# Parallel iterative procedures for a computational electromagnetic modeling based on a nonconforming mixed finite element method 

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#### Abstract

We present nonoverlapping domain decomposition methods for the approximation of both electromagnetic fields in a three-dimensional bounded domain satisfying absorbing boundary conditions. A Seidel-type domain decomposition iterative method is introduced based on a hybridization of a nonconforming mixed finite element method. Convergence results for the numerical procedure are proved by introducing a suitable pseudoenergy. The spectral radius of the iterative procedure is estimated and a method for choosing an optimal matching parameter is given. A red-black Seidel-type method which is readily parallelizable is also introduced and analyzed. Numerical experiments confirm that the presented algorithms are faster than the conventional Jacobi-type ones.


keyword: Nonoverlapping domain decomposition, Maxwell's equations, Transmission condition, Parallel iterative algorithm

## 1 Introduction

Computational electromagnetic modeling has many applications in engineering and industry and thus there have been increasing attention from both scientists and engineers: see, for instance, Ben Belgacem, Buffa, and Maday (2001); Boffi, Demkowicz, and Costabel (2003); Bouillault, Buffa, Maday, and Rapetti (2003); Gopalakrishnan, Pasciak, and Demkowicz (2004); Hiptmair and Schwab (2002); Hiptmair (2002); Hu and Zou (2003); Jose, Kanapady, and Tamma (2004); Monk (2003); Reitich and Tamma (2004); Toselli and Vasseur (2005); Volakis, Sertel, Jorgensen, and Kindt (2004), and the references therein

In this work we are particularly interested in the numer-

[^0]ical approximation of scattered electromagnetic fields $E=E(x, \omega)$ and $H=H(x, \omega)$ which satisfy the timeharmonic Maxwell's equations at a given angular frequency $\omega$. Absorbing boundary conditions are used to truncate the infinite, or at least huge, space domains into domains of reasonably smaller size (see Bendali and Halpern (1988); Berenger (1994, 1996); Feng (1999); Hanouzet and Sesqès (1990); Joly and Mercier (1989); Mur (1981); Sesquès (1990); Sheen (1997) and the references therein). However, the numerical treatment of three-dimensional time-harmonic Maxwell's equations still requires an efficient iterative procedure rather than a direct Gaussian elimination one for solving resulting matrix problems. Recently, several domain decomposition methods have been proposed to solve this problem; for instance, see Alonso and Valli (1997, 1999); Ben Belgacem, Buffa, and Maday (2001); Bouillault, Buffa, Maday, and Rapetti (2003); Collino, Delbue, Joly, and Piacentini (1997); Després (1991); Després, Joly, and Roberts (1992); Jr., Santos, and Sheen (2001); Hu and Zou (2003); Santos (1998); Santos and Sheen (1998); Toselli and Vasseur (2005).
In particular, following the idea in Lions (1988, 1990), Després (1991) and Després, Joly, and Roberts (1992) proposed an iterative method using nonoverlapping domain decompositions and a pseudo-energy with which convergence results were shown. The convergence results obtained by Després were weak in the sense no actual convergence rates are given. Later, Douglas Jr, Paes Leme, Roberts, and Wang (1993) applied this idea to second-order elliptic problems, and obtained, using mixed finite elements, an estimate for the spectral radius of the iterative operator which guarantees the actual rate of convergence of the scheme. By utilizing similar structures to mixed finite element spaces, such estimates for spectral radii have been obtained, based on hybridization of nonconforming finite elements, for elliptic, Helmholtz, and viscoelastic problems )Douglas Jr., Santos, Sheen, and Ye (1999); Ha, Santos, and Sheen (2002);

Jr., Santos, and Sheen (2001).
Recently a nonconforming mixed finite element method has been proposed in Douglas Jr., Santos, and Sheen (2000) to compute $E$ and $H$ simultaneously. Error estimates using the were also given. In that paper, Jacobitype domain decomposition iterative procedure based on hybridization has been analyzed with an actual rate convergence. The object of our current paper is to study Seidel-type domain decomposition iterative procedures for calculating both $E$ and $H$ satisfying the full Maxwell's equations, based on nonconforming mixed finite elements proposed in Douglas Jr., Santos, and Sheen (2000). In the domain decomposition methods the convergence speed depends on a choice of matching parameter $\beta$ between the tangential electric and magnetic fields that appear in the interface condition on the common boundaries of domains $\Omega_{j}$ and $\Omega_{k}$ :
$\left(\mathbf{v}_{j} \times H_{j}+\beta_{j k} \pi_{\tau} E_{j}\right)=-\mathbf{v}_{k} \times H_{k}+\beta_{j k} \pi_{\tau} E_{k} \quad$ on $\quad \Gamma_{j k} \forall k$.
In this paper we suggest a method of choice of matching parameter $\beta$. Since Seidel-type iterative method is not parallelizable directly, we employ red-black type iterative algorithm in numerical simulations.
Our approach is similar to the conforming version of domain decomposition scheme introduced in Collino and Joly (2000); Collino, Delbue, Joly, and Piacentini (1997), in the sense that the pseudo-energies calculated in our work are to be calculated on both sides of interfaces while those in Collino and Joly (2000); Collino, Delbue, Joly, and Piacentini (1997) are to be calculated on single sides of interfaces. We also give a simple proof of convergence of the iterative procedure which depends only on the resulting weak formulation of problems, and we will show that the speed of Seidel-type procedures is about twice of that of Jacobi-type ones. However, in Collino, Delbue, Joly, and Piacentini (1997) only one field, say $E$, of the two fields $E$ and $H$ is calculated in the second-order formulation which requires the calculation of $\mathbf{v} \times(\boldsymbol{\nabla} \times E)$ at interfaces to update; indeed, in Collino, Delbue, Joly, and Piacentini (1997) the problem is reformulated in the mixed variational form to compute $E$ and $\boldsymbol{v} \times(\boldsymbol{\nabla} \times E)$ in subdomains and their boundaries. We also suggest a method to find iteratively an optimal choice of the matching parameter between the tangential components of $E$ and $H$ on each subdomain interface. The use of nonlocal boundary conditions, their higher-order approximations, or perfectly matched layers and relaxation
as in Berenger (1994, 1996); Collino, Delbue, Joly, and Piacentini (1997) will be a subject of interest to be developed.
In the formulation of finite element methods for Maxwell's equations with absorbing boundary conditions, and domain decomposition methods for Maxwell's equations, the tangential traces of $\mathbf{H}(\operatorname{curl} ; \Omega)$ are involved in many places. Based on proper understanding the tangential traces of $\mathbf{H}(\mathbf{c u r l} ; \Omega)$, the integration by parts formula for functions in $\mathbf{H}(\mathbf{c u r l} ; \Omega)$ in the sense of Sobolev spaces has been studied intensively in recent literatures Buffa, Costabel, and Sheen (2002) (see also Buffa and P. Ciarlet Jr. (2001); Buffa, Costabel, and Sheen (2002); Buffa (2001).) We will briefly survey these results in Section 2.

The organization of the paper is as follows. In Section 2, Maxwell's equations are described with absorbing boundary conditions, and then a weak formulation is given. In the next section the nonconforming mixed finite element method and its hybridization are given. Then in Section 4 a Seidel-type iterative scheme is proposed, and a suitable pseudo-energy is defined with which convergence and an estimate for the spectral radius of the iterative procedure are shown. A suggestion is made for an optimal choice of matching parameter between the tangential components of $E$ and $H$ for the interface condition is given. Also a red-black Seidel type procedure is introduced which is readily parallelizable. Finally in Section 5 some results from numerical experiments are presented to compare the analyses given in the previous sections.

## 2 The time-harmonic Maxwell's equations

### 2.1 The model problem

Let $E$ and $H$ denote the electric and magnetic fields for a given angular frequency $\omega$. The time-harmonic Maxwell's equations are given by
$(i \omega \varepsilon+\sigma) E-\nabla \times H=F$,
$i \omega \mu H+\nabla \times E=G$,
where $\varepsilon, \mu$, and $\sigma$ denote the electric permittivity, magnetic permeability, and conductivity, respectively, which satisfy the following bound, for any $\xi \in \mathbf{C}^{3}$,
$0<\varepsilon_{*}|\xi|^{2} \leq\left|\sum_{j, k} \varepsilon_{j k} \xi_{j} \overline{\xi_{k}}\right| \leq \varepsilon^{*}|\xi|^{2}$,
$0 \leq \sigma_{*}|\xi|^{2} \leq\left|\sum_{j, k} \sigma_{j k} \xi_{j} \overline{\xi_{k}}\right| \leq \sigma^{*}|\xi|^{2}$,
$0<\mu_{*}|\xi|^{2} \leq\left|\sum_{j, k} \mu_{j k} \xi_{j} \bar{\xi}_{k}\right| \leq \mu^{*}|\xi|^{2}$.
The free space equations are reduced to a truncated compact domain, so that a practical computational procedure can be defined, with suitable absorbing boundary conditions imposed on the truncated boundary. Let $\Omega$ be a unit cube and $\Gamma=\partial \Omega$ be its boundary, and then we will consider the following boundary value problem

$$
\begin{align*}
& (i \omega \varepsilon+\sigma) E-\nabla \times H=F \quad \text { in } \quad \Omega  \tag{3a}\\
& i \omega \mu H+\nabla \times E=G \quad \text { in } \quad \Omega  \tag{3b}\\
& \alpha \pi_{\tau} E+\mathbf{v} \times H=0 \quad \text { on } \quad \Gamma \tag{3c}
\end{align*}
$$

where $\boldsymbol{v}$ denote the unit outer normal to $\Gamma, \pi_{\tau}(\varphi)=$ $\varphi-\boldsymbol{v}(\boldsymbol{v} \cdot \varphi)=-\mathbf{v} \times(\boldsymbol{v} \times \varphi)$ is the projection of the trace of $\varphi$ on $\Gamma$. Assume that, for $\alpha=\alpha_{R}-i \alpha_{I}$, there exist $\alpha_{R_{*}}, \alpha_{R^{*}}, \alpha_{I_{*}}$, and $\alpha_{I^{*}}$, with $0<\alpha_{R_{*}} \leq \alpha_{R^{*}}$ and $0 \leq \alpha_{I_{*}} \leq$ $\alpha_{I^{*}}$ such that
$0<\alpha_{R_{*}}|\xi|^{2} \leq\left|\sum_{j, k} \alpha_{R_{j k}} \xi_{j} \overline{\xi_{k}}\right| \leq \alpha_{R^{*}}|\xi|^{2}$,
$0 \leq \alpha_{I_{*}}|\xi|^{2} \leq\left|\sum_{j, k} \alpha_{I_{j k}} \xi_{j} \overline{\xi_{k}}\right| \leq \alpha_{I^{*}}|\xi|^{2}$.
Condition (3c) is a general form of absorbing boundary conditions such that electromagnetic fields arriving at $\Gamma$ at certain incident angles do not reflect Sheen (1997). For technical reasons to follow, we assume that $\alpha$ is a complex-valued scalar function which is Lipschitzcontinuous on $\Gamma$.

### 2.2 Function spaces and preliminaries

For an open set $\Omega$ and a real number $r$, let $\left(H^{r}(\Omega),\|\cdot\|_{r, \Omega}\right)$ indicate the usual complex-valued Sobolev space and its norm. In particular, $\left(H^{0}(\Omega),\|\cdot\|_{0, \Omega}\right)$ denotes the usual $L^{2}(\Omega)$-space and its norm, $\left(L^{2}(\Omega),\|\cdot\|_{0, \Omega}\right)$, with the associated inner product
$(\varphi, \psi)=\int_{\Omega} \varphi \bar{\psi} d \mathbf{x}$
Also, for a part $\gamma$ of the boundary $\partial \Omega$ of $\Omega$,
$\langle\varphi, \psi\rangle_{\gamma}=\int_{\gamma} \varphi \bar{\psi} d \Omega$
will mean the inner product on $L^{2}(\gamma)$ with associated norm $|\cdot|_{0, \gamma}$. The following spaces are standard:
$\mathbf{H}(\operatorname{curl} ; \Omega)=\left\{\varphi \in\left[L^{2}(\Omega)\right]^{3}: \boldsymbol{\nabla} \times \varphi \in\left[L^{2}(\Omega)\right]^{3}\right\}$,
$\mathbf{H}(\operatorname{div} ; \Omega)=\left\{\varphi \in\left[L^{2}(\Omega)\right]^{3}: \nabla \cdot \varphi \in\left[L^{2}(\Omega)\right]\right\}$,
equipped with the natural norms
$\|\varphi\|_{\mathbf{H}(\mathbf{c u r l} ; \Omega)}=\left(\|\varphi\|_{0, \Omega}^{2}+\|\boldsymbol{\nabla} \times \varphi\|_{0, \Omega}^{2}\right)^{1 / 2}$,
$\|\varphi\|_{\mathbf{H}(\mathrm{div} ; \Omega)}=\left(\|\varphi\|_{0, \Omega}^{2}+\|\nabla \cdot \varphi\|_{0, \Omega}^{2}\right)^{1 / 2}$.
In the following analysis proper meaning of the traces of $\mathbf{H}(\mathbf{c u r l} ; \Omega)$ and integration by parts with functions in the space $\mathbf{H}(\operatorname{curl} ; \Omega)$ are very important. For this, we will give a brief review of the characterization of the space of tangential traces and tangential components for vector fields in $\mathbf{H}(\mathbf{c u r l} ; \Omega)$ following Buffa, Costabel, and Sheen (2002). Set

$$
\begin{aligned}
& \mathbf{V}=\left[H^{1 / 2}(\Gamma)\right]^{3} ; \quad \mathbf{V}^{\prime}=\left[H^{-1 / 2}(\Gamma)\right]^{3} \\
& H_{*}^{-s}(\Gamma)=\left\{v \in H^{-s}(\Gamma) \mid\langle v, 1\rangle_{s, \Gamma}=0\right\}, s \in[0,1] \\
& \mathbf{L}_{t}^{2}(\Gamma)=\left\{\mathbf{v} \in\left[L^{2}(\Gamma)\right]^{3} \mid \mathbf{v} \cdot \mathbf{v}=0 \text { on } \Gamma\right\}
\end{aligned}
$$

where the space $\mathbf{L}_{t}^{2}(G)$ is identified with the space of fields belonging to the tangent bundle $T \Gamma$ of $\Gamma$ almost everywhere. Also, let $H^{3 / 2}(\Gamma)$ be the trace space of $H^{2}(\Omega)$ endowed with the norm

$$
\|\lambda\|_{3 / 2, \Gamma}=\inf _{\substack{\left.v \in H^{2}(\Omega) \\ v\right|_{G}=\lambda}}\left\{\|v\|_{2, \Omega}\right\}
$$

Then the space $H^{-3 / 2}(\Gamma)$ is defined as the dual space of $H^{3 / 2}(\Gamma)$ with the pivot space $L^{2}(\Gamma)$.

Definition 2.1. The "tangential component trace" mapping $\pi_{\tau}: \mathcal{D}(\bar{\Omega})^{3} \rightarrow \mathbf{L}_{t}^{2}(\Gamma)$ and the "tangential trace" mapping $\gamma_{\tau}: \mathcal{D}(\bar{\Omega})^{3} \rightarrow \mathbf{L}_{t}^{2}(\Gamma)$ are defined as $\mathbf{v} \mapsto-\mathbf{v} \times$ $(\mathbf{v} \times \mathbf{v})$ and $\mathbf{v} \mapsto-\mathbf{v} \times \mathbf{v}$, respectively.

Denoting by $\gamma$ the standard trace operator on the product space $\left[H^{1}(\Omega)\right]^{3} \rightarrow \mathbf{V}$ defined by $\gamma(\mathbf{v})=\left.\mathbf{v}\right|_{\Gamma}$ and by $\gamma^{-1}$ its right inverse, we shall abuse the notations $\pi_{\tau}$ and $\gamma_{\tau}$ for the composite operators $\pi_{\tau} \circ \gamma^{-1}$ and $\gamma_{\tau} \circ \gamma^{-1}$, respectively. Due to the density of $[D(\bar{\Omega})]^{3}$ in $\left[L^{2}(\Gamma)\right]^{3}$, the operators $\pi_{\tau}$ and $\gamma_{\tau}$ can be extended to $\left[L^{2}(\Gamma)\right]^{3}$ linearly and continuously. Now set
$\mathbf{V}_{\gamma}=\gamma_{\tau}(\mathbf{V}) ; \quad \mathbf{V}_{\pi}=\pi_{\tau}(\mathbf{V})$,
which are Hilbert spaces endowed with the norms

$$
\|\lambda\| \mathbf{v}_{\gamma}=\inf _{\substack{v \in \mathbf{V} \\ \gamma_{\tau}(v)=\lambda}}\{\|v\| \mathbf{v}\} ; \quad\|\lambda\|_{\mathbf{V}_{\pi}}=\inf _{\substack{v \in \mathbf{V} \\ \pi_{\tau}(v)=\lambda}}\{\|v\| \mathbf{V}\} .
$$

Let $i_{\pi}: \mathbf{L}_{t}^{2}(\Gamma) \rightarrow \mathbf{L}^{2}(\Gamma)$ and $i_{\gamma}: \mathbf{L}_{t}^{2}(\Gamma) \rightarrow \mathbf{L}^{2}(\Gamma)$ be the adjoint operators of $\pi_{\tau}$ and $\gamma_{\tau}$ respectively. These operators are the identifications of tangent fields with 3D vector fields. Thanks to the Lipschitz assumption, a local system of orthonormal coordinates $\left(\tau_{1}, \tau_{2}, \mathbf{v}\right)$ can be defined at almost every $\mathbf{x} \in \Gamma$. Here, $\tau_{1}$ and $\tau_{2}$ are two orthonormal vectors belonging to the tangent plane for almost every $\mathbf{x} \in \Gamma$, while $v$ is the outer normal to $\Omega$. Of course, the vectors $\tau_{1}$ and $\tau_{2}$ can also be considered as "tangent fields"(sections of the tangent bundle). The operators $i_{\pi}$ and $i_{\gamma}$ can be extended as isomorphisms in the following way:
$i_{\pi}: \mathbf{V}_{\pi}^{\prime} \rightarrow\left(\operatorname{ker}\left(\pi_{\tau}\right) \cap \mathbf{V}\right)^{0}$,
$i_{\gamma}: \mathbf{V}_{\gamma}^{\prime} \rightarrow\left(\operatorname{ker}\left(\gamma_{\tau}\right) \cap \mathbf{V}\right)^{0}$,
where ${ }^{0}$ denotes the polar set (or, the annihilator).
Tangential gradient and curl operators are defined as usual using a localization argument:
$\nabla_{\Gamma}: H^{1}(\Gamma) \rightarrow \mathbf{L}_{t}^{2}(\Gamma) ; \quad \operatorname{curl}_{\Gamma}: H^{1}(\Gamma) \rightarrow \mathbf{L}_{t}^{2}(\Gamma)$,
and the corresponding adjoint operators, which are linear and continuous, are defined

$$
\operatorname{div}_{\Gamma}: \mathbf{L}_{t}^{2}(\Gamma) \rightarrow H_{*}^{-1}(\Gamma) ; \quad \operatorname{curl}_{\Gamma}: \mathbf{L}_{t}^{2}(\Gamma) \rightarrow H_{*}^{-1}(\Gamma)
$$

respectively. Moreover, the following operators are continuous:
$\nabla_{\Gamma}: H^{3 / 2}(\Gamma) \rightarrow \mathbf{V}_{\pi} ; \quad \nabla_{\Gamma}: H^{1 / 2}(\Gamma) \rightarrow \mathbf{V}_{\gamma}^{\prime} ;$
$\operatorname{curl}_{\Gamma}: H^{3 / 2}(\Gamma) \rightarrow \mathbf{V}_{\gamma} ; \quad \operatorname{curl}_{\Gamma}: H^{1 / 2}(\Gamma) \rightarrow \mathbf{V}_{\pi}^{\prime} ;$
satisfying

$$
\|\lambda\|_{H^{1 / 2}(\Gamma) / \mathbf{C}} \leq C\|\nabla \lambda\|_{\mathbf{v}_{\gamma}^{\prime}} ; \quad\|\lambda\|_{H^{1 / 2}(\Gamma) / \mathbf{C}} \leq C\|\nabla \lambda\|_{\mathbf{V}_{\pi}^{\prime}}
$$

Therefore their adjoint operators $\operatorname{div}_{\Gamma}: \mathbf{V}_{\gamma} \rightarrow H_{*}^{-1 / 2}(\Gamma)$ and $\operatorname{curl}_{\Gamma}: \mathbf{V}_{\pi} \rightarrow H_{*}^{-1 / 2}(\Gamma)$ are continuous and surjective. Based on these tangential operators, the LaplaceBeltrami operator $\Delta_{\Gamma}: H^{1}(\Gamma) \rightarrow H_{*}^{-1}(\Gamma)$ is defined by $\Delta_{\Gamma} v=\operatorname{div}_{\Gamma} \nabla_{\Gamma} v$ for any $v \in H^{1}(\Gamma)$.

We are now ready to introduce the traces of $\mathbf{H}(\mathbf{c u r l} ; \Omega)$. Set
$\mathbf{H}^{-1 / 2}\left(\operatorname{div}_{\Gamma} ; \Gamma\right)=\left\{\lambda \in \mathbf{V}_{\pi}^{\prime} \mid \operatorname{div}_{\Gamma}(\lambda) \in H^{-1 / 2}(\Gamma)\right\}$,
$\mathbf{H}^{-1 / 2}\left(\operatorname{curl}_{\Gamma} ; \Gamma\right)=\left\{\lambda \in \mathbf{V}_{\gamma}^{\prime} \mid \operatorname{curl}_{\Gamma}(\lambda) \in H^{-1 / 2}(\Gamma)\right\}$
with the graph norms
$\|\mathbf{v}\|_{\mathbf{H}^{-1 / 2}\left(\operatorname{div}_{\Gamma} ; \Gamma\right)}=\sqrt{\|\mathbf{v}\|_{\mathbf{V}_{\pi}^{\prime}}^{2}+\left\|\operatorname{div}_{\Gamma}(\mathbf{v})\right\|_{-1 / 2, \Gamma}^{2}}$,
$\|\mathbf{v}\|_{\mathbf{H}^{-1 / 2}\left(\operatorname{curl}_{\Gamma} ; \Gamma\right)}=\sqrt{\|\mathbf{v}\|_{\mathbf{V}_{\gamma}^{\prime}}^{2}+\left\|\operatorname{curl}_{\Gamma}(\mathbf{v})\right\|_{-1 / 2, \Gamma}^{2}}$.
Then the following theorem holds:
Theorem 2.1. The operators $\gamma_{\tau}: \mathbf{H}(\operatorname{curl} ; \Omega) \rightarrow$ $\mathbf{H}^{-1 / 2}\left(\operatorname{div}_{\Gamma} ; \Gamma\right)$ and $\pi_{\tau}: \mathbf{H}(\mathbf{c u r l} ; \Omega) \rightarrow \mathbf{H}^{-1 / 2}\left(\operatorname{curl}_{\Gamma} ; \Gamma\right)$ are linear, continuous, and surjective.

The proof of the surjectivity, on the other hand, is based on the proof given by Tartar in Tartar (1997). Let

$$
\begin{align*}
T & :=\left\{\xi \in \mathbf{V}^{\prime} \mid \exists \eta \in H^{-1 / 2}(\Gamma): \forall \phi \in H^{2}(\Omega):\right. \\
& \left.\mathbf{v}^{\prime}\langle\xi, \gamma(\nabla \phi)\rangle \mathbf{v}=\langle\eta, \phi\rangle_{1 / 2, \Gamma}\right\} . \tag{6}
\end{align*}
$$

In Tartar (1997), the tangential trace operator is defined as " $\gamma_{\tau}: \mathbf{H}(\mathbf{c u r l} ; \Omega) \rightarrow T, \mathbf{v} \mapsto \mathbf{v} \times \mathbf{v}$ " and it is proven to be surjective by a localization argument. Here, our setting is different: the ranges of the operators $\pi_{\tau}$ and $\gamma_{\tau}$ defined above are Hilbert spaces of tangent fields. It is shown that the mapping $i_{\pi}$ defined in (5) is indeed an isomorphism between $T$ and $\mathbf{H}^{-1 / 2}\left(\operatorname{div}_{\Gamma} ; \Gamma\right)$, i.e.,
$i_{\pi}\left(\mathbf{H}^{-1 / 2}\left(\operatorname{div}_{\Gamma} ; \Gamma\right)\right) \equiv T$.
Set
$\mathcal{H}(\Gamma):=\left\{p \in H^{1}(\Gamma) / \mathbf{R} \mid \Delta_{\Gamma} p \in H_{\star}^{-1 / 2}(\Gamma)\right\}$.
The trace spaces of $\mathbf{H}(\mathbf{c u r l} ; \Omega)$ have the following Hodge-type decomposition results:

## Theorem 2.2.

$\mathbf{H}^{-1 / 2}\left(\operatorname{div}_{\Gamma} ; \Gamma\right)=\nabla_{\Gamma}(\mathcal{H}(\Gamma)) \oplus \operatorname{curl}_{\Gamma}\left(H^{1 / 2}(\Gamma)\right)$,
$\mathbf{H}^{-1 / 2}\left(\operatorname{curl}_{\Gamma} ; \Gamma\right)=\operatorname{curl}_{\Gamma}(\mathcal{H}(\Gamma)) \oplus \nabla_{\Gamma}\left(H^{1 / 2}(\Gamma)\right)$.

Based on the above decomposition results, a duality can be defined between $\mathbf{H}^{-1 / 2}\left(\operatorname{div}_{\Gamma} ; \Gamma\right)$ and $\mathbf{H}^{-1 / 2}\left(\right.$ curl $\left._{\Gamma} ; \Gamma\right)$ with the pivot space $\mathbf{L}_{t}^{2}(\Gamma)$.

Theorem 2.3. Let $\mathbf{u} \in \mathbf{H}^{-1 / 2}\left(\operatorname{div}_{\Gamma} ; \Gamma\right)$ and $\mathbf{v} \in$ $\mathbf{H}^{-1 / 2}\left(\right.$ curl $\left._{\Gamma} ; \Gamma\right)$ be decomposed as: $\mathbf{u}=\nabla_{\Gamma} \alpha_{u}+$ $\operatorname{curl}_{\Gamma} \beta_{\mathbf{u}}, \mathbf{v}=\nabla_{\Gamma} \beta_{\mathbf{v}}+\operatorname{curl}_{\Gamma} \alpha_{\mathbf{v}}$ with $\beta_{\mathbf{u}}, \beta_{\mathbf{v}} \in H^{1 / 2}(\Gamma)$ and $\alpha_{\mathbf{u}}, \alpha_{\mathbf{v}} \in \mathcal{H}(\Gamma)$. Then, we have

$$
\begin{equation*}
\gamma\langle\mathbf{u}, \mathbf{v}\rangle_{\pi}:=-\left\langle\Delta_{\Gamma} \alpha_{u}, \beta_{v}\right\rangle_{1 / 2, \Gamma}+\left\langle\Delta_{\Gamma} \alpha_{v}, \beta_{u}\right\rangle_{1 / 2, \Gamma} . \tag{7}
\end{equation*}
$$

Given $\mathbf{u} \in \mathbf{H}(\mathbf{c u r l} ; \Omega)$, recall the decompositions $\mathbf{u}=$ $\boldsymbol{\Phi}+\boldsymbol{\nabla} p$ with $\boldsymbol{\Phi} \in\left[H^{1}(\Omega)\right]^{3}, p \in H^{1}(\Omega)$. Based on these decompositions, the integration by parts formula holds:

Theorem 2.4. Given $\mathbf{u}, \mathbf{v} \in \mathbf{H}(\mathbf{c u r l} ; \Omega)$, let $\mathbf{u}=\boldsymbol{\Phi}+$ $\nabla p, \mathbf{v}=\Psi+\boldsymbol{\nabla} q$ with $\boldsymbol{\Phi}, \boldsymbol{\Psi} \in\left[H^{1}(\Omega)\right]^{3}, p, q \in H^{1}(\Omega)$.

$$
\begin{equation*}
\int_{\Omega}\{\nabla \times \mathbf{u} \cdot \overline{\mathbf{v}}-\mathbf{u} \cdot \boldsymbol{\nabla} \times \overline{\mathbf{v}}\} d \mathbf{x}=-\gamma\left(\gamma_{\tau}(\mathbf{u}), \pi_{\tau}(\mathbf{v})\right\rangle_{\pi} \tag{8}
\end{equation*}
$$

where the boundary term can be interpreted as

$$
\begin{aligned}
& \gamma\left\langle\gamma_{\tau}(\mathbf{u}), \pi_{\tau}(\mathbf{v})\right\rangle_{\pi}=\int_{\Gamma} \gamma_{\tau}(\boldsymbol{\Phi}) \cdot \overline{\pi_{\tau}(\boldsymbol{\Psi})}+\mathbf{v}_{\gamma}^{\prime}\left\langle\boldsymbol{\nabla}_{\Gamma} q, \gamma_{\tau}(\boldsymbol{\Psi})\right\rangle_{\mathbf{v}_{\gamma}} \\
& \quad+\mathbf{v}_{\pi}^{\prime}\left\langle\operatorname{curr}_{\Gamma} p, \pi_{\tau}(\boldsymbol{\Psi})\right\rangle_{\mathbf{V}_{\pi}} \\
& \quad=\int_{\Gamma} \gamma_{\tau}(\boldsymbol{\Phi}) \cdot \overline{\pi_{\tau}(\boldsymbol{\Psi})}-\left\langle\operatorname{div}_{\Gamma} \gamma_{\tau}(\boldsymbol{\Phi}), q\right\rangle_{1 / 2, \Gamma} \\
& \quad+\left\langle\operatorname{curl}_{\Gamma} \pi_{\tau}(\boldsymbol{\Psi}), p\right\rangle_{1 / 2, \Gamma} .
\end{aligned}
$$

In what follows we will use the notation $\mathbf{v} \times \varphi$ to denote $-\gamma_{\tau}(\varphi)$ if $\varphi$ is in the appropriate space $\mathbf{H}(\operatorname{curl} ; \Omega)$. Recall the classical Green's formula

$$
\begin{aligned}
& (\boldsymbol{\nabla} \times \varphi, \psi)-(\varphi, \boldsymbol{\nabla} \times \psi)=\langle\boldsymbol{v} \times \varphi, \psi\rangle_{\partial \Omega} \\
& \quad=\left\langle\mathbf{v} \times \varphi, \pi_{\tau} \psi\right\rangle_{\partial \Omega}
\end{aligned}
$$

for all $\varphi \in \mathbf{H}(\mathbf{c u r l} ; \Omega), \psi \in\left[H^{1}(\Omega)\right]^{3}$. If both $\varphi$ and $\psi$ belong to $\mathbf{H}(\mathbf{c u r l} ; \Omega), \boldsymbol{v} \times \varphi$ and $\pi_{\tau} \psi$ belong to $\left[H^{-1 / 2}(\partial \Omega)\right]^{3}$ and therefore Theorem 2.4 generalizes the above classical integration by formula.

Remark 2.1. In Sheen (1992) the integration by parts formula for a Lipschitz domain $\Omega$ is proved directly by using a density argument with a different interpretation of the boundary integral term in (8):

$$
\begin{align*}
& (\boldsymbol{\nabla} \times \varphi, \psi)-(\varphi, \boldsymbol{\nabla} \times \psi)=\langle\mathbf{v} \times \varphi, \psi\rangle_{\partial \Omega} \\
& \quad=\left\langle\mathbf{v} \times \varphi, \pi_{\tau} \psi\right\rangle_{\partial \Omega} \tag{9}
\end{align*}
$$

for all $\varphi, \psi \in \mathbf{H}(\mathbf{c u r l} ; \Omega)$, where, the boundary integral term $\left\langle\boldsymbol{v} \times \varphi, \pi_{\tau} \psi\right\rangle_{\partial \Omega}$ is understood as $\left\langle(\boldsymbol{v} \times \varphi) \cdot \overline{\pi_{\tau} \psi}, 1\right\rangle$, the duality pairing between $\mathbf{v} \times \varphi \cdot \overline{\pi_{\tau} \psi} \in \operatorname{Lip}(\partial \Omega)^{\prime}$
and $1 \in \operatorname{Lip}(\partial \Omega)$. Indeed, in Sheen (1992) the integration by parts formula (8) was proved for a general class of linear first order differential operators. $\mathcal{L}=\sum_{j=1}^{N} \mathbf{A}_{j}(x) \frac{\partial}{\partial x_{j}}$, where the $\mathbf{A}_{j}$ 's are $k \times k$ matrices with uniformly Lipschitz-continuous components on $\bar{\Omega}$ and $\mathcal{L}^{*}$ be the formal adjoint of $\mathcal{L}$ given by $\mathcal{L}^{*}=$ $-\sum_{j=1}^{N} \frac{\partial}{\partial x_{j}} \mathbf{A}_{j}^{*}(x)$, with $\mathbf{A}_{j}^{*}$ 's being the adjoint matrices of $\mathbf{A}_{j}$ 's. Then the following Hilbert space $H(\mathcal{L} ; \Omega)=$ $\left\{\mathbf{u} \in\left[L^{2}(\Omega)\right]^{k} ; \mathcal{L} \mathbf{u} \in\left[L^{2}(\Omega)\right]^{k}\right\}$ is endowed with the inner product and the norm
$(\mathbf{u}, \mathbf{v})_{H(\mathcal{L} ; \Omega)}=(\mathbf{u}, \mathbf{v})+(\mathcal{L} \mathbf{u}, \mathcal{L} \mathbf{v})$,
$\|\mathbf{u}\|_{H(\mathcal{L} ; \Omega)}=\left\{\|\mathbf{u}\|^{2}+\|\mathcal{L} \mathbf{u}\|^{2}\right\}^{\frac{1}{2}}$.
The analogues for $\mathcal{L}^{*}$ holds. Denoting $\mathbf{A}_{\boldsymbol{v}}=\sum_{j=1}^{N} \mathrm{v}_{j} \mathbf{A}_{j}$, the main result of Sheen (1992) is given as follows: the map $\{\mathbf{u}, \mathbf{v}\} \mapsto \mathbf{A}_{\mathbf{v}} \mathbf{u} \cdot \overline{\mathbf{v}}$ from $[\mathcal{D}(\bar{\Omega})]^{k} \times[\mathcal{D}(\bar{\Omega})]^{k}$ into $\operatorname{Lip}(\Gamma)^{\prime}$ can be extended by continuity to a continuous sesquilinearmap from $H(\mathcal{L} ; \Omega) \times H\left(\mathcal{L}^{*} ; \Omega\right)$ into $\operatorname{Lip}(\Gamma)^{\prime}$; moreover, for all $\mathbf{u} \in H(\mathcal{L} ; \Omega)$ and $\mathbf{v} \in H\left(\mathcal{L}^{*} ; \Omega\right)$, the following Green's formula holds:
$(\mathcal{L} \mathbf{u}, \mathbf{v})-\left(\mathbf{u}, \mathcal{L}^{*} \mathbf{v}\right)=\left\{A_{\mathbf{v}} \mathbf{u}, \mathbf{v}\right\}_{\Gamma} \equiv_{\operatorname{Lip}(\Gamma)^{\prime}}\left\langle\mathbf{A}_{\mathbf{v}} \mathbf{u} \cdot \overline{\mathbf{v}}, 1\right\rangle_{\operatorname{Lip}(\Gamma)}$,
where $_{\operatorname{Lip}(\Gamma)^{\prime}}\langle\cdot, \cdot\rangle_{\operatorname{Lip}(\Gamma)}$ denotes the duality paring between $\operatorname{Lip}(\Gamma)^{\prime}$ and $\operatorname{Lip}(\Gamma)$.

An immediate result of Theorem 2.3 is useful:
Proposition 2.1. Let $\mathbf{u} \in \mathbf{H}^{-1 / 2}\left(\operatorname{div}_{\Gamma} ; \Gamma\right) \cap$ $\mathbf{H}^{-1 / 2}\left(\right.$ curl $\left._{\Gamma} ; \Gamma\right)$ with $\mathbf{v} \cdot \mathbf{u}=0$. Then $\mathbf{u} \in \mathbf{L}_{t}^{2}(\Gamma)$.

Proof. Suppose $\mathbf{u} \in \mathbf{H}^{-1 / 2}\left(\operatorname{div}_{\Gamma} ; \Gamma\right) \cap \mathbf{H}^{-1 / 2}\left(\operatorname{curl}_{\Gamma} ; \Gamma\right)$. Then $\mathbf{u}$ has the following decomposition: $\mathbf{u}=\boldsymbol{\nabla}_{\Gamma} \alpha+$ $\operatorname{curl}_{\Gamma} \beta$ and $\mathbf{u}=\nabla_{\Gamma} \beta^{\prime}+\operatorname{curl}_{\Gamma} \alpha^{\prime}$ with $\alpha, \alpha^{\prime} \in \mathcal{H}(\Gamma)$ and $\beta, \beta^{\prime} \in H^{1 / 2}(\Gamma)$. Moreover, the duality (7) implies that

$$
\langle\mathbf{u}, \mathbf{u}\rangle_{0, \Gamma}=-\gamma\langle\mathbf{u}, \mathbf{u}\rangle_{\pi}=-\left\langle\Delta_{\Gamma} \alpha, \beta^{\prime}\right\rangle_{1 / 2, \Gamma}+\left\langle\Delta_{\Gamma} \alpha^{\prime}, \beta\right\rangle_{1 / 2, \Gamma},
$$

which is bounded by the definition of duality. Thus, $\mathbf{u} \in$ $\mathbf{L}_{t}^{2}(\Gamma)$.

### 2.3 Existence and uniqueness results

We begin with the following lemma:
Lemma 2.1. If $\{E, H\} \in[\mathbf{H}(\mathbf{c u r l} ; \Omega)]^{2}$ satisfies (3), the boundary terms $\alpha \pi_{\tau} E$ and $\boldsymbol{v} \times H$ belong to $\mathbf{L}_{t}^{2}(\Gamma)$.

Proof. Indeed, let $\widetilde{\alpha}$ be an invertible Lipschitz extension of $\alpha$ onto $\bar{\Omega}$ such that $\widetilde{\alpha}=\alpha$ on $\Gamma$. Then, taking the inner product of the equation (3a) with $\left(\widetilde{\alpha}^{*}\right)^{-1} E$ and applying (9) with using the boundary condition (3c), we obtain

$$
\begin{aligned}
0 & =\left(F,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right)-\left((i \omega \varepsilon+\sigma) E,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right) \\
& +\left(H, \nabla \times\left(\widetilde{\alpha}^{*}\right)^{-1} E\right) \\
& =\left(F,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right)-\left((i \omega \varepsilon+\sigma) E,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right) \\
& +\left(\nabla \times H,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right)-\left\langle\boldsymbol{v} \times H,\left(\alpha^{*}\right)^{-1} \pi_{\tau} E\right\rangle_{\Gamma} \\
& =\left(F,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right)-\left((i \omega \varepsilon+\sigma) E,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right) \\
& +\left(\nabla \times H,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right)+\left\langle\pi_{\tau} E, \pi_{\tau} E\right\rangle_{\Gamma} .
\end{aligned}
$$

Hence,

$$
\begin{aligned}
& \left|\pi_{\tau} E\right|_{0, \Gamma}^{2} \leq_{\gamma}\left\langle\pi_{\tau} E, \pi_{\tau} E\right\rangle_{\pi} \\
& \quad \leq\left|\left(F,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right)\right|+\left|\left((i \omega \varepsilon+\sigma) E,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right)\right| \\
& \quad+\left|\left(\nabla \times H,\left(\widetilde{\alpha}^{*}\right)^{-1} E\right)\right| \\
& \quad \leq C\left[\|E\|_{0, \Omega}^{2}+\|H\|_{\mathbf{H}(\mathbf{c u r l} ; \Omega)}^{2}+\|F\|_{0, \Omega}^{2}\right] .
\end{aligned}
$$

Therefore we see that $\pi_{\tau} E$ actually belongs to $\mathbf{L}_{t}^{2}(\Gamma)$. The boundary condition (3c) leads to $\mathbf{v} \times H \in \mathbf{L}_{t}^{2}(\Gamma)$.

The existence and uniqueness results for Problem (3) without the term $\varepsilon$ are given in Santos and Sheen (2000) with the use of integration by parts formula in the sense of Remark 2.1. However, it is straightforward to check all the arguments given there are still valid with the term $\varepsilon$ included and the use of integration by parts formula in the sense of Theorem 2.4:

Theorem 2.1. Assume that $\sigma_{*}>0$. Let $F, G \in\left[L^{2}(\Omega)\right]^{3}$ and $\omega \neq 0$. Then, there exists a unique electromagnetic field $\{E, H\} \in[\mathbf{H}(\mathbf{c u r l} ; \Omega)]^{2}$ satisfying (3) with $\pi_{\tau} E, \mathbf{v} \times$ $H \in \mathbf{L}_{t}^{2}(\Gamma)$. If, in addition, $F$ and $G$ belong to $H(\operatorname{div} ; \Omega)$ and $\varepsilon, \sigma$, and $\mu$ are Lipschitz-continuous on $\bar{\Omega}$, then $E$ and $H$ belong to $\left[H^{1 / 2}(\Omega)\right]^{3}$; more precisely, $E$ and $H$ belong to $[H(\mathbf{c u r l} ; \Omega) \cap H(\operatorname{div} ; \Omega)]$ with boundary values in $\left[L^{2}(\Gamma)\right]^{3}$

In getting the weak formulation we will see the cross term $\left\langle\alpha \pi_{\tau} E, \pi_{\tau} \phi\right\rangle_{\Gamma}$ with $E, \phi \in \mathbf{H}(\mathbf{c u r l} ; \Omega)$. In order to provide a meaning to this term, the correct test function space for the weak problem will be
$\mathbf{H}^{*}(\operatorname{curl} ; \Omega)=\left\{v \in \mathbf{H}(\operatorname{curl} ; \Omega): \alpha \pi_{\tau} v=\mathbf{v} \times \varphi\right.$
for some $\varphi \in \mathbf{H}(\operatorname{curl} ; \Omega)\}$.

By applying (9), a mixed weak formulation of Problem (1) follows immediately: find $\{E, H\} \in \mathbf{H}^{*}(\operatorname{curl} ; \Omega) \times$ $\left[L^{2}(\Omega)\right]^{3}$ such that

$$
\begin{align*}
& ((i \omega \varepsilon+\sigma) E, \varphi)-(H, \nabla \times \varphi)+\left\langle\alpha \pi_{\tau} E, \pi_{\tau} \varphi\right\rangle_{\Gamma}=(F, \varphi), \\
& \varphi \in H^{*}(\operatorname{curl} ; \Omega), \tag{10a}
\end{align*}
$$

$$
\begin{align*}
& i \omega(\mu H, \psi)+(\nabla \times E, \psi)=(G, \psi) \\
& \quad \psi \in\left[L^{2}(\Omega)\right]^{3} \tag{10b}
\end{align*}
$$

It is easy to see that if $\{E, H\}$ satisfies the above weak problem, it is then a solution of the differential equations (3) with the boundary condition (3c) in the sense of $\mathbf{H}^{-1 / 2}\left(\operatorname{curl}_{\Gamma} ; \Gamma\right)$.

## 3 The nonconforming method

In this section we summarize some results on nonconforming mixed finite element space based on threedimensional rectangular domain introduced in Douglas Jr., Santos, and Sheen (2000).

### 3.1 Nonconforming mixed finite element space

Let $\widehat{K}$ be the reference cube $[-1,1]^{3}$ and let $\widehat{Q}(\widehat{K})=\widehat{Q}_{x} \times$ $\widehat{Q}_{y} \times \widehat{Q}_{z}$, where

$$
\begin{aligned}
& \widehat{Q}_{x}=\operatorname{Span}\left\{1, y, z,\left(y^{2}-\frac{5}{3} y^{4}\right)-\left(z^{2}-\frac{5}{3} z^{4}\right)\right\} \\
& \widehat{Q}_{y}=\operatorname{Span}\left\{1, z, x,\left(z^{2}-\frac{5}{3} z^{4}\right)-\left(x^{2}-\frac{5}{3} x^{4}\right)\right\}
\end{aligned}
$$

$$
\widehat{Q}_{z}=\operatorname{Span}\left\{1, x, y,\left(x^{2}-\frac{5}{3} x^{4}\right)-\left(y^{2}-\frac{5}{3} y^{4}\right)\right\}
$$

Let $m_{i}, i=1, \cdots, 6$, be the centroid of the $i^{\text {th }}$ face of $\widehat{K}$. Then, for $\varphi \in \widehat{Q}(\widehat{K})$, we consider the following local degrees of freedom :
$\Sigma(\varphi)=\left\{\left(\pi_{\tau} \varphi\right)\left(m_{i}\right) \mid i=1, \cdots, 6\right\}$.
Then, a local interpolant $\widehat{\pi}: H^{2}(\widehat{K}) \rightarrow \widehat{Q}(\widehat{K})$ is defined as follows:
$\pi_{\tau}(\widehat{\pi} \varphi-\varphi)\left(m_{i}\right)=0, \quad i=1, \cdots, 6$.
Let $\widehat{S}(\widehat{K})=\widehat{S}_{x} \times \widehat{S}_{y} \times \widehat{S}_{z}$, where
$\widehat{S}_{x}=\operatorname{Span}\left\{1, y-\frac{10}{3} y^{3}, z-\frac{10}{3} z^{3}\right\}$,
$\widehat{S}_{y}=\operatorname{Span}\left\{1, z-\frac{10}{3} z^{3}, x-\frac{10}{3} x^{3}\right\}$,
$\widehat{S}_{z}=\operatorname{Span}\left\{1, x-\frac{10}{3} x^{3}, y-\frac{10}{3} y^{3}\right\}$,
and a local interpolant $\widehat{P}: \mathbf{H}(\mathbf{c u r l} ; \widehat{K}) \rightarrow \widehat{S}(\widehat{K})$ be defined by the following rules: for $l=x, y, z$,
$\int_{\widehat{K}}\left(\widehat{P} \psi_{l}-\psi_{l}\right) d x d y d z=0$,
$\int_{\widehat{K}} \boldsymbol{\nabla} \times\left(\widehat{P} \psi_{l}-\psi_{l}\right) d x d y d z=0$,
for all $\psi=\left(\psi_{x}, \psi_{y}, \psi_{z}\right)$, Note that (12) provides the nine degrees of freedom needed to determine an element in $\widehat{S}(\widehat{K})$ and that
$\nabla \times \widehat{Q}=\widehat{S}$.
The following lemma is trivial but useful Douglas Jr., Santos, and Sheen (2000).

Lemma 3.1. The degrees of freedom (11) and (12) determine, respectively, $\varphi \in \widehat{Q}(\widehat{K})$ and $\psi \in \widehat{S}(\widehat{K})$ uniquely.

The following proposition states an immediate but fundamental property of $\widehat{Q}$ and $\widehat{S}$ that is important in obtaining effective nonconforming methods Douglas Jr., Santos, and Sheen (2000); )Douglas Jr., Santos, Sheen, and Ye (1999).
Proposition 3.1. If an element of $\pi_{\tau} \widehat{Q}$ or $\pi_{\tau} \widehat{S}$ vanishes at the centroid of a face of $\widehat{K}$, it is orthogonal to constants on that face.

For $0<h<1$, let $\Omega$ be decomposed into nonoverlapping three-dimensional rectangular hexahedra $\left\{\Omega_{j}: j=\right.$ $1, \cdots, J\}$ with their edges bounded by $h$ :
$\bar{\Omega}=\cup_{j} \bar{\Omega}_{j}, \quad \Omega_{j} \cap \Omega_{k}=0, \quad j \neq k$.
Then $Q\left(\Omega_{j}\right)$ and $S\left(\Omega_{j}\right)$ are defined by scaling and translating from $\widehat{Q}$ and $\widehat{S}$. Let
$\Gamma_{j}=\partial \Omega_{j} \cap \Gamma, \quad \Gamma_{j k}=\partial \Omega_{j} \cap \partial \Omega_{k}=\Gamma_{k j}$.

### 3.2 The nonconforming mixed finite element method

Set
$\mathbf{V}^{h}=\left\{\varphi \in\left[L^{2}(\Omega)\right]^{3}: \varphi_{j}:=\varphi \mid \Omega_{j} \in Q\left(\Omega_{j}\right)\right.$,
$\varphi_{j}\left(m_{j k}\right)=\varphi_{k}\left(m_{j k}\right)$,
$m_{j k}$ being the centroid of $\left.\Gamma_{j k} \forall\{j, k\} \in I^{h}\right\}$,
$W^{h}=\left\{\psi \in\left[L^{2}(\Omega)\right]^{3}: \psi \mid \Omega_{j} \in S\left(\Omega_{j}\right)\right\}$,
where $I^{h}$ denotes the index set of all internal interfaces $\Gamma_{j k}$.
The nonconforming mixed finite element method is then defined as follows: find $\left\{E^{h}, H^{h}\right\} \in \mathbf{V}^{h} \times W^{h}$ such that

$$
\begin{align*}
& \left((i \omega \varepsilon+\sigma) E^{h}, \varphi\right)-\sum_{j}\left[\left(H^{h}, \boldsymbol{\nabla} \times \varphi\right)_{j}\right. \\
& \left.\quad+\sum_{k}\left\langle\mathbf{v}_{j} \times E_{j}^{h}, \pi_{\tau} \varphi\right\rangle_{\Gamma_{j k}}\right]+\left\langle\alpha \pi_{\tau} E^{h}, \pi_{\tau} \varphi\right\rangle_{\Gamma} \\
& \quad=(F, \varphi), \quad \varphi \in \mathbf{V}^{h}, \tag{13a}
\end{align*}
$$

$i \omega\left(\mu H^{h}, \psi\right)+\sum_{j}\left(\nabla \times E^{h}, \psi\right)_{j}=(G, \psi), \psi \in W^{h}$.
Denoting by $\|\cdot\|_{0, h}$ the broken $L(\boldsymbol{\Omega})$-norm, we then have the following a priori error estimate:

Theorem 3.1. Suppose that $\sigma_{*}>0$. Let $\{E, H\}$ and $\left\{E^{h}, H^{h}\right\}, 0<h<1$, be the solutions to (10) and (13), respectively. Then,

$$
\begin{align*}
& \left\|E-E^{h}\right\|_{0}+\left\|H-H^{h}\right\|_{0}+\left\|\nabla \times\left(E-E^{h}\right)\right\|_{0, h} \\
& \quad \leq C h^{1 / 2}\left(\|E\|_{2}+h^{1 / 2}\|H\|_{1}\right) \tag{14}
\end{align*}
$$

The proof of the above theorem is given in Douglas Jr., Santos, and Sheen (2000), where $\varepsilon^{*}=0$ is assumed, but including the case $\varepsilon_{*}>0$ does not change any argument given there.

### 3.3 Hybridization

Denote by $\langle\langle\cdot, \cdot\rangle\rangle_{\Gamma_{j k}}$ the approximation to $\langle\cdot, \cdot\rangle_{\Gamma_{j k}}$ obtained by using the mid-point rule on $\Gamma_{j k}$ so that
$\langle\langle u, v\rangle\rangle_{\Gamma_{j k}}=\left|\Gamma_{j k}\right|(u \bar{v})\left(m_{j k}\right)$,
where $m_{j k}$ and $\left|\Gamma_{j k}\right|$ denote the centroid and the measure of $\Gamma_{j k}$, respectively.
Following Arnold and Brezzi (1985); Fraeijs de Veubeke (1965); Fraeijs de Veubeke (1975), the hybridization of the procedure (13) is given by introducing the space $\tilde{\Lambda}^{h}$ of Lagrange multipliers associating its elements with $\mathbf{v}_{j} \times$ $E_{j}\left(m_{j k}\right)$ on $\Gamma_{j k}$. Thus set

$$
\begin{gathered}
\mathcal{N} C_{-1}^{h}=\left\{\varphi \in\left[L^{2}(\Omega)\right]^{3}:\left.\varphi\right|_{\Omega_{j}} \in Q\left(\Omega_{j}\right)\right\} \\
\tilde{\Lambda}^{h}=\left\{\tilde{\lambda}^{h}:\left.\tilde{\lambda}^{h}\right|_{\Gamma_{j k}}=\tilde{\lambda}_{j k} \in\left[P_{0}\left(\Gamma_{j k}\right)\right]^{2}\right. \\
\left.\tilde{\lambda}_{j k}+\tilde{\lambda}_{k j}=0 \text { for all } j k \in I^{h}\right\},
\end{gathered}
$$

and $P_{0}\left(\Gamma_{j k}\right)$ the set of constant functions defined on $\Gamma_{j k}$. The global hybridized nonconforming mixed finite element method is then defined as follows: find $\left\{\tilde{E}^{h}, \tilde{H}^{h}, \tilde{\lambda}^{h}\right\}$ $\in \mathcal{N} C_{-1}^{h} \times W^{h} \times \tilde{\Lambda}^{h}$ such that

$$
\begin{align*}
& \left((i \omega \varepsilon+\sigma) \tilde{E}^{h}, \varphi\right)-\sum_{j}\left[\left(\tilde{H}^{h}, \nabla \times \varphi\right)_{j}\right. \\
& \left.\quad+\sum_{k}\left\langle\left\langle\tilde{\lambda}^{h}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j k}}\right]+\left\langle\left\langle\alpha \pi_{\tau} \tilde{E}^{h}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma} \\
& \quad=(F, \varphi), \quad \varphi \in \mathcal{N} C_{-1}^{h}, \tag{15a}
\end{align*}
$$

$i \omega\left(\mu \tilde{H}^{h}, \psi\right)+\sum_{j}\left(\nabla \times \tilde{E}^{h}, \psi\right)_{j}=(G, \psi), \quad \psi \in W^{h}$,
$\sum_{j, k \in I^{h}}\left\langle\left\langle\theta, \pi_{\tau} \tilde{E}^{h}\right\rangle\right\rangle_{\Gamma_{j k}}=0, \quad \theta \in \tilde{\Lambda}^{h}$.
$i \omega\left(\mu H_{j}, \psi\right)_{j}+\left(\nabla \times E_{j}, \psi\right)_{j}=\left(G_{j}, \psi\right)_{j}, \psi \in\left[L^{2}\left(\Omega_{j}\right)\right]^{3}$.
Set $\mathbf{V}_{j}^{h}=\left.\mathcal{N} C_{-1}^{h}\right|_{\Omega_{j}}$ and $W_{j}^{h}=\left.W^{h}\right|_{\Omega_{j}}$. Also let $\Lambda^{h}$ be the set of Lagrange multipliers associated with $\left(\boldsymbol{v}_{j} \times\right.$ $\left.H_{j}\right)\left(m_{j k}\right)$ on $\Gamma_{j k}$ given in by
$\Lambda^{h}=\left\{\lambda^{h}\left|\lambda^{h}\right|_{\Gamma_{j k}}=\lambda_{j k}^{h} \in \Lambda_{j k} \quad \forall\{j, k\} \in I^{h}\right\}$,
where $\Lambda_{j k}=\left[P_{0}\left(\Gamma_{j k}\right)\right]^{2}$, and set
$\Lambda_{j}^{h}=\left\{\lambda_{j}^{h}\left|\lambda_{j}^{h}\right|_{j k}=\lambda_{j k}^{h} \in \Lambda_{j k} \quad \forall k\right\}$.
From now on associate $\left(\mathbf{v}_{j} \times H_{j}\right)\left(m_{j k}\right)$ on $\Gamma_{j k}$ with $\lambda_{j k}$.
We are now in a position to define a Seidel-type domain decomposition iterative procedure. Choose an initial guess $\left\{E_{j}^{h, 0}, H_{j}^{h, 0}, \lambda_{j}^{h, 0}\right\} \in \mathbf{V}_{j}^{h} \times W_{j}^{h} \times \Lambda_{j}$ for all $j$. Then, iteratively for $n \geq 1$, find $\left\{E_{j}^{h, n}, H_{j}^{h, n}, \lambda_{j}^{h, n}\right\} \in \mathbf{V}_{j}^{h} \times W_{j}^{h} \times$ $\Lambda_{j}^{h}$, for $j=1,2, \cdots$, fulfilling

$$
\begin{align*}
& \left((i \omega \varepsilon+\sigma) E_{j}^{h, n}, \varphi\right)_{j}-\left(H_{j}^{h, n}, \boldsymbol{\nabla} \times \varphi\right)_{j} \\
& \quad-\sum_{k}\left\langle\left\langle\lambda_{j k}^{h, n}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j k}}+\left\langle\left\langle\alpha \pi_{\tau} E_{j}^{h, n}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j}} \\
& \quad=\left(F_{j}, \varphi\right)_{j}, \quad \varphi \in \mathbf{V}_{j}^{h}, \tag{19a}
\end{align*}
$$

$$
\begin{align*}
& i \omega\left(\mu H_{j}^{h, n}, \psi\right)_{j}+\left(\nabla \times E_{j}^{h, n}, \psi\right)_{j}=\left(G_{j}, \Psi\right)_{j} \\
& \psi \in W_{j}^{h} \tag{19b}
\end{align*}
$$

$\lambda_{j k}^{h, n}=-\lambda_{k j}^{h, n^{*}}+\beta_{j k}\left(\pi_{\tau} E_{k}^{h, n^{*}}-\pi_{\tau} E_{j}^{h, n}\right)\left(m_{j k}\right)$,

$$
\begin{equation*}
\text { for all } k, \tag{19c}
\end{equation*}
$$

where $n^{*}$ is defined by
$n^{*}= \begin{cases}n-1, & j<k \\ n, & j>k\end{cases}$

Remark 4.1. It should be remarked that in the updating procedure (19c) data in use passed from neighboring subdomains are most up-to-date; in this sense (19) is said to be a Seidel-type scheme. Instead of the above, if $n^{*}$ were defined as $n-1$ for all $j, k$, the procedure (19) should be regarded as a Jacobi-type scheme in the sense that the updating procedure (19c) would use informations from the past step; this procedure has been studied in Douglas Jr., Santos, and Sheen (2000).

### 4.2 Convergence of the iterative procedure

The localized Problem (16)-(18) is weakly formulated as follows: find $\left\{\tilde{E}_{j}^{h}, \tilde{H}_{j}^{h}, \tilde{\lambda}_{j}^{h}\right\} \in \mathbf{V}_{j}^{h} \times W_{j}^{h} \times \Lambda_{j}^{h}$ such that

$$
\begin{align*}
& \left((i \omega \varepsilon+\sigma) \tilde{E}_{j}^{h}, \varphi\right)_{j}-\left(\tilde{H}_{j}^{h}, \boldsymbol{\nabla} \times \varphi\right)_{j}-\sum_{k}\left\langle\left\langle\tilde{\lambda}_{j k}^{h}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j k}} \\
& \quad+\left\langle\left\langle\alpha \pi_{\tau} \tilde{E}_{j}^{h}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j}}=\left(F_{j}, \varphi\right)_{j}, \quad \varphi \in \mathbf{V}_{j}^{h}, \tag{20a}
\end{align*}
$$

$i \omega\left(\mu \tilde{H}_{j}^{h}, \psi\right)_{j}+\left(\boldsymbol{\nabla} \times \tilde{E}_{j}^{h}, \psi\right)_{j}=\left(G_{j}, \psi\right)_{j}, \psi \in W_{j}^{h}$,
$\tilde{\lambda}_{j k}^{h}=-\tilde{\lambda}_{k j}^{h}+\beta_{j k}\left(\pi_{\tau} \tilde{E}_{k}^{h}-\pi_{\tau} \tilde{E}_{j}^{h}\right)\left(m_{j k}\right)$, for all $k$.
We will show the convergence of the solution $\left\{E_{j}^{h, n}, H_{j}^{h, n}, \lambda_{j}^{h, n}\right\}$ of Problem (19) to the solution $\left\{\tilde{E}_{j}^{h}, \tilde{H}_{j}^{h}, \tilde{\lambda}_{j}^{h}\right\}$ of Problem (20). We will restrict the convergence proof to the case $\beta_{j k}=\beta$. for all $j, k$. Set, for all $j$,

$$
\begin{equation*}
u_{j}^{n}=E_{j}^{h, n}-\tilde{E}_{j}^{h}, \quad v_{j}^{h}=H_{j}^{h, n}-\tilde{H}_{j}^{h}, \quad \theta_{j}^{n}=\lambda_{j}^{h, n}-\tilde{\lambda}_{j}^{h} . \tag{21}
\end{equation*}
$$

Then a subtraction of (20) from (19) gives the iteration error equations:

$$
\begin{align*}
& \left((i \omega \varepsilon+\sigma) u_{j}^{n}, \varphi\right)_{j}-\left(v_{j}^{n}, \boldsymbol{\nabla} \times \varphi\right)_{j}-\sum_{k}\left\langle\left\langle\theta_{j k}^{n}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j k}} \\
& \quad+\left\langle\left\langle\alpha \pi_{\tau} u_{j}^{n}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j}}=0, \varphi \in \mathbf{V}_{j}^{h}, \tag{22a}
\end{align*}
$$

$i \omega\left(\mu \nu_{j}^{n}, \psi\right)_{j}+\left(\nabla \times u_{j}^{n}, \psi\right)_{j}=0, \quad \psi \in W_{j}^{h}$,
$\theta_{j k}^{n}=-\theta_{k j}^{n^{*}}+\beta\left(\pi_{\tau} u_{k}^{n^{*}}-\pi_{\tau} u_{j}^{n}\right)\left(m_{j k}\right) \forall k$.
Choose $\varphi=u_{j}^{n}$ and $\psi=v_{j}^{n}$ in (22a) and (22b), respectively. Then we have

$$
\begin{gather*}
\sum_{k}\left\langle\left\langle\theta_{j k}^{n}, \pi_{\tau} u_{j}^{n}\right\rangle\right\rangle_{\Gamma_{j k}}=\left((i \omega \varepsilon+\sigma) u_{j}^{n}, u_{j}^{n}\right)_{j} \\
\quad-i \omega\left(\mu v_{j}^{n}, v_{j}^{n}\right)_{j}+\left\langle\left\langle\alpha \pi_{\tau} u_{j}^{n}, \pi_{\tau} u_{j}^{n}\right\rangle\right\rangle_{\Gamma_{j}} \tag{23}
\end{gather*}
$$

Taking the real and imaginary parts in the above equation, we get

$$
\begin{align*}
& \operatorname{Re} \sum_{k}\left\langle\left\langle\theta_{j k}^{n}, \pi_{\tau} u_{j}^{n}\right\rangle\right\rangle_{\Gamma_{j k}}=\left(\sigma u_{j}^{n}, u_{j}^{n}\right)_{j} \\
& \quad+\left\langle\left\langle\alpha_{R} \pi_{\tau} u_{j}^{n}, \pi_{\tau} u_{j}^{n}\right\rangle\right\rangle_{\Gamma_{j}} \tag{24}
\end{align*}
$$

and

$$
\begin{align*}
& \operatorname{Im} \sum_{k}\left\langle\left\langle\theta_{j k}^{n}, \pi_{\tau} u_{j}^{n}\right\rangle\right\rangle_{\Gamma_{j k}}=\omega\left(\varepsilon u_{j}^{n}, u_{j}^{n}\right)_{j} \\
& \quad-\omega\left(\mu v_{j}^{n}, v_{j}^{n}\right)_{j}-\left\langle\left\langle\alpha_{I} \pi_{\tau} u_{j}^{n}, \pi_{\tau} u_{j}^{n}\right\rangle\right\rangle_{\Gamma_{j}} . \tag{25}
\end{align*}
$$

We rearrange (22c) so that

$$
\begin{align*}
& \theta_{j k}^{n}=-\theta_{k j}^{n-1}+\beta\left(\pi_{\tau} u_{k}^{n-1}-\pi_{\tau} u_{j}^{n}\right)\left(m_{j k}\right) \\
& \quad j<k \tag{26a}
\end{align*}
$$

$$
\begin{align*}
\theta_{j k}^{n} & =-\theta_{k j}^{n}+\beta\left(\pi_{\tau} u_{k}^{n}-\pi_{\tau} u_{j}^{n}\right)\left(m_{j k}\right), \quad j>k,  \tag{26b}\\
& =\theta_{j k}^{n-1}-\beta\left(\pi_{\tau} u_{j}^{n-1}-\pi_{\tau} u_{k}^{n}\right)\left(m_{j k}\right) \\
& +\beta\left(\pi_{\tau} u_{k}^{n}-\pi_{\tau} u_{j}^{n}\right)\left(m_{j k}\right)  \tag{26c}\\
& =-\theta_{k j}^{n-1}+\beta\left(\pi_{\tau} u_{k}^{n-1}-2 \pi_{\tau} u_{j}^{n-1}\right)\left(m_{j k}\right) \\
& -\beta\left(\pi_{\tau} u_{j}^{n}-2 \pi_{\tau} u_{k}^{n}\right)\left(m_{j k}\right) . \tag{26d}
\end{align*}
$$

Here, (26c) follows from (26b) by applying (26a) and (26d) follows from (26c) by applying (26b) for $n-1$. Motivated by (26a) and (26d), we define the pseudoenergy for the iterative procedure (19):

$$
\begin{align*}
& R^{n}\left(\left\{u^{n}, v^{n}, \theta^{n}\right\}\right)=\sum_{j<k}\left|\theta_{j k}^{n}+\beta \pi_{\tau} u_{j}^{n}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2} \\
& \quad+\sum_{j>k}\left|\theta_{k j}^{n}+\beta \pi_{\tau} u_{k}^{n}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2} . \tag{27}
\end{align*}
$$

Denote $|||f|||=\sum_{j}\langle\langle f, f\rangle\rangle_{\Gamma_{j}}$ for a function defined on $\Gamma$. Then the following recurrence relation for the decay in pseudo-energy holds.

Theorem 4.2. With the pseudo-energy (27), we have the following general recurrence relation which is independent of differential or weak problems:
$R^{n}=R^{n-1}-8 \operatorname{Re} \sum_{j, k}\left\langle\theta_{j k}^{n-1}, \beta \pi_{\tau} u_{j}^{n-1}\right\rangle_{\Gamma_{j k}}$

Proof. Since by (26a), we have

$$
\begin{aligned}
R^{n} & =2 \sum_{j<k}\left|\theta_{j k}^{n}+\beta \pi_{\tau} u_{j}^{n}\right|_{0, \Gamma_{j k}} \\
& =2 \sum_{j<k}\left|\theta_{k j}^{n-1}-\beta \pi_{\tau} u_{k}^{n-1}\right|_{0, \Gamma_{j k}} \quad \text { by (26a) } \\
& =2 \sum_{j>k}\left|\theta_{j k}^{n-1}-\beta \pi_{\tau} u_{j}^{n-1}\right|_{0, \Gamma_{j k}} \\
& =2 \sum_{j>k}\left|-\theta_{k j}^{n-1}+\beta\left(\pi_{\tau} u_{k}^{n-1}-2 \pi_{\tau} u_{j}^{n-1}\right)\right|_{0, \Gamma_{j k}} \quad \text { by (26a) } \\
& =2 \sum_{j<k}\left|\theta_{j k}^{n-1}-\beta\left(\pi_{\tau} u_{j}^{n-1}-2 \pi_{\tau} u_{k}^{n-1}\right)\right|_{0, \Gamma_{j k}} \\
& =2 \sum_{j<k} \mid \theta_{j k}^{n-1}+\beta \pi_{\tau} u_{j}^{n-1} \\
& -\left.2 \beta\left(\pi_{\tau} u_{j}^{n-1}-\pi_{\tau} u_{k}^{n-1}\right)\right|_{0, \Gamma_{j k}} \\
& =R^{n-1}-8 \operatorname{Re} \sum_{j<k}\left\langle\theta_{j k}^{n-1}+\beta \pi_{\tau} u_{j}^{n-1},\right. \\
& \left.\beta\left(\pi_{\tau} u_{j}^{n-1}-\pi_{\tau} u_{k}^{n-1}\right)\right\rangle_{\Gamma_{j k}} \\
& +8 \sum_{j<k}\left|\beta\left(\pi_{\tau} u_{j}^{n-1}-\pi_{\tau} u_{k}^{n-1}\right)\right|_{0, \Gamma_{j k}} \\
& =R^{n-1}-8 \operatorname{Re} \sum_{j<k}\left\langle\theta_{j k}^{n-1}+\beta \pi_{\tau} u_{k}^{n-1},\right. \\
& \left.\beta\left(\pi_{\tau} u_{j}^{n-1}-\pi_{\tau} u_{k}^{n-1}\right)\right\rangle_{\Gamma_{j k}} \\
& =R^{n-1}-8 \operatorname{Re} \sum_{j, k}\left\langle\theta_{j k}^{n-1}, \beta \pi_{\tau} u_{j}^{n-1}\right\rangle_{\Gamma_{j k}}
\end{aligned}
$$

since

$$
\begin{aligned}
\operatorname{Re} & {\left[\sum_{j<k}\left\langle-\beta \pi_{\tau} u_{k}^{n-1}, \beta\left(\pi_{\tau} u_{j}^{n-1}-\pi_{\tau} u_{k}^{n-1}\right)\right\rangle_{\Gamma_{j k}}\right.} \\
& \left.+\sum_{j<k}\left\langle\theta_{j k}^{n-1}, \beta \pi_{\tau} u_{k}^{n-1}\right\rangle_{\Gamma_{j k}}\right] \\
& =\operatorname{Re} \sum_{j<k}\left\langle\beta \pi_{\tau} u_{k}^{n-1},-\beta\left(\pi_{\tau} u_{j}^{n-1}-\pi_{\tau} u_{k}^{n-1}\right)+\theta_{j k}^{n-1}\right\rangle_{\Gamma_{j k}} \\
& =\operatorname{Re} \sum_{j>k}\left\langle\beta \pi_{\tau} u_{j}^{n-1}, \beta\left(\pi_{\tau} u_{j}^{n-1}-\pi_{\tau} u_{k}^{n-1}\right)+\theta_{k j}^{n-1}\right\rangle_{\Gamma_{j k}} \\
& =-\operatorname{Re} \sum_{j>k}\left\langle\beta \pi_{\tau} u_{j}^{n-1}, \theta_{j k}^{n-1}\right\rangle_{\Gamma_{j k}} \quad \text { by }(26 a) \\
& =-\operatorname{Re} \sum_{j>k}\left\langle\theta_{j k}^{n-1}, \beta \pi_{\tau} u_{j}^{n-1}\right\rangle_{\Gamma_{j k}} .
\end{aligned}
$$

This proves (28).

From Theorem 4.2 and (23) it is immediate to have the following theorem.

Theorem 4.3. Let $\beta=\beta_{R}+i \beta_{I}$. Then, with the pseudoenergy (27), the following recurrence relation holds:

$$
\begin{align*}
R^{n} & =R^{n-1}-8 \beta_{R}\left[\left\|\sigma^{1 / 2} u^{n-1}\right\|^{2}+\left\|\mid \alpha_{R}^{1 / 2} \pi_{\tau} u^{n-1}\right\| \|^{2}\right] \\
& +8 \beta_{I}\left[\omega\left\|\mu^{1 / 2} \nu^{n-1}\right\|^{2}-\omega\left\|\varepsilon^{1 / 2} u^{n-1}\right\|^{2}\right. \\
& \left.+\left\|\alpha_{I}^{1 / 2} \pi_{\tau} u^{n-1}\right\| \|^{2}\right] \tag{29}
\end{align*}
$$

Remark 4.4. In the Jacobi-type iterative procedure given in Douglas Jr., Santos, and Sheen (2000) the pseudo-energy is given in the form
$R_{J}^{n}\left(\left\{u^{n}, v^{n}, \theta^{n}\right\}\right)=\sum_{j, k}\left|\theta_{j k}^{n}+\beta \pi_{\tau} u_{j}^{n}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2}$
instead of (27), the recurrence relation is given in the form:

$$
\begin{aligned}
R_{J}^{n} & =R_{J}^{n-1}-4 \beta_{R}\left[\left\|\sigma^{1 / 2} u^{n-1}\right\|^{2}+\left\|\alpha_{R}^{1 / 2} \pi_{\tau} u^{n-1}\right\|^{2}\right] \\
& +4 \beta_{I}\left[\omega\left\|\mu^{1 / 2} v^{n-1}\right\|^{2}-\omega\left\|\varepsilon^{1 / 2} u^{n-1}\right\|^{2}\right. \\
& \left.+\| \| \alpha_{I}^{1 / 2} \pi_{\tau} u^{n-1}\| \|^{2}\right]
\end{aligned}
$$

Thus the energy decays in the Seidel-type procedure is roughly as twice fast as that in the Jacobi-type one.

According to Theorem 4.3, we have the following convergence result.

Theorem 4.5. Assume that $\beta$ satisfies the following conditions:

$$
\begin{align*}
& \beta_{R} \sigma_{*}+\beta_{I} \omega \varepsilon^{*}>0, \beta_{R} \alpha_{R_{*}}-\beta_{I} \alpha_{I_{*}} \geq 0 \\
& \quad \beta_{R}>0, \beta_{I} \leq 0 \tag{30}
\end{align*}
$$

Then, the iteration error $\left\{u^{n}, v^{n}, \theta^{n}\right\}$ satisfying (22) converges to zero as $n$ tends to $\infty$, and $R^{n}$ tends to zero as well.

Proof. Due to (29), we have

$$
\begin{aligned}
R^{n} & =R^{0}-8 \sum_{k=1}^{n-1} \sum_{j}\left[\left(\beta_{R} \sigma+\beta_{I} \omega \varepsilon\right)\left(\sigma u_{j}^{k}, u_{j}^{k}\right)_{j}\right. \\
& -\beta_{I} \omega \mu\left(v_{j}^{k}, v_{j}^{k}\right)_{j} \\
& \left.+\left(\beta_{R} \alpha_{R}-\beta_{I} \alpha_{I}\right)\left\langle\left\langle\pi_{\tau} u_{j}^{k}, \pi_{\tau} u_{j}^{k}\right\rangle\right\rangle_{\Gamma_{j}}\right]
\end{aligned}
$$

Thus, under Condition (30), we see that $R^{n}$ is a nonnegative nonincreasing sequence, and therefore,

$$
\begin{aligned}
& \sum_{k=1}^{\infty}\left[\beta_{R} \sum_{j}\left\{\left(\sigma u_{j}^{k}, u_{j}^{k}\right)_{j}+\left\langle\left\langle\alpha_{R} \pi_{\tau} u_{j}^{k}, \pi_{\tau} u_{j}^{k}\right\rangle\right\rangle_{\Gamma_{j}}\right\}\right. \\
& \quad-\beta_{I} \sum_{j}\left\{\omega\left(\mu v_{j}^{k}, v_{j}^{k}\right)_{j}-\omega\left(\varepsilon u_{j}^{k}, u_{j}^{k}\right)_{j}\right. \\
& \left.\left.\quad+\left\langle\left\langle\alpha_{I} \pi_{\tau} u_{j}^{k}, \pi_{\tau} u_{j}^{k}\right\rangle\right\rangle_{\Gamma_{j}}\right\}\right]<\infty .
\end{aligned}
$$

Suppose that (30) holds. Then we have that $u_{j}^{n} \rightarrow 0$ in $\left[L^{2}\left(\Omega_{j}\right)\right]^{3}$ as $n$ tends to $\infty$, which implies that $\nabla \times u_{j}^{n} \rightarrow 0$ in $\left[H^{-1}\left(\Omega_{j}\right)\right]^{3}$ as $n$ tends to $\infty$. Due to finite dimensionality of $\mathbf{V}_{j}^{h}, \nabla \times u_{j}^{n} \rightarrow 0$ in $\left[L^{2}\left(\Omega_{j}\right)\right]^{3}$ as $n$ tends to $\infty$. Hence, (22b) leads to that $v_{j}^{n} \rightarrow 0$ in $\left[L^{2}\left(\Omega_{j}\right)\right]^{3}$ as $n$ tends to $\infty$. Next, as $n$ tends to $\infty$, (22a) tends to $-\sum_{k}\left\langle\left\langle\theta_{j k}^{n}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j k}}=0$ for all $\varphi \in \mathbf{V}_{j}^{h}$, in particular for $\varphi=\tilde{\varphi}_{j k}$ such that $\pi_{\tau} \tilde{\varphi}\left(m_{j k}\right)=\theta_{j k}^{n}$ and vanishes at the other five centroids of $\partial \Omega_{j}$. This implies that $\theta_{j k}^{n}$ converges to zero as $n$ tends to $\infty$ for all $\{j, k\} \in I^{h}$. The convergence of $R^{n}$ to zero follows from the formula (27). The proof is complete.

### 4.3 An estimate for the spectral radius of the iterative procedure

Let $T_{F, G}=T_{F, G}(\beta, \cdot)=: \mathcal{N} C_{-1}^{h} \times W^{h} \times \Lambda^{h} \rightarrow$ $\mathcal{N} C_{-1}^{h} \times W^{h} \times \Lambda^{h}$ be the affine map such that, for any $(U, V, \theta) \in \mathcal{N}\left(C_{-1}^{h} \times W^{h} \times \Lambda^{h}, \quad\{E, H, \lambda\}=\right.$ $T_{F, G}\{U, V, \theta\}=T_{F, G}(\beta ;\{U, V, \theta\})$ is the solution of the following problem: for $j=1,2, \cdots$,

$$
\begin{align*}
& \left((i \omega \varepsilon+\sigma) E_{j}, \varphi\right)_{j}-\left(H_{j}, \boldsymbol{\nabla} \times \varphi\right)_{j}-\sum_{j}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j k}}  \tag{34}\\
& \quad+\left\langle\left\langle\alpha \pi_{\tau} E_{j}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j}}=\left(F_{j}, \varphi\right)_{j}, \quad \varphi \in \mathbf{V}_{j}^{h}, \tag{31a}
\end{align*}
$$

$i \omega\left(\mu H_{j}, \psi\right)_{j}+\left(\nabla \times E_{j}, \psi\right)_{j}=\left(G_{j}, \psi\right)_{j}, \psi \in W_{j}^{h}$,
$\lambda_{j k}= \begin{cases}-\theta_{k j}+\beta_{j k}\left(\pi_{\tau} U_{k}-\pi_{\tau} E_{j}\right)\left(m_{j k}\right), & j<k, \\ -\lambda_{k j}+\beta_{j k}\left(\pi_{\tau} E_{k}-\pi_{\tau} E_{j}\right)\left(m_{j k}\right), & j>k .\end{cases}$

$$
\begin{align*}
& R\left(T_{0,0}\{E, H, \lambda\}\right)=R(\{E, H, \lambda\}) \\
& \quad-8\left[\beta_{R} \operatorname{Re} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}\right.  \tag{31b}\\
& \left.\quad+\beta_{I} \operatorname{Im} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}\right] . \tag{35}
\end{align*}
$$

On the other hand, Theorem 4.3, incorporated with (24) and (25), gives the equality:

Thus, a combination of (34) and (35) gives

$$
\begin{align*}
& |\gamma|^{2}=1-\frac{8}{R(\{E, H, \lambda\})}\left[\beta_{R} \operatorname{Re} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}\right. \\
& \left.\quad+\beta_{I} \operatorname{Im} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}\right] . \tag{36}
\end{align*}
$$

Obviously, $|\gamma|<1$ if $\beta$ satisfies

$$
\begin{align*}
& \beta_{R} \operatorname{Re} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}} \\
& \quad+\beta_{I} \operatorname{Im} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}>0 . \tag{37}
\end{align*}
$$

The next step is to derive a bound for $R(\{E, H, \lambda\})$, hopefully, in certain norms of $E, H$, and $\lambda$. For this, we observe that, since $\{E, H, \lambda\}$ is an eigenvector of $T_{0,0}$ with the associate eigenvalue $\gamma$, (31) with $F=G=0$ leads to

$$
\begin{align*}
& \left((i \omega \varepsilon+\sigma) E_{j}, \varphi\right)_{j}-\left(H_{j}, \boldsymbol{\nabla} \times \varphi\right)_{j}-\sum_{k}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j k}} \\
& \quad+\left\langle\left\langle\alpha \pi_{\tau} E_{j}, \pi_{\tau} \varphi\right\rangle\right\rangle_{\Gamma_{j}}=0, \quad \varphi \in \mathbf{V}_{j}^{h},  \tag{38a}\\
& i \omega\left(\mu H_{j}, \psi\right)_{j}+\left(\boldsymbol{\nabla} \times E_{j}, \psi\right)_{j}=0, \quad \psi \in W_{j}^{h},  \tag{38b}\\
& \gamma \lambda_{j k}=-\lambda_{k j}+\beta\left(\pi_{\tau} E_{k}-\gamma \pi_{\tau} E_{j}\right)\left(m_{j k}\right), \quad j<k,  \tag{42a}\\
& \lambda_{j k}=-\lambda_{k j}+\beta\left(\pi_{\tau} E_{k}-\pi_{\tau} E_{j}\right)\left(m_{j k}\right), \quad j>k \tag{38d}
\end{align*}
$$

Due to (27) and (38),

$$
\begin{align*}
& R(\{E, H, \lambda\})=\sum_{j<k}\left|\frac{1}{\gamma}\left(\lambda_{k j}-\beta \pi_{\tau} E_{k}\left(m_{j k}\right)\right)\right|_{0, \Gamma_{j k}}^{2}  \tag{42b}\\
& \quad+\sum_{j>k}\left|\lambda_{k j}+\beta \pi_{\tau} E_{k}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2} \\
& \quad=\sum_{j<k}\left|\lambda_{j k}+\beta \pi_{\tau} E_{j}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2} \\
& \quad+\sum_{j>k}\left|\frac{1}{\gamma}\left(\lambda_{j k}-\beta \pi_{\tau} E_{j}\left(m_{j k}\right)\right)\right|_{0, \Gamma_{j k}}^{2} \\
& \quad \leq \frac{1}{|\gamma|^{2}}\left\{\sum_{j<k}\left|\lambda_{j k}+\beta \pi_{\tau} E_{j}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2}\right. \\
& \left.\quad+\sum_{j>k}\left|\lambda_{j k}-\beta \pi_{\tau} E_{j}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2}\right\} \\
& \quad \leq \frac{2}{|\gamma|^{2}} \sum_{j, k \in I^{h}}\left\{\left|\lambda_{j k}\right|_{0, \Gamma_{j k}}^{2}\right. \\
& \left.\quad+|\beta|^{2}\left|\pi_{\tau} E_{j}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2}\right\}
\end{align*}
$$

Observe that $\operatorname{Re} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}>0$ and that (37) is equivalent to $M(\beta)>0$.
The maximum $M^{*}$ of $M(\beta)$ is taken where $\beta$ is chosen optimally such that
$|\boldsymbol{\beta}|^{2}=\frac{\sum_{j, k \in I^{h}}\left|\lambda_{j k}\right|_{0, \Gamma_{j k}}^{2}}{\sum_{j, k \in I^{h}} \mid \pi_{\tau} E_{j}\left(\left.m_{j k}\right|_{0, \Gamma_{j k}} ^{2}\right.}$
and
$\beta_{I} \operatorname{Re} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}=\beta_{R} \operatorname{Im} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}$.

In this case,
$M^{*}=\frac{\left|\sum_{j, k \in l^{l}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}\right|}{2\left[\sum_{j, k \in I^{n}}\left|\lambda_{j k}\right|_{0, \Gamma_{j k}}^{2}\right]^{1 / 2}\left[\sum_{j, k \in I^{h}}\left|\pi_{\tau} E_{j}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2}\right]^{1 / 2}}$,
and therefore, $|\gamma|$ satisfies
$|\gamma|^{2} \leq \frac{1}{1+4 M^{*}}$.
We summarize the above result in the following theorem, which suggests a method of the choice of the optimal $\beta$.

Theorem 4.7. Let $\{E, H, \lambda\}$ be an eigenvector of $T_{0,0}$ with the associate eigenvalue $\gamma$. Assume that (37) holds.

Then $|\gamma|<1$. Moreover, $\beta$, with positive real part, is optimal if it satisfies (42); in this case, the following bound for $|\gamma|$ holds:

$$
\begin{align*}
|\gamma|^{2} & \leq 1 /\left\{1+2\left|\sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}\right|\right. \\
& {\left.\left[\sum_{j, k \in I^{h}}\left|\lambda_{j k}\right|_{0, \Gamma_{j k}}^{2} \sum_{j, k \in I^{h}}\left|\pi_{\tau} E_{j}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2}\right]^{-1 / 2}\right\} . } \tag{43}
\end{align*}
$$

Remark 4.8. Notice that Conditions (37) and (42b) are equivalent to

$$
\begin{aligned}
& \beta_{R}\left\{\left\|\sigma^{1 / 2} E\right\|^{2}+\| \| \alpha_{R}^{1 / 2} \pi_{\tau} E \mid \|^{2}\right\} \\
& \quad-\beta_{I}\left\{\omega\left\|\mu^{1 / 2} H\right\|^{2}-\omega\left\|\varepsilon^{1 / 2} E\right\|^{2}+\mid\left\|\alpha_{I}^{1 / 2} \pi_{\tau} E\right\| \|^{2}\right\}>0
\end{aligned}
$$

and

$$
\begin{aligned}
& \beta_{I}\left\{\left\|\sigma^{1 / 2} E\right\|^{2}+\| \| \alpha_{R}^{1 / 2} \pi_{\tau} E \|^{2}\right\}=-\beta_{R}\left\{\omega\left\|\mu^{1 / 2} H\right\|^{2}\right. \\
& \left.\quad-\omega\left\|\varepsilon^{1 / 2} E\right\|^{2}+\left\|\alpha_{I}^{1 / 2} \pi_{\tau} E\right\| \|^{2}\right\}
\end{aligned}
$$

respectively.
Remark 4.9. We remark that as $\left|\sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}\right|$ vanishes, so do $E$ and $H$, and $|\gamma|$ tends to 1 .

Remark 4.10. With minor modifications in the above proof of Theorem 4.7, $\beta$ 's can be chosen such that they vary according to $j$ and $k$. Also, practically it suffices to update $\beta$ only once after certain number of iterations, instead of updating $\beta$ at every iteration. However, the true solution $\left\{\tilde{E}^{h}, \tilde{H}^{h}, \tilde{\lambda}^{h}\right\}$ is not known in the actual computations, and hence the use of optimal choice of $\beta$ as in Theorem 4.7, or its approximation, still needs further investigations.
Remark 4.11. It should be observed that the optimal $\beta$ in Theorem 4.7 depends on the values of $\pi_{\tau} E$ and $\lambda$ on the interfaces $\Gamma_{j k}$ 's. Theorem 4.12 establishes the spectral radius estimate and choice of $\beta$ to guarantee an optimal order of convergence. In actual computations, however, Theorem 4.7 suggests a much more efficient choice of $\beta$ than this, although it is heuristic: for this, see the next section on numerical results.

Instead of varying $\beta$ as in Theorem 4.7, consider the case in which $\beta$ is independent of solutions. For the following theorem, set

$$
\begin{aligned}
& \zeta=\max _{j}\left\{\frac{h_{\max \left(\Omega_{j}\right)}}{h_{\min \left(\Omega_{j}\right)}}\right\}, \\
& h_{\max }=\max _{j}\left\{h_{\max }\left(\Omega_{j}\right)\right\}, \quad h_{\min }=\min _{j}\left\{h_{\min }\left(\Omega_{j}\right)\right\},
\end{aligned}
$$

where $h_{\min }\left(\Omega_{j}\right)$ and $h_{\max }\left(\Omega_{j}\right)$ are the minimum and maximum edge length of $\Omega_{j}$, respectively. The following theorem gives an asymptotic estimate for the spectral radius of $T_{0,0}$.

Theorem 4.12. Let $\rho\left(T_{0,0}\right)$ be the spectral radius of $T_{0,0}$. Suppose $\mu_{*}>0$. If $\beta$ is chosen such that $\beta_{R} \alpha_{R_{*}}-\beta_{I} \alpha_{I_{*}}>$ 0 and $\beta_{R} \sigma_{*}+\beta_{I} \omega \varepsilon^{*}>0$ with $\beta_{R}>0, \beta_{I} \leq 0$, then
$\left|\rho\left(T_{0,0}\right)\right| \approx 1-C h_{\min }$.
Proof. We begin with finding a lower bound for $M(\beta)$ given by (41), which will give an upper bound for $|\gamma|$ by (40). For this, let $\tilde{\Omega}_{j}$ be an arbitrary element among $\Omega_{j}$ 's and $\tilde{\Gamma}_{j k}$ an interior face of $\tilde{\Omega}_{j}$ common with another element $\Omega_{k}$. Choose $\tilde{\varphi}$ in (38) such that
$\pi_{\tau} \tilde{\varphi}\left(m_{j k}\right)=\left\{\begin{array}{lll}\tilde{\lambda}_{j k} & \text { on } & \tilde{\Gamma}_{j k}, \\ 0 & \text { on } & \partial \Omega_{j} \backslash \tilde{\Gamma}_{j k} .\end{array}\right.$
Recall the estimates given in Douglas Jr., Santos, and Sheen (2000):
$\|\tilde{\varphi}\|_{0, \Omega_{j}}^{2} \leq C_{1}\left\langle\left\langle\tilde{\lambda}_{j k}, \tilde{\lambda}_{j k}\right\rangle\right\rangle_{\Gamma_{j k}}$,
$\|\nabla \times \tilde{\varphi}\|_{0, \Omega_{j}}^{2} \leq C_{2}\left\langle\left\langle\tilde{\lambda}_{j k}, \tilde{\lambda}_{j k}\right\rangle\right\rangle_{\Gamma_{j k}}$,
$\left\langle\left\langle\pi_{\tau} E_{j}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}} \leq C_{3}\left\|E_{j}\right\|_{0, \Omega_{j}}^{2}$,
where

$$
\begin{equation*}
C_{1}=C \max _{j} \frac{h_{\max }\left(\Omega_{j}\right)^{2}}{h_{\min }\left(\Omega_{j}\right)}, \quad C_{2}=C \max _{j} \frac{h_{\max }\left(\Omega_{j}\right)^{2}}{h_{\min }\left(\Omega_{j}\right)^{3}} \tag{46}
\end{equation*}
$$

$C_{3}=C \max _{j} \frac{h_{\max }\left(\Omega_{j}\right)}{h_{\min }\left(\Omega_{j}\right)^{2}}$.
It follows from (38a) and (45) that

$$
\begin{aligned}
& \left\langle\left\langle\tilde{\lambda}_{j k}, \tilde{\lambda}_{j k}\right\rangle\right\rangle_{\Gamma_{j k}}=\left((i \omega \varepsilon+\sigma) E_{j}, \tilde{\varphi}\right)_{j}-\left(H_{j}, \nabla \times \tilde{\varphi}\right)_{j} \\
& \quad \leq\left[\sqrt{C_{1}}\left(\omega\left\|\varepsilon E_{j}\right\|_{0, \Omega_{j}}+\left\|\sigma E_{j}\right\|_{0, \Omega_{j}}\right)\right. \\
& \left.\left.\quad+\sqrt{C_{2}}\left\|H_{j}\right\|_{0, \Omega_{j}}\right]\left\langle\left\langle\tilde{\lambda}_{j k}, \tilde{\lambda}_{j k}\right\rangle\right\rangle\right\rangle_{\Gamma_{j k}}^{1 / 2},
\end{aligned}
$$

from which it is immediate that

$$
\begin{aligned}
& \left\langle\left\langle\lambda_{j k}, \lambda_{j k}\right\rangle\right\rangle_{\Gamma_{j k}} \leq C_{1}\left(\omega^{2}\left\|\varepsilon E_{j}\right\|_{0, \Omega_{j}}^{2}+\left\|\sigma E_{j}\right\|_{0, \Omega_{j}}^{2}\right) \\
& \quad+C_{2}\left\|H_{j}\right\|_{0, \Omega_{j}}^{2}
\end{aligned}
$$

Assume that $\beta_{R} \sigma_{*}+\beta_{I} \omega \varepsilon^{*}>0$ and $\mu_{*}>0$. Choose $\beta_{R}>$ 0 and $\beta_{I} \leq 0$. Then, a combination of (41), (24), (25), $(45 \mathrm{c}),(47)$, and (46) results in

$$
\begin{array}{ll}
M(\beta)=\left[\beta_{R} \operatorname{Re} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}\right. & \begin{array}{l}
\text { The maximum of the last quantity is taken if we choose } \\
\\
\beta \text { such that }
\end{array} \\
\left.\quad+\beta_{I} \operatorname{Im} \sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}\right] & |\beta|=\sqrt{h_{\max } h_{\min }\left(\omega^{2} \varepsilon^{* 2}+\sigma^{* 2}\right)+\zeta}  \tag{47}\\
/\left[\sum_{j, k \in I^{h}}\left\{\left|\lambda_{j k}\right|_{0, \Gamma_{j k}}^{2}+|\beta|^{2}\left|\pi_{\tau} E_{j}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2}\right\}\right] & \beta_{R}\left(\sigma_{*}+\alpha_{R} \frac{h_{\min }}{C h_{\max }^{2}}\right)=-\beta_{I}\left(\omega \varepsilon^{*}+\omega \mu_{*}-\alpha_{I} \frac{h_{\min }}{C h_{\max }^{2}}\right)
\end{array}
$$

$$
=\left[\beta_{R}\left(\sigma\|E\|^{2}+\alpha_{R}\| \| \pi_{\tau} E\| \|^{2}\right)\right.
$$

$$
\left.+\beta_{I}\left(\omega \varepsilon\|E\|^{2}-\omega \mu\|H\|^{2}-\alpha_{I}\| \| \pi_{\tau} E\| \|^{2}\right)\right]
$$

$$
/\left[\sum_{j, k \in I^{h}}\left\{\left|\lambda_{j k}\right|_{0, \Gamma_{j k}}^{2}+|\beta|^{2}\left|\pi_{\tau} E_{j}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2}\right\}\right]
$$

$$
=\left[\left(\beta_{R} \sigma+\beta_{I} \omega \varepsilon\right)\|E\|^{2}-\beta_{I} \omega \mu\|H\|^{2}\right.
$$

$$
\left.+\left(\beta_{R} \alpha_{R}+\beta_{I} \alpha_{I}\right)\| \| \pi_{\tau} E \mid \|^{2}\right]
$$

$$
/\left[\sum_{j, k \in I^{h}}\left\{\left|\lambda_{j k}\right|_{0, \Gamma_{j k}}^{2}+|\beta|^{2}\left|\pi_{\tau} E_{j}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2}\right\}\right]
$$

$$
\geq\left[\left(\beta_{R} \sigma_{*}+\beta_{I} \omega \varepsilon^{*}\right)\|E\|^{2}-\beta_{I} \omega \mu_{*}\|H\|^{2}\right.
$$

$$
\left.+\left(\beta_{R} \alpha_{R}-\beta_{I} \alpha_{I}\right) \frac{h_{\min }}{C h_{\max }^{2}}\|E\|^{2}\right]
$$

$$
/\left[C_{1}\left(\omega^{2} \varepsilon^{* 2}+\sigma^{* 2}\right)\|E\|^{2}+C_{2}\|H\|^{2}\right.
$$

$$
\left.+|\beta|^{2} C_{3}\|E\|^{2}\right]
$$

$$
=\left[\left(\beta_{R} \sigma_{*}+\beta_{I} \omega \varepsilon^{*}+\left(\beta_{R} \alpha_{R}-\beta_{I} \alpha_{I}\right) \frac{h_{\min }}{C h_{\max }^{2}}\right) \times\right.
$$

$$
\left.\times\|E\|^{2}-\beta_{I} \omega \mu_{*}\|H\|^{2}\right]
$$

$$
/\left[C_{1}\left(\omega^{2} \varepsilon^{* 2}+\sigma^{* 2}\right)\|E\|^{2}+C_{2}\|H\|^{2}\right.
$$

$$
\left.+|\beta|^{2} C_{3}\|E\|^{2}\right]
$$

$$
\begin{aligned}
& \left.\left.\times \frac{h_{\min }}{C h_{\max }^{2}}-\beta_{I} \omega \mu_{*}\right\}\left\{\|E\|^{2}+\|H\|^{2}\right\}\right] \\
& /\left[\left[C_{1} C_{3}^{-1}\left(\omega^{2} \varepsilon^{* 2}+\sigma^{* 2}\right)+C_{2} C_{3}^{-1}+|\beta|^{2}\right]\right. \\
& \left.\left\{\|E\|^{2}+\|H\|^{2}\right\}\right] \\
& \geq \frac{h_{\min }}{C \zeta}\left[\operatorname { m i n } \left\{\beta_{R} \sigma_{*}+\beta_{I} \omega \varepsilon^{*}+\left(\beta_{R} \alpha_{R}-\beta_{I} \alpha_{I}\right) \times\right.\right. \\
& \left.\left.\times \frac{h_{\min }}{C h_{\max }^{2}}-\beta_{I} \omega \mu_{*}\right\}\right] \\
& /\left[\left[C_{1} C_{3}^{-1}\left(\omega^{2} \varepsilon^{* 2}+\sigma^{* 2}\right)+C_{2} C_{3}^{-1}+|\beta|^{2}\right]\right]
\end{aligned}
$$

With this $\beta$, we have that $M(\beta) \geq \frac{h_{\min }}{C}$, and thus $|\gamma| \approx$ $1-C h_{\text {min }}$.
This completes the proof.

Remark 4.13. Also observe that the results corresponding to Theorem 4.7 and Theorem 4.12 for the Jacobi-type procedure are immediate. Replacing only (40) by

$$
|\gamma|^{2} \leq \frac{1}{1+4 M(\beta)}
$$

instead of (43), one has

$$
\begin{aligned}
& |\gamma|^{2} \leq 1 /\left[1+\left|\sum_{j, k \in I^{h}}\left\langle\left\langle\lambda_{j k}, \pi_{\tau} E_{j}\right\rangle\right\rangle_{\Gamma_{j k}}\right|\right. \\
& \left.\quad\left\{\sum_{j, k \in I^{h}}\left|\lambda_{j k}\right|_{0, \Gamma_{j k}}^{2} \sum_{j, k \in I^{h}}\left|\pi_{\tau} E_{j}\left(m_{j k}\right)\right|_{0, \Gamma_{j k}}^{2}\right\}^{-1 / 2}\right]
\end{aligned}
$$

and then the estimate for the spectral radius will be accordingly affected. Indeed instead of (44),
$|\gamma| \approx 1-C h_{\min }$
with the same constant $C$ as in (44).

$$
\geq \frac{1}{C_{3}}\left[\operatorname { m i n } \left\{\beta_{R} \sigma_{*}+\beta_{I} \omega \varepsilon^{*}+\left(\beta_{R} \alpha_{R}-\beta_{I} \alpha_{I}\right) \times\right.\right.
$$

### 4.4 The red-black Seidel-type procedure

In order to supplement the lack in parallelism of Seideltype iterative procedure, a red-black Seidel-type procedure is proposed in what follows. Let the set of subdomain indices be divided into the two parts $I_{R}$ and $I_{B}$, so that
$\bar{\Omega}=\left[\cup_{j \in I_{R}} \bar{\Omega}_{j}\right] \bigcup\left[\cup_{j \in I_{B}} \bar{\Omega}_{j}\right], \quad \Omega_{j} \cap_{j \neq k} \Omega_{k}=\emptyset$,
and each element $\Omega_{j}, j \in I_{R}$, is not adjacent to any element $\Omega_{k}, k \in I_{R}$. The red-black Seidel-type iterative procedure corresponding to Problem (19) is based on the the following updating scheme for $\lambda_{j k}^{h, n} \forall k \forall j$ :
$\lambda_{j k}^{h, n}=\left\{\begin{array}{c}-\lambda_{k j}^{h, n-1}+\beta_{j k}\left(\pi_{\tau} E_{k}^{h, n-1}-\pi_{\tau} E_{j}^{h, n}\right)\left(m_{j k}\right) \\ \forall k, \forall j \in I_{R}, \\ -\lambda_{k j}^{h, n}+\beta_{j k}\left(\pi_{\tau} E_{k}^{h, n}-\pi_{\tau} E_{j}^{h, n}\right)\left(m_{j k}\right) \\ \forall k, \forall j \in I_{B} .\end{array}\right.$

Then the red-black Seidel-type iterative procedure consists of the alteration of the two substeps: the first substep of updating (48a) and solving (19a) and (19b) for all $j \in I_{R}$, and the second substep of updating (48b) and solving (19a) and (19b) for all $j \in I_{B}$.
Introduce the pseudo-energy for red-black Seidel-type iterative procedure in the form

$$
\begin{aligned}
& R^{n}\left(\left\{u^{n}, v^{n}, \theta^{n}\right\}\right)=\sum_{j \in I_{R}} \sum_{k}\left|\theta_{j k}^{n}+\beta \pi_{\tau} u_{j}^{n}\left(m_{j k}\right)\right|^{2} \\
& \quad+\sum_{j \in I_{B}} \sum_{k}\left|\theta_{k j}^{n}+\beta \pi_{\tau} u_{k}^{n}\left(m_{j k}\right)\right|^{2} .
\end{aligned}
$$

With this pseudo-energy, all the results in the previous subsections for Seidel-type iterative procedure remain valid for red-black Seidel-type iterative procedure.

## 5 Numerical Results

Let the computational domain $\Omega$ be given as $[0,4 \times$ $\left.10^{3} \mathrm{~m}\right]^{3}$. For the physical parameters in our numerical simulation we use the data following Coggon (1971). Let the frequency $\omega$ be 10 Hz , and the electric permittivity and magnetic permeability contrasts $\varepsilon / \varepsilon_{0}, \mu / \mu_{0}$ be 10 and 1 , where $\varepsilon_{0}$ is the electric permittivity of the vacuum and $\mu_{0}$ the magnetic permeability of the free space. Let the


Figure 1: The logarithmic relative $L^{2}$-error with the fixed $\beta$ value
conductivity $\sigma$ be $0.01 \mathrm{~S} / \mathrm{m}$ (Siemens $/$ meter). The coefficient matrix $\alpha$ in the absorbing boundary condition (3c) is given by the diagonal matrix $(1-i) \sqrt{\frac{\sigma}{2 \omega \mu}} I$. In the rest of the section some numerical results are illustrated to compare the classical Gauss-Jacobi type algorithm with the presented Gauss-Seidel type and red-black type ones with various updating procedures.

### 5.1 The case of homogeneous problems

Setting $F=G=0$, we investigate in the convergence of the solutions of (3) to zeros. Denote by $\left\{u^{h, n}, v^{h, n}, \theta^{h, n}\right\}$ the numerical solution at the $n$th step error equations (22), and we calculate
$\log _{10}\left[\frac{\left\|u^{h, n}\right\|_{0}+\left\|v^{h, n}\right\|_{0}+\left\|\boldsymbol{\nabla} \times u^{h, n}\right\|_{0, h}}{\left\|u^{h, 1}\right\|_{0}+\left\|v^{h, 1}\right\|_{0}+\left\|\boldsymbol{\nabla} \times u^{h, 1}\right\|_{0, h}}\right]$
as logarithmic relative $L^{2}$-errors. Here, $u^{h, n}, \nu^{h, n}$ are the solutions at the $n$th iterative step. For the numerical computations in this subsection, the domain $\Omega$ is divided into the $80 \times 80 \times 80$ congruent subdomains.
We consider the case of updating procedure with a fixed $\beta$ satisfying (47). Figure 1 shows the comparison of the Gauss-Jacobi, Gauss-Seidel, and red-black type iterative procedures. For all cases the speed of convergence of the Gauss-Seidel and Red-Black type procedures are faster (approximately by a factor of two) than that of the GaussJacobi type one.

### 5.2 The case of nonhomogeneous problems

Let us take $F$ and $G$ such that the exact solutions of (3) are
$E_{e x}=(\phi(x, y, z), 0,0)$
$H_{e x}=2000 \times \alpha\left(0, \frac{\partial \phi(x, y, z)}{\partial z},-\frac{\partial \phi(x, y, z)}{\partial y}\right)$
where
$\phi(x, y, z)=\left(\frac{x}{2000}+1\right)\left(\frac{x}{2000}-1\right) \frac{y}{2000}\left(\frac{z}{2000}-1\right)$.
As a synthetic model for nonhomogeneous case, we test the model in COMMEMI project Zhdanov, Varentsov, Weaver, Golubev, and Krylov (1997). It consists of a conductive block of $2 \mathrm{~S} / \mathrm{m}$ embedded in a homogeneous body with conductivity $0.01 \mathrm{~S} / \mathrm{m}$. The anomaly size is $1 \mathrm{~km} \times 1 \mathrm{~km} \times 2 \mathrm{~km}$ in the computational domain $4 \mathrm{~km} \times 4 \mathrm{~km} \times 4 \mathrm{~km}$.


Figure 2: Two-dimensional slice of synthetic anomaly model at $y=0$

We define the relative $L^{2}$-error $\mathcal{E}_{\text {rel }}^{h, n}$ by the quotient
$\frac{\left\|E_{e x}-E^{h, n}\right\|_{0}+\left\|H_{e x}-H^{h, n}\right\|_{0}+\left\|\nabla \times\left(E_{e x}-E^{h, n}\right)\right\|_{0, h}}{\left\|E_{e x}\right\|_{0}+\left\|H_{e x}\right\|_{0}+\left\|\nabla \times E_{e x}\right\|_{0, h}}$,
the $n-1$ st step to the current $n$th step is less than 0.001 ; that is,
$1-\frac{\mathcal{E}_{\text {rel }}^{h, n}}{\mathcal{E}_{\text {rel }}^{h, n-1}} \leq 0.001$.
Although the theoretical convegence given in Theorem 3.1 is suboptimal, we remark that an optimal convergence rate of $O(h)$ can be observed in Table 1.

Table 1 : Relative $L^{2}$ errors and reduction ratio $r$

| $\mathrm{h}(\mathrm{km})$ | Num. of Grids | Rel. $L^{2}$-Error | Red. Ratio $r$ |
| :--- | :---: | :---: | :---: |
| $1 / 4$ | $16 \times 16 \times 16$ | 0.084063 |  |
| $1 / 8$ | $32 \times 32 \times 32$ | 0.042073 | 0.998 |
| $1 / 16$ | $64 \times 64 \times 64$ | 0.021045 | 0.999 |
| $1 / 32$ | $128 \times 128 \times 128$ | 0.010546 | 0.996 |

Figures 3 and 4 plot horizontal slices of the real part of the electric field $E_{1}^{h, N}$ across the source location for the homogeneous and nonhomogeneous cases, respectively. For these figures the stopping was made when

$$
\begin{aligned}
& {\left[\left\|E^{h, n+1}-E^{h, n}\right\|_{0}+\left\|H^{h, n+1}-H^{h, n}\right\|_{0}\right.} \\
& \left.\quad+\left\|\boldsymbol{\nabla} \times\left(E^{h, n+1}-E^{h, n}\right)\right\|_{0, h}\right] \\
& \quad /\left[\left\|E^{h, n}\right\|_{0}+\left\|H^{h, n}\right\|_{0}+\left\|\boldsymbol{\nabla} \times E^{h, n}\right\|_{0, h}\right]
\end{aligned}
$$

was 0.00017 . From the snapshots in Figures 3 and 4 one may observe that the absorbing boundary conditions (3c) perform well.

### 5.3 Calculus of the spectral radius

In this subsection we will compute and confirm numerically the actual rate of convergence of the domain decomposition iteration that is estimated by the spectral radius given in (44) in $\S 4.3$. For this, let $E, E^{h, n}$ and $E^{h, \infty}$ be the exact, $n$-th iterative, and the limiting iterative solutions, respectively. The $n$-th error $e^{h, n}=\left\|E-E^{h, n}\right\|_{L^{2}(\Omega)}$ is then estimated by

$$
\begin{aligned}
e^{h, n} & \leq\left\|E-E^{h, \infty}\right\|_{L^{2}(\Omega)}+\left\|E^{h, \infty}-E^{h, n}\right\|_{L^{2}(\Omega)} \\
& \left.\leq C_{1}\|E\|_{L^{2}(\Omega)}\right)^{\alpha_{1}}+C_{2}\left(1-C_{3} h^{\alpha_{2}}\right)^{n} .
\end{aligned}
$$



Figure 3 : The real part of $e_{1}^{h, N}$ for homogeneous model

Since $\alpha_{1}$ is at least 0.5 by Theorem 3.1 but nearly 1 by the numerical experiment given in $\S 5.2$, we will concentrate on calculating only $C_{3}$ and $\alpha_{2}$ which relate with the iterative procedure. By using the following relations
$e^{h, n}-e^{h, n+1} \approx C_{2}\left(1-C_{3} h^{\alpha_{2}}\right)^{n} C_{3} h^{\alpha_{2}}$
$e^{h, n+1}-e^{h, n+2} \approx C_{2}\left(1-C_{3} h^{\alpha_{2}}\right)^{n+1} C_{3} h^{\alpha_{2}}$,
we have

$$
\begin{aligned}
& \frac{e^{h, n+1}-e^{h, n+2}}{e^{h, n}-e^{h, n+1}} \approx 1-C_{3} h^{\alpha_{2}} \\
& \frac{e^{2 h, n+1}-e^{2 h, n+2}}{e^{2 h, n}-e^{2 h, n+1}} \approx 1-C_{3} 2^{\alpha_{2}} h^{\alpha_{2}} .
\end{aligned}
$$

It seems to be improper to take a fixed number $n$, and hence we obtain the values by taking averages through $n=1$ to a certain number (=50). In Table 2, we present numerical calculus with $\beta$ satisfying (47) for the spectral radius coefficients $C_{3}$ and $\alpha_{2}$ :
$\rho\left(T_{0,0}\right) \leq 1-C_{3} h^{\alpha_{2}}$.


Figure 4 : The real part of $e_{1}^{h, N}$ for nonhomogeneous model

Table 2 : Numerical calculus of coefficients in the spectral radius estimation

|  | Jacobi-type | Red-Black |
| :---: | :---: | :---: |
| $C_{3}$ | 0.537675 | 1.152814 |
| $\alpha_{2}$ | 1.078837 | 1.156271 |

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