.97e-3</td>
<td>3.45e-2</td>
<td>0.2549</td>
<td>3.96</td>
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<td>1.53e-2</td>
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<td>23</td>
<td>0.2556</td>
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<td>20</td>
</tr>
<tr>
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<td>8</td>
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<td>8.62e-3</td>
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<td>4.08</td>
<td>23</td>
<td>0.2082</td>
<td>3.22</td>
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</tr>
</tbody>
</table>

Table 8.3
Solving 3D problem using continuous pressure discretization, $\alpha = 1$ in (6.6) and (6.7).

<table>
<thead>
<tr>
<th>$\frac{H}{h}$</th>
<th>#sub</th>
<th>$|u - u_h|_2$</th>
<th>$|p - p_h|_2$</th>
<th>$\lambda_{\min}$</th>
<th>$\lambda_{\max}$</th>
<th>iter</th>
<th>$\lambda_{\min}$</th>
<th>$\lambda_{\max}$</th>
<th>iter</th>
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<tbody>
<tr>
<td>$3^3$</td>
<td>3</td>
<td>2.36e-2</td>
<td>1.10e-1</td>
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<td>3.19</td>
<td>20</td>
<td>0.2467</td>
<td>3.02</td>
<td>18</td>
</tr>
<tr>
<td>4</td>
<td>4</td>
<td>1.24e-2</td>
<td>6.15e-2</td>
<td>0.2509</td>
<td>3.72</td>
<td>21</td>
<td>0.2510</td>
<td>3.15</td>
<td>19</td>
</tr>
<tr>
<td>$6^3$</td>
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<td>5.12e-3</td>
<td>2.72e-2</td>
<td>0.2543</td>
<td>5.77</td>
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<td>0.2590</td>
<td>3.28</td>
<td>19</td>
</tr>
<tr>
<td>$8^3$</td>
<td>8</td>
<td>3.09e-3</td>
<td>1.53e-2</td>
<td>0.2548</td>
<td>8.20</td>
<td>28</td>
<td>0.2794</td>
<td>3.34</td>
<td>18</td>
</tr>
</tbody>
</table>

In Table 8.3, the minimum eigenvalue is independent of the mesh size for both preconditioners. The maximum eigenvalue is independent of the number of subdomains for fixed $H/h$; we have to admit that due to the limit of computing resource, we are not able to experiment with smaller mesh size to show the asymptotic behavior of convergence rate more accurately. For fixed number of subdomains, $\lambda_{\max}$ depends on $H/h$ and its least squares fits by $C_1\alpha + C_2(H/h)(1 + \log(H/h))$ for the lumped preconditioner and by $C_1\alpha + C_2(1 + \log(H/h))^2$ for the Dirichlet preconditioner, as guided by Theorem 7.6, are shown in Figure 8.2. Moreover, the convergence rate of the algorithm using the Dirichlet preconditioner is only slightly better than using the lumped preconditioner. The reason is that the upper eigenvalue bound in Theorem 7.6 depends on both $\alpha$ and $\Phi(H/h)$, and in this case $\alpha = 1$ dominates when $H/h$ is small. Therefore, even though using the Dirichlet preconditioner can reduce $\Phi(H/h)$ compared with using the lumped preconditioner, this improvement cannot show up in
Table 8.4
Solving 3D problem using continuous pressure discretization, $\alpha = 1/2$ in (6.6) and (6.7).

<table>
<thead>
<tr>
<th>$H/h$</th>
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<th></th>
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<tr>
<td></td>
<td></td>
<td>$\lambda_{\text{min}}$</td>
<td>$\lambda_{\text{max}}$</td>
<td>iteration</td>
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<tr>
<td></td>
<td></td>
<td>0.0395</td>
<td>7.20</td>
<td>59</td>
<td></td>
</tr>
<tr>
<td>3</td>
<td>$3^3$</td>
<td>0.0394</td>
<td>8.15</td>
<td>66</td>
<td></td>
</tr>
<tr>
<td>4</td>
<td>$4^3$</td>
<td>0.0393</td>
<td>8.85</td>
<td>70</td>
<td></td>
</tr>
<tr>
<td>6</td>
<td>$6^3$</td>
<td>0.0393</td>
<td>9.09</td>
<td>72</td>
<td></td>
</tr>
<tr>
<td>8</td>
<td>$8^3$</td>
<td>0.0392</td>
<td>11.70</td>
<td>73</td>
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</tbody>
</table>

<table>
<thead>
<tr>
<th>$# \text{sub}$</th>
<th>$H/h$</th>
<th>$\lambda_{\text{min}}$</th>
<th>$\lambda_{\text{max}}$</th>
<th>iteration</th>
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<tbody>
<tr>
<td>3</td>
<td>3</td>
<td>0.0387</td>
<td>5.15</td>
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<tr>
<td>$3^3$</td>
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<td>6</td>
<td>0.0397</td>
<td>11.70</td>
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<tr>
<td></td>
<td>8</td>
<td>0.0397</td>
<td>16.52</td>
<td>73</td>
<td></td>
</tr>
</tbody>
</table>

Table 8.3. What shows in Table 8.3 for $\lambda_{\text{max}}$ is essentially its dependence on $\alpha$. Only for larger $H/h$, e.g., for $H/h = 6$ and $H/h = 8$ in Table 8.3, this improvement on $\lambda_{\text{max}}$ by using the Dirichlet preconditioner becomes visible.

To experiment the case when $\alpha$ is less dominant in the upper eigenvalue bound, we take $\alpha = 1/2$ in Table 8.4. Consistent with Theorem 7.6, the lower eigenvalue bounds in Table 8.4 become half of those in Table 8.3 and they are also independent of the mesh size. The upper eigenvalue bounds show the improvement by using the Dirichlet preconditioner. The least squares fits of $\lambda_{\text{max}}$ with respect to $H/h$ are also shown in Figure 8.2.

We also test using an enriched coarse space in the algorithm with continuous pressures on the three-dimensional problem. Besides the subdomain vertex nodal basis functions and subdomain edge-cutoff functions corresponding to each velocity component, one face-cutoff function representing the normal velocity component on each subdomain face is also included in the coarse space. The performance of the algorithm is listed in Table 8.5. Compared with Table 8.4, the iteration counts are reduced for both preconditioners, and the condition number bound is improved a little bit in the lumped preconditioner case, while the improvement is too small to be seen in the Dirichlet preconditioner case, cf. Remark 6.3.

At the end, we test the algorithm solving the three-dimensional problem with a mixed boundary condition where the normal derivative of each velocity component equals zero at the bottom face of the cube and the velocity equals zero at the other faces. We can see from Table 8.6 that the performance is only a bit worse than that for solving the Dirichlet boundary condition problem shown in Table 8.4.

REFERENCES

Fig. 8.2. Least squares fits of $\lambda_{\text{max}}$ in Table 8.3 (top row) and in Table 8.4 (bottom row) with respect to $H/h$ by $C_1 \alpha + C_2 (H/h)(1 + \log (H/h))$ for the lumped preconditioner (left) and by $C_1 \alpha + C_2 (1 + \log (H/h))^2$ for the Dirichlet preconditioner (right).

<table>
<thead>
<tr>
<th>$H/h$</th>
<th>#sub</th>
<th>lumped</th>
<th>Dirichlet</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td></td>
<td>$\lambda_{\text{min}}$</td>
<td>$\lambda_{\text{max}}$</td>
</tr>
<tr>
<td>4</td>
<td>$3^3$</td>
<td>0.0395</td>
<td>5.07</td>
</tr>
<tr>
<td></td>
<td>$4^3$</td>
<td>0.0394</td>
<td>5.51</td>
</tr>
<tr>
<td></td>
<td>$6^3$</td>
<td>0.0393</td>
<td>6.18</td>
</tr>
<tr>
<td></td>
<td>$8^3$</td>
<td>0.0393</td>
<td>6.43</td>
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<td></td>
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<td>$\lambda_{\text{max}}$</td>
</tr>
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<td>3</td>
<td>0.0388</td>
<td>4.47</td>
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<td></td>
<td>$4^3$</td>
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<tr>
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<td>$6^3$</td>
<td>0.0396</td>
<td>7.10</td>
</tr>
<tr>
<td></td>
<td>$8^3$</td>
<td>0.0397</td>
<td>9.67</td>
</tr>
</tbody>
</table>


Table 8.6

Solving 3D problem with mixed boundary conditions, $\alpha = 1/2$ in (6.6) and (6.7).

<table>
<thead>
<tr>
<th>$H/h$</th>
<th>#sub</th>
<th>lumped $\lambda_{\text{min}}$</th>
<th>lumped $\lambda_{\text{max}}$</th>
<th>lumped iter</th>
<th>Dirichlet $\lambda_{\text{min}}$</th>
<th>Dirichlet $\lambda_{\text{max}}$</th>
<th>Dirichlet iter</th>
</tr>
</thead>
<tbody>
<tr>
<td>3</td>
<td>3</td>
<td>0.0395</td>
<td>7.32</td>
<td>71</td>
<td>0.0394</td>
<td>4.89</td>
<td>62</td>
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<tr>
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<td>0.0394</td>
<td>8.22</td>
<td>80</td>
<td>0.0389</td>
<td>5.01</td>
<td>66</td>
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<tr>
<td>6</td>
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<td>0.0393</td>
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<td>86</td>
<td>0.0297</td>
<td>5.10</td>
<td>69</td>
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<tr>
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<td>$8^3$</td>
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<td>9.10</td>
<td>89</td>
<td>0.0237</td>
<td>5.13</td>
<td>71</td>
</tr>
</tbody>
</table>

#sub $H/h$ $\lambda_{\text{min}}$ $\lambda_{\text{max}}$ iter $\lambda_{\text{min}}$ $\lambda_{\text{max}}$ iter

| 3    | 3    | 0.0391                        | 5.22                         | 67          | 0.0396                        | 4.35                         | 64           |
| 3    | 4    | 0.0395                        | 7.32                         | 71          | 0.0394                        | 4.89                         | 62           |
| 6    | 4    | 0.0397                        | 11.91                        | 81          | 0.0397                        | 5.11                         | 62           |
| 8    | 8    | 0.0397                        | 16.81                        | 90          | 0.0397                        | 5.17                         | 61           |


