

# A PRIORI AND A POSTERIORI ERROR ANALYSIS OF AN AUGMENTED MIXED FINITE ELEMENT METHOD FOR INCOMPRESSIBLE ELASTICITY

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ABSTRACT. In this paper we extend recent results on the a priori and a posteriori error analysis of an augmented mixed finite element method for the linear elasticity problem, to the case of incompressible materials. Similarly as before, the present approach is based on the introduction of the Galerkin least-squares type terms arising from the constitutive and equilibrium equations, and from the relations defining the pressure in terms of the stress tensor and the rotation in terms of the displacement, all them multiplied by stabilization parameters. We show that these parameters can be suitably chosen so that the resulting augmented variational formulation is defined by a strongly coercive bilinear form, whence the associated Galerkin scheme becomes well posed for any choice of finite element subspaces. Next, we derive a reliable and efficient residual-based a posteriori error estimator for the augmented mixed finite element scheme. Finally, several numerical results confirming the theoretical properties of this estimator, and illustrating the capability of the corresponding adaptive algorithm to localize the singularities and the large stress regions of the solution, are also reported.

## 1. INTRODUCTION

The stabilization of dual-mixed variational formulations through the application of diverse procedures has been widely investigated during the last two decades. In particular, the augmented variational formulations, also known as Galerkin least-squares methods, and which go back to [14] and [15], have already been extended in different directions. Some applications to elasticity problems can be found in [17] and [9], and a non-symmetric variant was considered in [13] for the Stokes problem. In addition, stabilized mixed finite element methods for related problems, including Darcy and incompressible flows, can be seen in [2], [6], [16], [20], [21], and [23]. For an abstract framework concerning the stabilization of general mixed finite element methods, we refer to [8].

On the other hand, a new stabilized mixed finite element method for plane linear elasticity with homogeneous Dirichlet boundary conditions was presented and analyzed in [18]. The approach there is based on the introduction of suitable Galerkin least-squares terms arising from the constitutive and equilibrium equations, and from the relation defining the rotation in terms of the displacement. It is shown that the resulting continuous and discrete augmented formulations are well posed, and that the latter becomes locking-free. Moreover, since the augmented variational formulation is strongly coercive, arbitrary finite element subspaces can be utilized in the discrete scheme, which constitutes one of its main advantages. In particular, Raviart-Thomas spaces of lowest order for the stress tensor, piecewise linear elements for the displacement, and piecewise constants for the rotation can be used. The corresponding

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1991 *Mathematics Subject Classification.* 65N15, 65N30, 65N50, 74B05.

*Key words and phrases.* mixed finite element, incompressible flow, a posteriori error estimator.

This research was partially supported by CONICYT, Chile through the FONDAP Program in Applied Mathematics, by the Dirección de Investigación of the Universidad de Concepción through the Advanced Research Groups Program, and by Fundación Andes, Chile, through the project C-14040.

extension to the case of non-homogeneous Dirichlet boundary conditions was provided recently in [19]. In addition, a residual based a posteriori error analysis yielding a reliable and efficient estimator for the augmented method from [18], is provided in the recent work [5]. A posteriori error analyses of the traditional mixed finite element methods for the elasticity problem can be seen in [10] and the references therein.

The purpose of the present paper is to extend the results from [18] and [5] to the case of incompressible elasticity. The rest of this work is organized as follows. In Section 2 we describe the boundary value problem of interest, establish its dual-mixed variational formulation, and prove that it is well-posed. Then, in Sections 3 and 4 we introduce and analyze the continuous and discrete augmented formulations, respectively. Next, in Section 5 we develop the residual-based a posteriori error analysis of our augmented mixed finite element method. Finally, several numerical results confirming the reliability and efficiency of the estimator are provided in Section 6. The capability of the corresponding adaptive algorithm to localize the singularities and the large stress regions of the solution is also illustrated here.

We end this section with some notations to be used below. Given any Hilbert space  $U$ ,  $U^2$  and  $U^{2 \times 2}$  denote, respectively, the space of vectors and square matrices of order 2 with entries in  $U$ . In particular,  $\mathbf{I}$  is the identity matrix of  $\mathbb{R}^{2 \times 2}$ , and given  $\boldsymbol{\tau} := (\tau_{ij})$ ,  $\boldsymbol{\zeta} := (\zeta_{ij}) \in \mathbb{R}^{2 \times 2}$ , we write as usual  $\boldsymbol{\tau}^\mathbf{t} := (\tau_{ji})$ ,  $\text{tr}(\boldsymbol{\tau}) := \sum_{i=1}^2 \tau_{ii}$ ,  $\boldsymbol{\tau}^\mathbf{d} := \boldsymbol{\tau} - \frac{1}{2} \text{tr}(\boldsymbol{\tau}) \mathbf{I}$ , and  $\boldsymbol{\tau} : \boldsymbol{\zeta} := \sum_{i,j=1}^2 \tau_{ij} \zeta_{ij}$ . Also, in what follows we utilize the standard terminology for Sobolev spaces and norms, employ  $\mathbf{0}$  to denote a generic null vector, and use  $C$  and  $c$ , with or without subscripts, bars, tildes or hats, to denote generic constants independent of the discretization parameters, which may take different values at different places.

## 2. THE PROBLEM AND ITS DUAL-MIXED FORMULATION

Let  $\Omega$  be a bounded and simply connected polygonal domain in  $\mathbb{R}^2$  with boundary  $\Gamma$ . Our goal is to determine the displacement  $\mathbf{u}$ , the stress tensor  $\boldsymbol{\sigma}$ , and the pressure-like unknown  $p$  of a linear incompressible material occupying the region  $\Omega$ , under the action of an external force. In other words, given a volume force  $\mathbf{f} \in [L^2(\Omega)]^2$ , we seek a symmetric tensor field  $\boldsymbol{\sigma}$ , a vector field  $\mathbf{u}$  and a scalar field  $p$  such that

$$\begin{aligned} \boldsymbol{\sigma} &= 2\mu \boldsymbol{\varepsilon}(\mathbf{u}) - p\mathbf{I} \quad \text{in } \Omega, & \mathbf{div}(\boldsymbol{\sigma}) &= -\mathbf{f} \quad \text{in } \Omega, \\ \mathbf{div}(\mathbf{u}) &= 0 \quad \text{in } \Omega, & \mathbf{u} &= \mathbf{0} \quad \text{on } \Gamma, \end{aligned} \tag{2.1}$$

where  $\boldsymbol{\varepsilon}(\mathbf{u}) := \frac{1}{2}(\nabla \mathbf{u} + (\nabla \mathbf{u})^\mathbf{t})$  is the linearized strain tensor,  $\mu$  is the shear modulus, and  $\mathbf{div}$  stands for the usual divergence operator  $\text{div}$  acting along each row of the tensor.

Since  $\text{tr}(\boldsymbol{\varepsilon}(\mathbf{u})) = \mathbf{div}(\mathbf{u})$  in  $\Omega$ , we find from the first equation in (2.1) that the incompressibility condition  $\mathbf{div}(\mathbf{u}) = 0$  in  $\Omega$  can be stated in terms of the stress tensor and the pressure as follows

$$p + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}) = 0 \quad \text{in } \Omega. \tag{2.2}$$

Next, we choose to impose weakly the symmetry of  $\boldsymbol{\sigma}$  through the introduction of the infinitesimal rotation tensor  $\boldsymbol{\gamma} := \frac{1}{2}(\nabla \mathbf{u} - (\nabla \mathbf{u})^\mathbf{t})$  as a further unknown (see [1] and [24]), which yields

$$\frac{1}{2\mu}(\boldsymbol{\sigma} + p\mathbf{I}) = \boldsymbol{\varepsilon}(\mathbf{u}) = \nabla \mathbf{u} - \boldsymbol{\gamma} \quad \text{in } \Omega. \tag{2.3}$$

Note that (2.2) and (2.3) imply the modified constitutive equation

$$\frac{1}{2\mu} \boldsymbol{\sigma}^\mathbf{d} = \boldsymbol{\varepsilon}(\mathbf{u}) \quad \text{in } \Omega. \tag{2.4}$$

Then, testing equations (2.3) and (2.2) and weakly taking care of the equilibrium equation of (2.1) and the symmetry of  $\boldsymbol{\sigma}$  gives rise to the problem: Find  $(\boldsymbol{\sigma}, p, (\mathbf{u}, \boldsymbol{\gamma}))$  in  $H(\mathbf{div}; \Omega) \times L^2(\Omega) \times Q$  such that

$$\begin{aligned} \frac{1}{2\mu} \int_{\Omega} \boldsymbol{\sigma} : \boldsymbol{\tau} + \frac{1}{2\mu} \int_{\Omega} p \operatorname{tr}(\boldsymbol{\tau}) + \frac{1}{2\mu} \int_{\Omega} q \operatorname{tr}(\boldsymbol{\sigma}) + \frac{1}{\mu} \int_{\Omega} p q + \int_{\Omega} \mathbf{u} \cdot \mathbf{div}(\boldsymbol{\tau}) + \int_{\Omega} \boldsymbol{\gamma} : \boldsymbol{\tau} &= 0, \\ \int_{\Omega} \mathbf{v} \cdot \mathbf{div}(\boldsymbol{\sigma}) + \int_{\Omega} \boldsymbol{\eta} : \boldsymbol{\sigma} &= - \int_{\Omega} \mathbf{f} \cdot \mathbf{v}, \end{aligned}$$

for all  $(\boldsymbol{\tau}, q, (\mathbf{v}, \boldsymbol{\eta})) \in H(\mathbf{div}; \Omega) \times L^2(\Omega) \times Q$ , where

$$H(\mathbf{div}; \Omega) := \{ \boldsymbol{\tau} \in [L^2(\Omega)]^{2 \times 2} : \mathbf{div}(\boldsymbol{\tau}) \in [L^2(\Omega)]^2 \} \quad \text{and} \quad Q := [L^2(\Omega)]^2 \times [L^2(\Omega)]_{\text{asym}}^{2 \times 2},$$

with

$$[L^2(\Omega)]_{\text{asym}}^{2 \times 2} := \{ \boldsymbol{\eta} \in [L^2(\Omega)]^{2 \times 2} : \boldsymbol{\eta} + \boldsymbol{\eta}^{\dagger} = \mathbf{0} \}.$$

Now, noting that

$$\boldsymbol{\sigma} : \boldsymbol{\tau} + p \operatorname{tr}(\boldsymbol{\sigma}) + q \operatorname{tr}(\boldsymbol{\tau}) + 2pq = \boldsymbol{\sigma}^{\text{d}} : \boldsymbol{\tau}^{\text{d}} + 2 \left( p + \frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma}) \right) \left( q + \frac{1}{2} \operatorname{tr}(\boldsymbol{\tau}) \right),$$

the last system can be written in the more compact form: Find  $(\boldsymbol{\sigma}, p, (\mathbf{u}, \boldsymbol{\gamma}))$  in  $H(\mathbf{div}; \Omega) \times L^2(\Omega) \times Q$  such that

$$\begin{aligned} \frac{1}{2\mu} \int_{\Omega} \boldsymbol{\sigma}^{\text{d}} : \boldsymbol{\tau}^{\text{d}} + \frac{1}{\mu} \int_{\Omega} \left( p + \frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma}) \right) \left( q + \frac{1}{2} \operatorname{tr}(\boldsymbol{\tau}) \right) + \int_{\Omega} \mathbf{u} \cdot \mathbf{div}(\boldsymbol{\tau}) + \int_{\Omega} \boldsymbol{\gamma} : \boldsymbol{\tau} &= 0, \\ \int_{\Omega} \mathbf{v} \cdot \mathbf{div}(\boldsymbol{\sigma}) + \int_{\Omega} \boldsymbol{\eta} : \boldsymbol{\sigma} &= - \int_{\Omega} \mathbf{f} \cdot \mathbf{v}, \end{aligned} \tag{2.5}$$

for all  $(\boldsymbol{\tau}, q, (\mathbf{v}, \boldsymbol{\eta})) \in H(\mathbf{div}; \Omega) \times L^2(\Omega) \times Q$ . At this point we observe that for any  $c \in \mathbb{R}$ ,  $(c\mathbf{I}, -c, (\mathbf{0}, \mathbf{0}))$  is a solution of the homogeneous version of system (2.5). Hence, in order to avoid this non-uniqueness we consider the decomposition

$$H(\mathbf{div}; \Omega) = H_0 \oplus \mathbb{R}\mathbf{I}, \tag{2.6}$$

where  $H_0 := \left\{ \boldsymbol{\tau} \in H(\mathbf{div}; \Omega) : \int_{\Omega} \operatorname{tr}(\boldsymbol{\tau}) = 0 \right\}$ , and require from now on that  $\boldsymbol{\sigma} \in H_0$ .

The following lemma guarantees that the test space can also be restricted to  $H_0$ .

**LEMMA 2.1.** *Any solution of (2.5) with  $\boldsymbol{\sigma} \in H_0$  is also solution of: Find  $(\boldsymbol{\sigma}, p, (\mathbf{u}, \boldsymbol{\gamma})) \in H_0 \times L^2(\Omega) \times Q$  such that*

$$\begin{aligned} \frac{1}{2\mu} \int_{\Omega} \boldsymbol{\sigma}^{\text{d}} : \boldsymbol{\tau}^{\text{d}} + \frac{1}{\mu} \int_{\Omega} \left( p + \frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma}) \right) \left( q + \frac{1}{2} \operatorname{tr}(\boldsymbol{\tau}) \right) + \int_{\Omega} \mathbf{u} \cdot \mathbf{div}(\boldsymbol{\tau}) + \int_{\Omega} \boldsymbol{\gamma} : \boldsymbol{\tau} &= 0, \\ \int_{\Omega} \mathbf{v} \cdot \mathbf{div}(\boldsymbol{\sigma}) + \int_{\Omega} \boldsymbol{\eta} : \boldsymbol{\sigma} &= - \int_{\Omega} \mathbf{f} \cdot \mathbf{v}, \end{aligned} \tag{2.7}$$

for all  $(\boldsymbol{\tau}, q, (\mathbf{v}, \boldsymbol{\eta})) \in H_0 \times L^2(\Omega) \times Q$ . Conversely, any solution of (2.7) is also a solution of (2.5).

*Proof.* It is immediate that any solution of (2.5) with  $\boldsymbol{\sigma} \in H_0$  is also a solution of (2.7). Conversely, let  $(\boldsymbol{\sigma}, p, (\mathbf{u}, \boldsymbol{\gamma}))$  be a solution of (2.7). Because of (2.6) it suffices to prove that  $(\boldsymbol{\sigma}, p, (\mathbf{u}, \boldsymbol{\gamma}))$  also satisfies (2.5) if tested with  $(\mathbf{I}, 0, (\mathbf{0}, \mathbf{0}))$ . This requires that  $\int_{\Omega} (p + \frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma}))$  vanishes which can be seen to be true by selecting  $(\boldsymbol{\tau}, q, (\mathbf{v}, \boldsymbol{\eta})) = (\mathbf{0}, 1, (\mathbf{0}, \mathbf{0})) \in H_0 \times L^2(\Omega) \times Q$  in (2.7).  $\square$

Furthermore, we now let  $H := H_0 \times L^2(\Omega)$ , consider a constant  $\kappa_0 > 0$ , and introduce a generalized version of (2.7): Find  $((\boldsymbol{\sigma}, p), (\mathbf{u}, \boldsymbol{\gamma}))$  in  $H \times Q$  such that

$$\begin{aligned} a((\boldsymbol{\sigma}, p), (\boldsymbol{\tau}, q)) + b(\boldsymbol{\tau}, (\mathbf{u}, \boldsymbol{\gamma})) &= 0 & \forall (\boldsymbol{\tau}, q) \in H, \\ b(\boldsymbol{\sigma}, (\mathbf{v}, \boldsymbol{\eta})) &= - \int_{\Omega} \mathbf{f} \cdot \mathbf{v} & \forall (\mathbf{v}, \boldsymbol{\eta}) \in Q, \end{aligned} \quad (2.8)$$

where  $a : H \times H \rightarrow \mathbb{R}$  and  $b : H_0 \times Q \rightarrow \mathbb{R}$  are the bounded bilinear forms defined by

$$a((\boldsymbol{\zeta}, r), (\boldsymbol{\tau}, q)) := \frac{1}{2\mu} \int_{\Omega} \boldsymbol{\zeta}^{\text{d}} : \boldsymbol{\tau}^{\text{d}} + \frac{\kappa_0}{\mu} \int_{\Omega} \left( r + \frac{1}{2} \text{tr}(\boldsymbol{\zeta}) \right) \left( q + \frac{1}{2} \text{tr}(\boldsymbol{\tau}) \right) \quad (2.9)$$

and

$$b(\boldsymbol{\zeta}, (\mathbf{v}, \boldsymbol{\eta})) := \int_{\Omega} \mathbf{v} \cdot \text{div}(\boldsymbol{\zeta}) + \int_{\Omega} \boldsymbol{\eta} : \boldsymbol{\zeta} \quad (2.10)$$

for  $(\boldsymbol{\zeta}, r)$ ,  $(\boldsymbol{\tau}, q)$  in  $H$  and  $(\mathbf{v}, \boldsymbol{\eta})$  in  $Q$ . Note that (2.7) corresponds to (2.8) with  $\kappa_0 = 1$ .

In order to show that the formulations (2.8) are independent of  $\kappa_0 > 0$ , we prove next that they are all equivalent to the simplified version arising after replacing the incompressibility condition (2.2) into (2.8) (equivalently, taking  $\kappa_0 = 0$  in (2.8)), that is: Find  $(\boldsymbol{\sigma}, (\mathbf{u}, \boldsymbol{\gamma})) \in H_0 \times Q$  such that

$$\begin{aligned} a_0(\boldsymbol{\sigma}, \boldsymbol{\tau}) + b(\boldsymbol{\tau}, (\mathbf{u}, \boldsymbol{\gamma})) &= 0 & \forall \boldsymbol{\tau} \in H_0, \\ b(\boldsymbol{\sigma}, (\mathbf{v}, \boldsymbol{\eta})) &= - \int_{\Omega} \mathbf{f} \cdot \mathbf{v} & \forall (\mathbf{v}, \boldsymbol{\eta}) \in Q, \end{aligned} \quad (2.11)$$

where  $a_0 : H_0 \times H_0 \rightarrow \mathbb{R}$  is the bounded bilinear form defined by

$$a_0(\boldsymbol{\zeta}, \boldsymbol{\tau}) := \frac{1}{2\mu} \int_{\Omega} \boldsymbol{\zeta}^{\text{d}} : \boldsymbol{\tau}^{\text{d}} \quad \forall (\boldsymbol{\zeta}, \boldsymbol{\tau}) \in H_0 \times H_0.$$

LEMMA 2.2. *Problems (2.8) and (2.11) are equivalent. Indeed,  $((\boldsymbol{\sigma}, p), (\mathbf{u}, \boldsymbol{\gamma})) \in H \times Q$  is a solution of (2.8) if and only if  $(\boldsymbol{\sigma}, (\mathbf{u}, \boldsymbol{\gamma})) \in H_0 \times Q$  is a solution of (2.11) and  $p = -\frac{1}{2} \text{tr}(\boldsymbol{\sigma})$ .*

*Proof.* It suffices to take  $\boldsymbol{\tau} = \mathbf{0}$  in (2.8) and then use that the traces of the tensor-valued functions in  $H(\text{div}; \Omega)$  live in  $L^2(\Omega)$  as the pressure test functions do.  $\square$

The following lemmata will be useful in order to prove well-posedness of (2.8) and (2.11).

LEMMA 2.3. *There exists a positive constant  $\beta$ , depending only on  $\Omega$  such that*

$$\sup_{\substack{\boldsymbol{\tau} \in H(\text{div}; \Omega) \\ \boldsymbol{\tau} \neq \mathbf{0}}} \frac{\int_{\Omega} \mathbf{v} \cdot \text{div}(\boldsymbol{\tau}) + \int_{\Omega} \boldsymbol{\eta} : \boldsymbol{\tau}}{\|\boldsymbol{\tau}\|_{H(\text{div}; \Omega)}} \geq \beta \|(\mathbf{v}, \boldsymbol{\eta})\|_Q \quad (2.12)$$

for all  $(\mathbf{v}, \boldsymbol{\eta})$  in  $Q$ .

*Proof.* See Lemma 4.3 in [4] for a detailed proof.  $\square$

LEMMA 2.4. *There exists  $c_1 > 0$ , depending only on  $\Omega$ , such that*

$$c_1 \|\boldsymbol{\tau}\|_{[L^2(\Omega)]^{2 \times 2}}^2 \leq \|\boldsymbol{\tau}^{\text{d}}\|_{[L^2(\Omega)]^{2 \times 2}}^2 + \|\text{div}(\boldsymbol{\tau})\|_{[L^2(\Omega)]^2}^2 \quad \forall \boldsymbol{\tau} \in H_0, \quad (2.13)$$

*Proof.* See Lemma 3.1 in [3] or Proposition 3.1 of Chapter IV in [7].  $\square$

We are now in a position to state the following theorem.

**THEOREM 2.5.** *Problem (2.11) has a unique solution  $(\boldsymbol{\sigma}, (\mathbf{u}, \boldsymbol{\gamma})) \in H_0 \times Q$ . Moreover, there exists a positive constant  $C$ , depending only on  $\Omega$ , such that*

$$\|(\boldsymbol{\sigma}, (\mathbf{u}, \boldsymbol{\gamma}))\|_{H(\mathbf{div}; \Omega) \times Q} \leq C \|\mathbf{f}\|_{[L^2(\Omega)]^2}.$$

*Proof.* It suffices to prove that the bilinear forms  $a_0$  and  $b$  satisfy the hypotheses of the Babuška-Brezzi theory. Indeed, given  $(\mathbf{v}, \boldsymbol{\eta})$  in  $Q$  it is easy to see that

$$\sup_{\substack{\boldsymbol{\tau} \in H_0 \\ \boldsymbol{\tau} \neq 0}} \frac{\int_{\Omega} \mathbf{v} \cdot \mathbf{div}(\boldsymbol{\tau}) + \int_{\Omega} \boldsymbol{\eta} : \boldsymbol{\tau}}{\|\boldsymbol{\tau}\|_{H(\mathbf{div}; \Omega)}} = \sup_{\substack{\boldsymbol{\tau} \in H(\mathbf{div}; \Omega) \\ \boldsymbol{\tau} \neq 0}} \frac{\int_{\Omega} \mathbf{v} \cdot \mathbf{div}(\boldsymbol{\tau}) + \int_{\Omega} \boldsymbol{\eta} : \boldsymbol{\tau}}{\|\boldsymbol{\tau}\|_{H(\mathbf{div}; \Omega)}}, \quad (2.14)$$

which, together with Lemma 2.3, proves the continuous inf-sup condition for  $b$ . Now, let  $V$  be the kernel of the operator induced by  $b$ , that is

$$\begin{aligned} V &:= \{ \boldsymbol{\tau} \in H_0 : b(\boldsymbol{\tau}, (\mathbf{v}, \boldsymbol{\eta})) = 0 \quad \forall (\mathbf{v}, \boldsymbol{\eta}) \in Q \} \\ &= \{ \boldsymbol{\tau} \in H_0 : \mathbf{div}(\boldsymbol{\tau}) = 0 \quad \text{and} \quad \boldsymbol{\tau} = \boldsymbol{\tau}^t \quad \text{in} \quad \Omega \}. \end{aligned}$$

It follows, applying Lemma 2.4, that for each  $\boldsymbol{\tau} \in V$  there holds

$$a_0(\boldsymbol{\tau}, \boldsymbol{\tau}) = \frac{1}{2\mu} \|\boldsymbol{\tau}^d\|_{[L^2(\Omega)]^{2 \times 2}}^2 \geq \frac{c_1}{2\mu} \|\boldsymbol{\tau}\|_{[L^2(\Omega)]^{2 \times 2}}^2 = \frac{c_1}{2\mu} \|\boldsymbol{\tau}\|_{H(\mathbf{div}; \Omega)}^2,$$

which shows that the bilinear form  $a_0$  is strongly coercive in  $V$ . Finally, a straightforward application of the classical result given by Theorem 1.1 in Chapter II of [7] completes the proof.  $\square$

**THEOREM 2.6.** *Problem (2.8) has a unique solution  $((\boldsymbol{\sigma}, p), (\mathbf{u}, \boldsymbol{\gamma})) \in H \times Q$ , independent of  $\kappa_0$ , and there holds  $p = -\frac{1}{2} \text{tr}(\boldsymbol{\sigma})$ . Moreover, there exists a constant  $C > 0$ , depending only on  $\Omega$ , such that*

$$\|((\boldsymbol{\sigma}, p), (\mathbf{u}, \boldsymbol{\gamma}))\|_{H \times Q} \leq C \|\mathbf{f}\|_{[L^2(\Omega)]^2}.$$

*Proof.* It is a direct consequence of Lemma 2.2, which gives the equivalence between (2.8) and (2.11), and Theorem 2.5, which yields the well-posedness of (2.11).  $\square$

### 3. THE AUGMENTED DUAL-MIXED VARIATIONAL FORMULATIONS

In the following we enrich the formulations (2.8) and (2.11) with residuals arising from the modified constitutive equation (2.4), the equilibrium equation, and the relation defining the rotation as a function of the displacement. More precisely, as in [18] we subtract the second from the first equation in both (2.8) and (2.11) and then add the Galerkin least-squares terms given by

$$\kappa_1 \int_{\Omega} \left( \boldsymbol{\varepsilon}(\mathbf{u}) - \frac{1}{2\mu} \boldsymbol{\sigma}^d \right) : \left( \boldsymbol{\varepsilon}(\mathbf{v}) + \frac{1}{2\mu} \boldsymbol{\tau}^d \right) = 0, \quad (3.1)$$

$$\kappa_2 \int_{\Omega} \mathbf{div}(\boldsymbol{\sigma}) \cdot \mathbf{div}(\boldsymbol{\tau}) = -\kappa_2 \int_{\Omega} \mathbf{f} \cdot \mathbf{div}(\boldsymbol{\tau}), \quad (3.2)$$

and

$$\kappa_3 \int_{\Omega} \left( \boldsymbol{\gamma} - \frac{1}{2} (\nabla \mathbf{u} - (\nabla \mathbf{u})^t) \right) : \left( \boldsymbol{\eta} + \frac{1}{2} (\nabla \mathbf{v} - (\nabla \mathbf{v})^t) \right) = 0, \quad (3.3)$$

for all  $(\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta}) \in H_0 \times [H_0^1(\Omega)]^2 \times [L^2(\Omega)]_{\text{asym}}^{2 \times 2}$ , where  $(\kappa_1, \kappa_2, \kappa_3)$  is a vector of positive parameters to be specified later. We notice that (3.1) and (3.3) implicitly require now the displacement  $\mathbf{u}$  to live in the smaller space  $[H_0^1(\Omega)]^2$ .

In this way, instead of (2.8) we propose the following augmented variational formulation: Find  $(\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma}) \in \mathbf{H} := H_0 \times L^2(\Omega) \times [H_0^1(\Omega)]^2 \times [L^2(\Omega)]_{\text{asym}}^{2 \times 2}$  such that

$$A((\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma}), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})) = F(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}) \quad \forall (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}) \in \mathbf{H}, \quad (3.4)$$

where the bilinear form  $A : \mathbf{H} \times \mathbf{H} \longrightarrow \mathbb{R}$  and the functional  $F : \mathbf{H} \longrightarrow \mathbb{R}$  are defined by

$$\begin{aligned} A((\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma}), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})) &:= a((\boldsymbol{\sigma}, p), (\boldsymbol{\tau}, q)) + b(\boldsymbol{\tau}, (\mathbf{u}, \boldsymbol{\gamma})) - b(\boldsymbol{\sigma}, (\mathbf{v}, \boldsymbol{\eta})) \\ &+ \kappa_1 \int_{\Omega} \left( \boldsymbol{\varepsilon}(\mathbf{u}) - \frac{1}{2\mu} \boldsymbol{\sigma}^d \right) : \left( \boldsymbol{\varepsilon}(\mathbf{v}) + \frac{1}{2\mu} \boldsymbol{\tau}^d \right) + \kappa_2 \int_{\Omega} \mathbf{div}(\boldsymbol{\sigma}) \cdot \mathbf{div}(\boldsymbol{\tau}) \\ &+ \kappa_3 \int_{\Omega} \left( \boldsymbol{\gamma} - \frac{1}{2} (\nabla \mathbf{u} - (\nabla \mathbf{u})^t) \right) : \left( \boldsymbol{\eta} + \frac{1}{2} (\nabla \mathbf{v} - (\nabla \mathbf{v})^t) \right) \end{aligned} \quad (3.5)$$

and

$$F(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}) := \int_{\Omega} \mathbf{f} \cdot (\mathbf{v} - \kappa_2 \mathbf{div}(\boldsymbol{\tau})) . \quad (3.6)$$

Similarly, instead of (2.11) we propose: Find  $(\boldsymbol{\sigma}, \mathbf{u}, \boldsymbol{\gamma}) \in \mathbf{H}_0 := H_0 \times [H_0^1(\Omega)]^2 \times [L^2(\Omega)]_{\text{asym}}^{2 \times 2}$  such that

$$A_0((\boldsymbol{\sigma}, \mathbf{u}, \boldsymbol{\gamma}), (\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta})) = F_0(\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta}) \quad \forall (\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta}) \in \mathbf{H}_0 , \quad (3.7)$$

where the bilinear form  $A_0 : \mathbf{H}_0 \times \mathbf{H}_0 \longrightarrow \mathbb{R}$  and the functional  $F_0 : \mathbf{H}_0 \longrightarrow \mathbb{R}$  are defined by

$$\begin{aligned} A_0((\boldsymbol{\sigma}, \mathbf{u}, \boldsymbol{\gamma}), (\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta})) &:= a_0(\boldsymbol{\sigma}, \boldsymbol{\tau}) + b(\boldsymbol{\tau}, (\mathbf{u}, \boldsymbol{\gamma})) - b(\boldsymbol{\sigma}, (\mathbf{v}, \boldsymbol{\eta})) \\ &+ \kappa_1 \int_{\Omega} \left( \boldsymbol{\varepsilon}(\mathbf{u}) - \frac{1}{2\mu} \boldsymbol{\sigma}^d \right) : \left( \boldsymbol{\varepsilon}(\mathbf{v}) + \frac{1}{2\mu} \boldsymbol{\tau}^d \right) + \kappa_2 \int_{\Omega} \mathbf{div}(\boldsymbol{\sigma}) \cdot \mathbf{div}(\boldsymbol{\tau}) \\ &+ \kappa_3 \int_{\Omega} \left( \boldsymbol{\gamma} - \frac{1}{2} (\nabla \mathbf{u} - (\nabla \mathbf{u})^t) \right) : \left( \boldsymbol{\eta} + \frac{1}{2} (\nabla \mathbf{v} - (\nabla \mathbf{v})^t) \right) \end{aligned} \quad (3.8)$$

and

$$F_0(\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta}) := \int_{\Omega} \mathbf{f} \cdot (\mathbf{v} - \kappa_2 \mathbf{div}(\boldsymbol{\tau})) . \quad (3.9)$$

The analogue of Lemma 2.2 is given now.

**LEMMA 3.1.** *Problems (3.4) and (3.7) are equivalent. Indeed,  $(\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma}) \in \mathbf{H}$  is a solution of (3.4) if and only if  $(\boldsymbol{\sigma}, \mathbf{u}, \boldsymbol{\gamma}) \in \mathbf{H}_0$  is a solution of (3.7) and  $p = -\frac{1}{2} \text{tr}(\boldsymbol{\sigma})$ .*

*Proof.* It suffices to take  $(\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta}) = (\mathbf{0}, \mathbf{0}, \mathbf{0})$  in (3.4) and then use again that the traces of the tensor-valued functions in  $H(\mathbf{div}; \Omega)$  live in  $L^2(\Omega)$  as the pressure test functions do.  $\square$

In what follows we aim to show the well-posedness of (3.7). The main idea is to choose the vector of parameters  $(\kappa_1, \kappa_2, \kappa_3)$  such that  $A_0$  be strongly coercive on  $\mathbf{H}_0$  with respect to the norm  $\|\cdot\|_{\mathbf{H}_0}$  defined by

$$\|(\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta})\|_{\mathbf{H}_0} := \left\{ \|\boldsymbol{\tau}\|_{H(\mathbf{div}; \Omega)}^2 + |\mathbf{v}|_{[H^1(\Omega)]^2}^2 + \|\boldsymbol{\eta}\|_{[L^2(\Omega)]^{2 \times 2}}^2 \right\}^{1/2} .$$

We first notice, after simple computations, that

$$\int_{\Omega} \left( \boldsymbol{\varepsilon}(\mathbf{v}) - \frac{1}{2\mu} \boldsymbol{\tau}^d \right) : \left( \boldsymbol{\varepsilon}(\mathbf{v}) + \frac{1}{2\mu} \boldsymbol{\tau}^d \right) = \|\boldsymbol{\varepsilon}(\mathbf{v})\|_{[L^2(\Omega)]^{2 \times 2}}^2 - \frac{1}{4\mu^2} \|\boldsymbol{\tau}^d\|_{[L^2(\Omega)]^{2 \times 2}}^2 ,$$

and that

$$\begin{aligned} \int_{\Omega} \left( \boldsymbol{\eta} - \frac{1}{2} (\nabla \mathbf{v} - (\nabla \mathbf{v})^t) \right) : \left( \boldsymbol{\eta} + \frac{1}{2} (\nabla \mathbf{v} - (\nabla \mathbf{v})^t) \right) \\ = \|\boldsymbol{\eta}\|_{[L^2(\Omega)]^{2 \times 2}}^2 + \|\boldsymbol{\varepsilon}(\mathbf{v})\|_{[L^2(\Omega)]^{2 \times 2}}^2 - |\mathbf{v}|_{[H^1(\Omega)]^2}^2 , \end{aligned}$$

which gives

$$\begin{aligned} A_0((\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta}), (\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta})) &= \frac{1}{2\mu} \left( 1 - \frac{\kappa_1}{2\mu} \right) \|\boldsymbol{\tau}^d\|_{[L^2(\Omega)]^{2 \times 2}}^2 + \kappa_2 \|\mathbf{div}(\boldsymbol{\tau})\|_{[L^2(\Omega)]^2}^2 \\ &+ (\kappa_1 + \kappa_3) \|\boldsymbol{\varepsilon}(\mathbf{v})\|_{[L^2(\Omega)]^{2 \times 2}}^2 - \kappa_3 |\mathbf{v}|_{[H^1(\Omega)]^2}^2 + \kappa_3 \|\boldsymbol{\eta}\|_{[L^2(\Omega)]^{2 \times 2}}^2 . \end{aligned} \quad (3.10)$$

Now, Korn's first inequality (see, e.g., Theorem 10.1 in [22]) establishes that

$$\|\boldsymbol{\varepsilon}(\mathbf{v})\|_{[L^2(\Omega)]^2}^2 \geq \frac{1}{2} |\mathbf{v}|_{[H^1(\Omega)]^2}^2 \quad \forall \mathbf{v} \in [H_0^1(\Omega)]^2, \quad (3.11)$$

and hence (3.10) yields

$$\begin{aligned} A_0((\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta}), (\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta})) &\geq \frac{1}{2\mu} \left(1 - \frac{\kappa_1}{2\mu}\right) \|\boldsymbol{\tau}^d\|_{[L^2(\Omega)]^{2 \times 2}}^2 + \kappa_2 \|\mathbf{div}(\boldsymbol{\tau})\|_{[L^2(\Omega)]^2}^2 \\ &+ \frac{(\kappa_1 - \kappa_3)}{2} |\mathbf{v}|_{[H^1(\Omega)]^2}^2 + \kappa_3 \|\boldsymbol{\eta}\|_{[L^2(\Omega)]^{2 \times 2}}^2. \end{aligned}$$

Then, choosing  $\kappa_1$  and  $\kappa_2$  such that

$$0 < \kappa_1 < 2\mu \quad \text{and} \quad 0 < \kappa_2,$$

and applying Lemma 2.4, we deduce that

$$A_0((\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta}), (\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta})) \geq \alpha_2 \|\boldsymbol{\tau}\|_{H(\mathbf{div}; \Omega)}^2 + \frac{(\kappa_1 - \kappa_3)}{2} |\mathbf{v}|_{[H^1(\Omega)]^2}^2 + \kappa_3 \|\boldsymbol{\eta}\|_{[L^2(\Omega)]^{2 \times 2}}^2,$$

where

$$\alpha_2 := \min \left\{ c_1 \alpha_1, \frac{\kappa_2}{2} \right\}, \quad \alpha_1 := \min \left\{ \frac{1}{2\mu} \left(1 - \frac{\kappa_1}{2\mu}\right), \frac{\kappa_2}{2} \right\},$$

and  $c_1$  is the constant that appears in Lemma 2.4. In addition, choosing the parameter  $\kappa_3$  such that  $0 < \kappa_3 < \kappa_1$ , we find that

$$A_0((\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta}), (\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta})) \geq \alpha \|(\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta})\|_{\mathbf{H}_0}^2 \quad \forall (\boldsymbol{\tau}, \mathbf{v}, \boldsymbol{\eta}) \in \mathbf{H}_0, \quad (3.12)$$

where

$$\alpha := \min \left\{ \alpha_2, \frac{(\kappa_1 - \kappa_3)}{2}, \kappa_3 \right\}.$$

As a consequence of the above analysis, we obtain the following main results.

**THEOREM 3.2.** *Assume that there hold*

$$0 < \kappa_1 < 2\mu, \quad 0 < \kappa_2, \quad \text{and} \quad 0 < \kappa_3 < \kappa_1.$$

*Then, the augmented variational formulation (3.7) has a unique solution  $(\boldsymbol{\sigma}, \mathbf{u}, \boldsymbol{\gamma}) \in \mathbf{H}_0$ . Moreover, there exists a positive constant  $C$ , depending only on  $\mu$  and  $(\kappa_1, \kappa_2, \kappa_3)$ , such that  $\|(\boldsymbol{\sigma}, \mathbf{u}, \boldsymbol{\gamma})\|_{\mathbf{H}_0} \leq C \|F_0\|_{\mathbf{H}'_0} \leq C \|\mathbf{f}\|_{[L^2(\Omega)]^2}$ .*

*Proof.* It is clear from (3.8) and (3.12) that  $A_0$  is bounded and strongly coercive on  $\mathbf{H}_0$  with constants depending on  $\mu$  and  $(\kappa_1, \kappa_2, \kappa_3)$ . Also, the linear functional  $F_0$  (cf. (3.9)) is clearly continuous with norm bounded by  $(1 + \kappa_2) \|\mathbf{f}\|_{[L^2(\Omega)]^2}$ . Therefore, the assertion is a simple consequence of the Lax-Milgram Lemma.  $\square$

**THEOREM 3.3.** *Assume that there hold*

$$0 < \kappa_1 < 2\mu, \quad 0 < \kappa_2, \quad \text{and} \quad 0 < \kappa_3 < \kappa_1.$$

*Then the augmented variational formulation (3.4) has a unique solution  $(\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma}) \in \mathbf{H}$ , independent of  $\kappa_0$ , and there holds  $p = -\frac{1}{2} \text{tr}(\boldsymbol{\sigma})$ . Moreover, there exists a positive constant  $C$ , depending only on  $\mu$  and  $(\kappa_1, \kappa_2, \kappa_3)$ , such that  $\|(\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma})\|_{\mathbf{H}} \leq C \|\mathbf{f}\|_{[L^2(\Omega)]^2}$ .*

*Proof.* It is a direct consequence of Lemma 3.1 and Theorem 3.2.  $\square$

We end this section by emphasizing that the introduction of the augmented formulations (3.4) and (3.7) is motivated by the possibility of using arbitrary finite element subspaces in the definition of the associated Galerkin schemes. This is certainly guaranteed by the strong coerciveness of the resulting bilinear form, as already proved for  $A_0$  (cf. (3.12)) and as will be proved for  $A$  in the next section. We also remark here that at first glance it could seem, due to Lemmata 2.2 and 3.1, that there is actually no need of considering the continuous variational formulations (2.8) and (3.4) since the equivalent ones, given respectively by (2.11) and (3.7), are clearly simpler. Nevertheless, as we show below in Section 4, the main interest in (2.8) and particularly in the corresponding augmented formulation (3.4) lies in the associated Galerkin scheme, which provides more flexibility for choosing the pressure finite element subspace.

#### 4. THE AUGMENTED MIXED FINITE ELEMENT METHODS

We now let  $H_{0,h}^\sigma$ ,  $H_h^p$ ,  $H_{0,h}^u$  and  $H_h^\gamma$  be arbitrary finite element subspaces of  $H_0$ ,  $L^2(\Omega)$ ,  $[H_0^1(\Omega)]^2$  and  $[L^2(\Omega)]_{\text{asym}}^{2 \times 2}$ , respectively, and define

$$\mathbf{H}_h := H_{0,h}^\sigma \times H_h^p \times H_{0,h}^u \times H_h^\gamma \quad \text{and} \quad \mathbf{H}_{0,h} := H_{0,h}^\sigma \times H_{0,h}^u \times H_h^\gamma.$$

In addition, let  $\kappa_0$ ,  $\kappa_1$ ,  $\kappa_2$ , and  $\kappa_3$  be given positive parameters. Then, the Galerkin schemes associated with (3.4) and (3.7) read: Find  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_h$  such that

$$A((\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h), (\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h)) = F(\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \quad \forall (\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \in \mathbf{H}_h, \quad (4.1)$$

and: Find  $(\boldsymbol{\sigma}_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_{0,h}$  such that

$$A_0((\boldsymbol{\sigma}_h, \mathbf{u}_h, \boldsymbol{\gamma}_h), (\boldsymbol{\tau}_h, \mathbf{v}_h, \boldsymbol{\eta}_h)) = F_0(\boldsymbol{\tau}_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \quad \forall (\boldsymbol{\tau}_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \in \mathbf{H}_{0,h}. \quad (4.2)$$

The following theorem provides the unique solvability, stability, and convergence of (4.2).

**THEOREM 4.1.** *Assume that the parameters  $\kappa_1$ ,  $\kappa_2$ , and  $\kappa_3$  satisfy the assumptions of Theorem 3.2 and let  $\mathbf{H}_{0,h}$  be any finite element subspace of  $\mathbf{H}_0$ . Then, the Galerkin scheme (4.2) has a unique solution  $(\boldsymbol{\sigma}_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_{0,h}$ , and there exist positive constants  $C$ ,  $\tilde{C}$ , independent of  $h$ , such that*

$$\|(\boldsymbol{\sigma}_h, \mathbf{u}_h, \boldsymbol{\gamma}_h)\|_{\mathbf{H}_0} \leq C \sup_{\substack{(\boldsymbol{\tau}_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \in \mathbf{H}_{0,h} \\ (\boldsymbol{\tau}_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \neq 0}} \frac{|F_0(\boldsymbol{\tau}_h, \mathbf{v}_h, \boldsymbol{\eta}_h)|}{\|(\boldsymbol{\tau}_h, \mathbf{v}_h, \boldsymbol{\eta}_h)\|_{\mathbf{H}_0}} \leq C \|\mathbf{f}\|_{[L^2(\Omega)]^2},$$

and

$$\|(\boldsymbol{\sigma}, \mathbf{u}, \boldsymbol{\gamma}) - (\boldsymbol{\sigma}_h, \mathbf{u}_h, \boldsymbol{\gamma}_h)\|_{\mathbf{H}_0} \leq \tilde{C} \inf_{(\boldsymbol{\tau}_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \in \mathbf{H}_{0,h}} \|(\boldsymbol{\sigma}, \mathbf{u}, \boldsymbol{\gamma}) - (\boldsymbol{\tau}_h, \mathbf{v}_h, \boldsymbol{\eta}_h)\|_{\mathbf{H}_0}. \quad (4.3)$$

*Proof.* Since  $A_0$  is bounded and strongly coercive on  $\mathbf{H}_0$  (cf. (3.8) and (3.12)) with constants depending on  $\mu$  and  $(\kappa_1, \kappa_2, \kappa_3)$ , the proof follows from a straightforward application of the Lax-Milgram Lemma and Cea's estimate.  $\square$

In order to define an explicit finite element subspace of  $\mathbf{H}_0$ , we now let  $\{\mathcal{T}_h\}_{h>0}$  be a regular family of triangulations of the polygonal region  $\bar{\Omega}$  by triangles  $T$  of diameter  $h_T$  such that  $\bar{\Omega} = \cup \{T : T \in \mathcal{T}_h\}$  and define  $h := \max \{h_T : T \in \mathcal{T}_h\}$ . Given an integer  $\ell \geq 0$  and a subset  $S$  of  $\mathbb{R}^2$ , we denote by  $\mathbb{P}_\ell(S)$  the space of polynomials of total degree at most  $\ell$  defined on  $S$ . Also, for each  $T \in \mathcal{T}_h$  we define the local Raviart-Thomas space of order zero

$$\mathbb{RT}_0(T) := \text{span} \left\{ \begin{pmatrix} 1 \\ 0 \end{pmatrix}, \begin{pmatrix} 0 \\ 1 \end{pmatrix}, \begin{pmatrix} x_1 \\ x_2 \end{pmatrix} \right\} \subseteq [\mathbb{P}_1(T)]^2,$$



where  $\begin{pmatrix} x_1 \\ x_2 \end{pmatrix}$  is a generic vector of  $\mathbb{R}^2$ , and let  $\tilde{H}_h^\sigma$  be the corresponding global space, that is

$$\tilde{H}_h^\sigma := \{ \boldsymbol{\tau}_h \in H(\mathbf{div}; \Omega) : \boldsymbol{\tau}_h|_T \in [\mathbb{RT}_0(T)^t]^2 \quad \forall T \in \mathcal{T}_h \}. \quad (4.4)$$

Then we let  $\tilde{\mathbf{H}}_{0,h} := \tilde{H}_{0,h}^\sigma \times \tilde{H}_{0,h}^{\mathbf{u}} \times \tilde{H}_h^\gamma$ , where

$$\tilde{H}_{0,h}^\sigma := \left\{ \boldsymbol{\tau}_h \in \tilde{H}_h^\sigma : \int_{\Omega} \operatorname{tr}(\boldsymbol{\tau}_h) = 0 \right\}, \quad (4.5)$$

$$\tilde{H}_{0,h}^{\mathbf{u}} := \{ \mathbf{v}_h \in [C(\bar{\Omega})]^2 : \mathbf{v}_h|_T \in [\mathbb{P}_1(T)]^2 \quad \forall T \in \mathcal{T}_h, \quad \mathbf{v}_h = 0 \text{ on } \partial\Omega \}, \quad (4.6)$$

and

$$\tilde{H}_h^\gamma := \{ \boldsymbol{\eta}_h \in [L^2(\Omega)]_{\text{asym}}^{2 \times 2} : \boldsymbol{\eta}_h|_T \in [\mathbb{P}_0(T)]^{2 \times 2} \quad \forall T \in \mathcal{T}_h \}. \quad (4.7)$$

The approximation properties of these subspaces are given as follows (see [7], [11], [18]):

(AP $_{0,h}^\sigma$ ) For each  $r \in (0, 1]$  and for each  $\boldsymbol{\tau} \in [H^r(\Omega)]^{2 \times 2} \cap H_0$  with  $\mathbf{div}(\boldsymbol{\tau}) \in [H^r(\Omega)]^2$  there exists  $\boldsymbol{\tau}_h \in \tilde{H}_{0,h}^\sigma$  such that

$$\| \boldsymbol{\tau} - \boldsymbol{\tau}_h \|_{H(\mathbf{div}; \Omega)} \leq C h^r \left\{ \| \boldsymbol{\tau} \|_{[H^r(\Omega)]^{2 \times 2}} + \| \mathbf{div}(\boldsymbol{\tau}) \|_{[H^r(\Omega)]^2} \right\}.$$

(AP $_{0,h}^{\mathbf{u}}$ ) For each  $r \in [1, 2]$  and for each  $\mathbf{v} \in [H^{1+r}(\Omega)]^2 \cap [H_0^1(\Omega)]^2$  there exists  $\mathbf{v}_h \in \tilde{H}_{0,h}^{\mathbf{u}}$  such that

$$\| \mathbf{v} - \mathbf{v}_h \|_{[H^1(\Omega)]^2} \leq C h^r \| \mathbf{v} \|_{[H^{1+r}(\Omega)]^2}.$$

(AP $_h^\gamma$ ) For each  $r \in [0, 1]$  and for each  $\boldsymbol{\eta} \in [H^r(\Omega)]^{2 \times 2} \cap [L^2(\Omega)]_{\text{asym}}^{2 \times 2}$  there exists  $\boldsymbol{\eta}_h \in \tilde{H}_h^\gamma$  such that

$$\| \boldsymbol{\eta} - \boldsymbol{\eta}_h \|_{[L^2(\Omega)]^{2 \times 2}} \leq C h^r \| \boldsymbol{\eta} \|_{[H^r(\Omega)]^{2 \times 2}}.$$

Then, we have the following result providing the rate of convergence of (4.2) with  $\mathbf{H}_{0,h} = \tilde{\mathbf{H}}_{0,h}$ .

**THEOREM 4.2.** *Let  $(\boldsymbol{\sigma}, \mathbf{u}, \boldsymbol{\gamma}) \in \mathbf{H}_0$  and  $(\boldsymbol{\sigma}_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \tilde{\mathbf{H}}_{0,h}$  be the unique solutions of the continuous and discrete augmented formulations (3.7) and (4.2), respectively. Assume that  $\boldsymbol{\sigma} \in [H^r(\Omega)]^{2 \times 2}$ ,  $\mathbf{div}(\boldsymbol{\sigma}) \in [H^r(\Omega)]^2$ ,  $\mathbf{u} \in [H^{1+r}(\Omega)]^2$ , and  $\boldsymbol{\gamma} \in [H^r(\Omega)]^{2 \times 2}$ , for some  $r \in (0, 1]$ . Then there exists  $C > 0$ , independent of  $h$ , such that*

$$\begin{aligned} & \| (\boldsymbol{\sigma}, \mathbf{u}, \boldsymbol{\gamma}) - (\boldsymbol{\sigma}_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \|_{\mathbf{H}_0} \leq \\ & C h^r \left\{ \| \boldsymbol{\sigma} \|_{[H^r(\Omega)]^{2 \times 2}} + \| \mathbf{div}(\boldsymbol{\sigma}) \|_{[H^r(\Omega)]^2} + \| \mathbf{u} \|_{[H^{1+r}(\Omega)]^2} + \| \boldsymbol{\gamma} \|_{[H^r(\Omega)]^{2 \times 2}} \right\}. \end{aligned}$$

*Proof.* It follows from the Cea estimate (4.3) and the approximation properties (AP $_{0,h}^\sigma$ ), (AP $_{0,h}^{\mathbf{u}}$ ), and (AP $_h^\gamma$ ).  $\square$

We now go back to the general situation and state the discrete analogue of Lemma 3.1, which gives a sufficient condition for the equivalence between (4.1) and (4.2).

**LEMMA 4.3.** *Assume that the pressure finite element subspace  $H_h^p$  contains the traces of the members of the stress tensor finite element subspace  $H_{0,h}^\sigma$ , that is,*

$$\operatorname{tr}(H_{0,h}^\sigma) \subseteq H_h^p, \quad (4.8)$$

*Then, problems (4.1) and (4.2) are equivalent:  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_h$  is a solution of (4.1) if and only if  $(\boldsymbol{\sigma}_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_{0,h}$  is a solution of (4.2) and  $p_h = -\frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma}_h)$ .*

*Proof.* Let  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_h$  be a solution of (4.1). It is clear from (4.8) that  $p_h + \frac{1}{2}\text{tr}(\boldsymbol{\sigma}_h)$  belongs to  $H_h^p$ . Then, taking  $(\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h) = (\mathbf{0}, p_h + \frac{1}{2}\text{tr}(\boldsymbol{\sigma}_h), \mathbf{0}, \mathbf{0}) \in \mathbf{H}_h$ , we find from (4.1) that

$$\frac{\kappa_0}{\mu} \int_{\Omega} \left( p_h + \frac{1}{2}\text{tr}(\boldsymbol{\sigma}_h) \right)^2 = 0,$$

which yields  $p_h = -\frac{1}{2}\text{tr}(\boldsymbol{\sigma}_h)$ . Conversely, given  $(\boldsymbol{\sigma}_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_{0,h}$  a solution of (4.2), we let  $p_h := -\frac{1}{2}\text{tr}(\boldsymbol{\sigma}_h)$  and see that  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_h$  becomes a solution of (4.1).  $\square$

A particular example of finite element subspaces satisfying (4.8) is given by (cf. (4.5))

$$H_{0,h}^{\boldsymbol{\sigma}} := \tilde{H}_{0,h}^{\boldsymbol{\sigma}} \quad \text{and} \quad H_h^p := \{q_h \in L^2(\Omega) : q_h|_T \in \mathbb{P}_1(T) \quad \forall T \in \mathcal{T}_h\}.$$

Anyway, it becomes clear from Lemma 4.3 that the augmented scheme (4.1) makes sense only for pressure finite element subspaces not satisfying the condition (4.8). According to the above, we now aim to show that (4.1) is well-posed when an arbitrary finite element subspace  $\mathbf{H}_h$  of  $\mathbf{H}$  is considered. The idea, similarly as for  $A_0$ , is to choose  $\kappa_0, \kappa_1, \kappa_2$ , and  $\kappa_3$  such that  $A$  be strongly coercive on  $\mathbf{H}$  with respect to the norm  $\|\cdot\|_{\mathbf{H}}$  defined by

$$\|(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})\|_{\mathbf{H}} := \left\{ \|\boldsymbol{\tau}\|_{H(\text{div};\Omega)}^2 + \|q\|_{L^2(\Omega)}^2 + |\mathbf{v}|_{[H^1(\Omega)]^2}^2 + \|\boldsymbol{\eta}\|_{[L^2(\Omega)]^{2 \times 2}}^2 \right\}^{1/2}.$$

In fact, we first notice that

$$\begin{aligned} A((\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})) &= \frac{1}{2\mu} \left( 1 - \frac{\kappa_1}{2\mu} \right) \|\boldsymbol{\tau}^d\|_{[L^2(\Omega)]^{2 \times 2}}^2 + \frac{\kappa_0}{\mu} \left\| q + \frac{1}{2}\text{tr}(\boldsymbol{\tau}) \right\|_{L^2(\Omega)}^2 \\ &+ \kappa_2 \|\mathbf{div}(\boldsymbol{\tau})\|_{[L^2(\Omega)]^2}^2 + (\kappa_1 + \kappa_3) \|\boldsymbol{\varepsilon}(\mathbf{v})\|_{[L^2(\Omega)]^{2 \times 2}}^2 - \kappa_3 |\mathbf{v}|_{[H^1(\Omega)]^2}^2 + \kappa_3 \|\boldsymbol{\eta}\|_{[L^2(\Omega)]^{2 \times 2}}^2, \end{aligned}$$

which, using again Korn's first inequality, employing the estimate

$$\left\| q + \frac{1}{2}\text{tr}(\boldsymbol{\tau}) \right\|_{L^2(\Omega)}^2 \geq \frac{1}{2} \|q\|_{L^2(\Omega)}^2 - \left\| \frac{1}{2}\text{tr}(\boldsymbol{\tau}) \right\|_{L^2(\Omega)}^2 \geq \frac{1}{2} \|q\|_{L^2(\Omega)}^2 - \frac{1}{2} \|\boldsymbol{\tau}\|_{[L^2(\Omega)]^{2 \times 2}}^2,$$

and taking  $\kappa_0 > 0$ , yields

$$\begin{aligned} A((\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})) &\geq \frac{1}{2\mu} \left( 1 - \frac{\kappa_1}{2\mu} \right) \|\boldsymbol{\tau}^d\|_{[L^2(\Omega)]^{2 \times 2}}^2 - \frac{\kappa_0}{2\mu} \|\boldsymbol{\tau}\|_{[L^2(\Omega)]^{2 \times 2}}^2 \\ &+ \kappa_2 \|\mathbf{div}(\boldsymbol{\tau})\|_{[L^2(\Omega)]^2}^2 + \frac{\kappa_0}{2\mu} \|q\|_{L^2(\Omega)}^2 + \frac{(\kappa_1 - \kappa_3)}{2} |\mathbf{v}|_{[H^1(\Omega)]^2}^2 + \kappa_3 \|\boldsymbol{\eta}\|_{[L^2(\Omega)]^{2 \times 2}}^2. \end{aligned}$$

Then, choosing  $\kappa_1$  and  $\kappa_2$  such that

$$0 < \kappa_1 < 2\mu \quad \text{and} \quad 0 < \kappa_2,$$

and applying Lemma 2.4, we deduce that

$$\begin{aligned} A((\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})) &\geq \left( c_1 \alpha_1 - \frac{\kappa_0}{2\mu} \right) \|\boldsymbol{\tau}\|_{[L^2(\Omega)]^{2 \times 2}}^2 + \frac{\kappa_2}{2} \|\mathbf{div}(\boldsymbol{\tau})\|_{[L^2(\Omega)]^2}^2 \\ &+ \frac{\kappa_0}{2\mu} \|q\|_{L^2(\Omega)}^2 + \frac{(\kappa_1 - \kappa_3)}{2} |\mathbf{v}|_{[H^1(\Omega)]^2}^2 + \kappa_3 \|\boldsymbol{\eta}\|_{[L^2(\Omega)]^{2 \times 2}}^2, \end{aligned}$$

where  $c_1$  is the constant from Lemma 2.4 and

$$\alpha_1 := \min \left\{ \frac{1}{2\mu} \left( 1 - \frac{\kappa_1}{2\mu} \right), \frac{\kappa_2}{2} \right\}.$$

Hence, choosing the parameters  $\kappa_0$  and  $\kappa_3$  such that

$$0 < \kappa_0 < 2\mu c_1 \alpha_1 \quad \text{and} \quad 0 < \kappa_3 < \kappa_1,$$

we find that

$$A((\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})) \geq \alpha \|(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})\|_{\mathbf{H}}^2 \quad \forall (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}) \in \mathbf{H}, \quad (4.9)$$

where

$$\alpha := \min \left\{ \alpha_2, \frac{\kappa_0}{2\mu}, \frac{(\kappa_1 - \kappa_3)}{2}, \kappa_3 \right\} \quad \text{and} \quad \alpha_2 := \min \left\{ c_1 \alpha_1 - \frac{\kappa_0}{2\mu}, \frac{\kappa_2}{2} \right\}.$$

We are now in a position to establish the following result.

**THEOREM 4.4.** *Assume that there hold*

$$0 < \kappa_0 < 2\mu c_1 \alpha_1, \quad 0 < \kappa_1 < 2\mu, \quad 0 < \kappa_2, \quad \text{and} \quad 0 < \kappa_3 < \kappa_1.$$

*In addition, let  $\mathbf{H}_h$  be any finite element subspace of  $\mathbf{H}$ . Then, the Galerkin scheme (4.1) has a unique solution  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_h$ , and there exist positive constants  $C, \tilde{C}$ , independent of  $h$ , such that*

$$\|(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h)\|_{\mathbf{H}} \leq C \sup_{\substack{(\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \in \mathbf{H}_h \\ (\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \neq 0}} \frac{|F(\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h)|}{\|(\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h)\|_{\mathbf{H}}} \leq C \|\mathbf{f}\|_{[L^2(\Omega)]^2},$$

and

$$\|(\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma}) - (\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h)\|_{\mathbf{H}} \leq \tilde{C} \inf_{(\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \in \mathbf{H}_h} \|(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) - (\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h)\|_{\mathbf{H}}.$$

*Proof.* Since  $A$  is bounded and strongly coercive on  $\mathbf{H}$  (cf. (3.5) and (4.9)) with constants depending on  $\mu$  and  $(\kappa_0, \kappa_1, \kappa_2, \kappa_3)$ , the proof follows from a straightforward application of the Lax-Milgram Lemma, and Cea's estimate.  $\square$

In order to consider an explicit Galerkin scheme (4.1), we now let

$$\tilde{H}_h^p := \{q_h \in L^2(\Omega) : q_h|_T \in \mathbb{P}_0(T) \quad \forall T \in \mathcal{T}_h\},$$

and define

$$\tilde{\mathbf{H}}_h := \tilde{H}_{0,h}^{\boldsymbol{\sigma}} \times \tilde{H}_h^p \times \tilde{H}_{0,h}^{\mathbf{u}} \times \tilde{H}_h^{\boldsymbol{\gamma}}, \quad (4.10)$$

where  $\tilde{H}_{0,h}^{\boldsymbol{\sigma}}$ ,  $\tilde{H}_{0,h}^{\mathbf{u}}$ , and  $\tilde{H}_h^{\boldsymbol{\gamma}}$  are given, respectively, by (4.5), (4.6), and (4.7).

The approximation property of  $\tilde{H}_h^p$  is given as follows (see [7], [11]):

(AP $_h^p$ ) *For each  $r \in [0, 1]$  and for each  $q \in H^r(\Omega)$  there exists  $q_h \in \tilde{H}_h^p$  such that*

$$\|q - q_h\|_{L^2(\Omega)} \leq C h^r \|q\|_{H^r(\Omega)}.$$

Then, we have the following theorem providing the rate of convergence of (4.1) with  $\mathbf{H}_h = \tilde{\mathbf{H}}_h$ .

**THEOREM 4.5.** *Let  $(\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma}) \in \mathbf{H}$  and  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \tilde{\mathbf{H}}_h$  be the unique solutions of the continuous and discrete augmented formulations (3.4) and (4.1), respectively. Assume that  $\boldsymbol{\sigma} \in [H^r(\Omega)]^{2 \times 2}$ ,  $\mathbf{div}(\boldsymbol{\sigma}) \in [H^r(\Omega)]^2$ ,  $\mathbf{u} \in [H^{r+1}(\Omega)]^2$ , and  $\boldsymbol{\gamma} \in [H^r(\Omega)]^{2 \times 2}$ , for some  $r \in (0, 1]$ . Then there exists  $C > 0$ , independent of  $h$ , such that*

$$\|(\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma}) - (\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h)\|_{\mathbf{H}} \leq C h^r \left\{ \|\boldsymbol{\sigma}\|_{[H^r(\Omega)]^{2 \times 2}} + \|\mathbf{div}(\boldsymbol{\sigma})\|_{[H^r(\Omega)]^2} + \|\mathbf{u}\|_{[H^{r+1}(\Omega)]^2} + \|\boldsymbol{\gamma}\|_{[H^r(\Omega)]^{2 \times 2}} \right\}.$$

*Proof.* We first notice, according to Theorem 3.3 and the hypothesis on  $\boldsymbol{\sigma}$ , that  $p = -\frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma})$  belongs to  $H^r(\Omega)$  and that  $\|p\|_{H^r(\Omega)} \leq C \|\boldsymbol{\sigma}\|_{[H^r(\Omega)]^{2 \times 2}}$ . Then, the proof follows from the Cea estimate from Theorem 4.4 and the approximation properties  $(\text{AP}_{0,h}^\boldsymbol{\sigma})$ ,  $(\text{AP}_h^p)$ ,  $(\text{AP}_{0,h}^{\mathbf{u}})$ , and  $(\text{AP}_h^\gamma)$ .  $\square$

At this point we would like to emphasize the main features of our augmented Galerkin schemes (4.1) and (4.2), as compared to each other, besides the fact that both of them can be implemented with any finite element subspace of  $\mathbf{H}$  and  $\mathbf{H}_0$ , respectively. In fact, it is important to notice on one hand that (4.2) allows an explicit and simple definition of the whole vector of parameters  $(\kappa_1, \kappa_2, \kappa_3)$  (cf. Theorem 3.2), whereas the choice of  $\kappa_0$  in (4.1) depends on the unknown constant  $c_1$  from Lemma 2.4. On the other hand, it is clear that (4.1) provides more flexibility for approximating the pressure since the corresponding finite element subspace  $H_h^p$  can be chosen arbitrarily, whereas (4.2) needs a postprocess to compute  $p_h$  in terms of  $\boldsymbol{\sigma}_h$ , either simply as  $p_h := -\frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma}_h)$  or projecting  $-\frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma}_h)$  onto some finite element subspace.

We end this section by mentioning that a useful discussion on the actual implementation of augmented Galerkin schemes of the present kind can be seen in [18].

## 5. A RESIDUAL BASED A POSTERIORI ERROR ESTIMATOR

In this section we derive a residual based a posteriori error estimator for (4.1), much in the spirit of [5]. The analysis for (4.2) is contained in what follows, and hence we omit details.

First we introduce several notations. Given  $T \in \mathcal{T}_h$ , we let  $\mathcal{E}(T)$  be the set of its edges, and let  $\mathcal{E}_h$  be the set of all edges of the triangulation  $\mathcal{T}_h$ . Then we write  $\mathcal{E}_h = \mathcal{E}_{h,\Omega} \cup \mathcal{E}_{h,\Gamma}$ , where  $\mathcal{E}_{h,\Omega} := \{e \in \mathcal{E}_h : e \subseteq \Omega\}$  and  $\mathcal{E}_{h,\Gamma} := \{e \in \mathcal{E}_h : e \subseteq \Gamma\}$ . In what follows,  $h_e$  stands for the length of the edge  $e$ . Further, given  $\boldsymbol{\tau} \in [L^2(\Omega)]^{2 \times 2}$  such that  $\boldsymbol{\tau}|_T \in C(T)$  on each  $T \in \mathcal{T}_h$ , an edge  $e \in \mathcal{E}(T) \cap \mathcal{E}_{h,\Omega}$ , and the unit tangential vector  $\mathbf{t}_T$  along  $e$ , we let  $J[\boldsymbol{\tau}\mathbf{t}_T]$  be the corresponding jump across  $e$ , that is,  $J[\boldsymbol{\tau}\mathbf{t}_T] := (\boldsymbol{\tau}|_T - \boldsymbol{\tau}|_{T'})|_e \mathbf{t}_T$ , where  $T'$  is the other triangle of  $\mathcal{T}_h$  having  $e$  as an edge. Abusing notation, when  $e \in \mathcal{E}_{h,\Gamma}$ , we also write  $J[\boldsymbol{\tau}\mathbf{t}_T] := \boldsymbol{\tau}|_e \mathbf{t}_T$ . We recall here that  $\mathbf{t}_T := (-\nu_2, \nu_1)^\mathbf{t}$ , where  $\boldsymbol{\nu}_T := (\nu_1, \nu_2)^\mathbf{t}$  is the unit outward vector normal to  $\partial T$ . Analogously, we define the normal jumps  $J[\boldsymbol{\tau}\boldsymbol{\nu}_T]$ . In addition, given scalar, vector and tensor valued fields  $v$ ,  $\boldsymbol{\varphi} := (\varphi_1, \varphi_2)$  and  $\boldsymbol{\tau} := (\tau_{ij})$ , respectively, we let

$$\operatorname{curl}(v) := \begin{pmatrix} -\frac{\partial v}{\partial x_2} \\ \frac{\partial v}{\partial x_1} \end{pmatrix}, \quad \underline{\operatorname{curl}}(\boldsymbol{\varphi}) := \begin{pmatrix} \operatorname{curl}(\varphi_1)^\mathbf{t} \\ \operatorname{curl}(\varphi_2)^\mathbf{t} \end{pmatrix}, \quad \text{and} \quad \operatorname{curl}(\boldsymbol{\tau}) := \begin{pmatrix} \frac{\partial \tau_{12}}{\partial x_1} - \frac{\partial \tau_{11}}{\partial x_2} \\ \frac{\partial \tau_{22}}{\partial x_1} - \frac{\partial \tau_{21}}{\partial x_2} \end{pmatrix}.$$

Then, letting  $(\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma}) \in \mathbf{H}$  and  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_h$  be the unique solutions of the continuous and discrete augmented formulations (3.4) and (4.1), respectively, we define for each  $T \in \mathcal{T}_h$  a local error indicator  $\theta_T$  as follows:

$$\begin{aligned} \theta_T^2 &:= \|\mathbf{f} + \operatorname{div}(\boldsymbol{\sigma}_h)\|_{[L^2(T)]^2}^2 + \|\boldsymbol{\sigma}_h - \boldsymbol{\sigma}_h^\mathbf{t}\|_{[L^2(T)]^{2 \times 2}}^2 + \|\boldsymbol{\gamma}_h - \frac{1}{2}(\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^\mathbf{t})\|_{[L^2(T)]^{2 \times 2}}^2 \\ &+ h_T^2 \|\operatorname{curl}\left(\frac{1}{2\mu} \boldsymbol{\sigma}_h^\mathbf{d} + \boldsymbol{\gamma}_h\right)\|_{[L^2(T)]^2}^2 + h_T^2 \|\underline{\operatorname{curl}}\left(p_h + \frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma}_h)\right)\|_{[L^2(T)]^2}^2 \\ &+ h_T^2 \|\operatorname{curl}\left(\boldsymbol{\varepsilon}(\mathbf{u}_h)^\mathbf{d} - \frac{1}{2\mu} \boldsymbol{\sigma}_h^\mathbf{d}\right)\|_{[L^2(T)]^2}^2 \\ &+ \sum_{e \in \mathcal{E}(T)} h_e \|J\left[\left(\frac{1}{2\mu} \boldsymbol{\sigma}_h^\mathbf{d} - \nabla \mathbf{u}_h + \boldsymbol{\gamma}_h\right) \mathbf{t}_T\right]\|_{[L^2(e)]^2}^2 \\ &+ \sum_{e \in \mathcal{E}(T)} h_e \|J\left[\left(p_h + \frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma}_h)\right) \mathbf{t}_T\right]\|_{[L^2(e)]^2}^2 + \sum_{e \in \mathcal{E}(T)} h_e \|J\left[\left(\boldsymbol{\varepsilon}(\mathbf{u}_h)^\mathbf{d} - \frac{1}{2\mu} \boldsymbol{\sigma}_h^\mathbf{d}\right) \mathbf{t}_T\right]\|_{[L^2(e)]^2}^2 \end{aligned}$$

$$\begin{aligned}
& + h_T^2 \|\mathbf{div}\left(\boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2}(\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^\dagger)^\mathbf{d}\right)\|_{[L^2(T)]^2}^2 \\
& + h_T^2 \|\mathbf{div}\left(\boldsymbol{\gamma}_h - \frac{1}{2}(\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^\dagger)\right)\|_{[L^2(T)]^2}^2 \\
& + \sum_{e \in \mathcal{E}(T) \cap \mathcal{E}_{h,\Omega}} h_e \|J\left[\left(\boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2}(\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^\dagger)^\mathbf{d}\right) \boldsymbol{\nu}_T\right]\|_{[L^2(e)]^2}^2 \\
& + \sum_{e \in \mathcal{E}(T) \cap \mathcal{E}_{h,\Omega}} h_e \|J\left[\left(\boldsymbol{\gamma}_h - \frac{1}{2}(\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^\dagger)\right) \boldsymbol{\nu}_T\right]\|_{[L^2(e)]^2}^2.
\end{aligned} \tag{5.1}$$

The residual character of each term on the right hand side of (5.1) is quite clear. As usual the expression  $\boldsymbol{\theta} := \left\{ \sum_{T \in \mathcal{T}_h} \theta_T^2 \right\}^{1/2}$  is employed as the global residual error estimator.

The following theorem is the main result of this section.

**THEOREM 5.1.** *Let  $(\boldsymbol{\sigma}, p, \mathbf{u}, \boldsymbol{\gamma}) \in \mathbf{H}$  and  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \mathbf{H}_h$  be the unique solutions of (3.4) and (4.1), respectively. Then there exist positive constants  $C_{\text{eff}}$  and  $C_{\text{rel}}$ , independent of  $h$ , such that*

$$C_{\text{eff}} \boldsymbol{\theta} \leq \|(\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h)\|_{\mathbf{H}} \leq C_{\text{rel}} \boldsymbol{\theta}. \tag{5.2}$$

The efficiency of the global error estimator (lower bound in (5.2)) is proved below in Subsection 5.2 and the reliability of the global error estimator (upper bound in (5.2)) is derived now.

**5.1. Reliability.** We begin with the following preliminary estimate.

**LEMMA 5.2.** *There exists  $C > 0$ , independent of  $h$ , such that*

$$\begin{aligned}
& C \|(\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h)\|_{\mathbf{H}} \leq \\
& \sup_{\substack{(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}) \in \mathbf{H}\{\mathbf{0}\} \\ \mathbf{div}(\boldsymbol{\tau}) = \mathbf{0}}} \frac{A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}))}{\|(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})\|_{\mathbf{H}}} + \|\mathbf{f} + \mathbf{div}(\boldsymbol{\sigma}_h)\|_{[L^2(\Omega)]^2}
\end{aligned} \tag{5.3}$$

*Proof.* Let us define  $\boldsymbol{\sigma}^* = \boldsymbol{\varepsilon}(\mathbf{z})$ , where  $\mathbf{z} \in [H_0^1(\Omega)]^2$  is the unique solution of the boundary value problem:  $-\mathbf{div}(\boldsymbol{\varepsilon}(\mathbf{z})) = \mathbf{f} + \mathbf{div}(\boldsymbol{\sigma}_h)$  in  $\Omega$ ,  $\mathbf{z} = \mathbf{0}$  on  $\Gamma$ . It follows that  $\boldsymbol{\sigma}^* \in H_0$  and the corresponding continuous dependence result establishes the existence of  $c > 0$  such that

$$\|\boldsymbol{\sigma}^*\|_{H(\mathbf{div}; \Omega)} \leq c \|\mathbf{f} + \mathbf{div}(\boldsymbol{\sigma}_h)\|_{[L^2(\Omega)]^2}. \tag{5.4}$$

In addition,  $\mathbf{div}(\boldsymbol{\sigma} - \boldsymbol{\sigma}_h - \boldsymbol{\sigma}^*) = -\mathbf{f} - \mathbf{div}(\boldsymbol{\sigma}_h) + (\mathbf{f} + \mathbf{div}(\boldsymbol{\sigma}_h)) = \mathbf{0}$  in  $\Omega$ . Let  $\alpha$  and  $M$  be the coercivity and boundedness constants of  $A$ . Then, using the coercivity of  $A$  we find that

$$\begin{aligned}
& \alpha \|(\boldsymbol{\sigma} - \boldsymbol{\sigma}_h - \boldsymbol{\sigma}^*, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h)\|_{\mathbf{H}}^2 \\
& \leq A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h - \boldsymbol{\sigma}^*, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\boldsymbol{\sigma} - \boldsymbol{\sigma}_h - \boldsymbol{\sigma}^*, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h)) \\
& \leq A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\boldsymbol{\sigma} - \boldsymbol{\sigma}_h - \boldsymbol{\sigma}^*, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h)) \\
& \quad - A((\boldsymbol{\sigma}^*, 0, \mathbf{0}, \mathbf{0}), (\boldsymbol{\sigma} - \boldsymbol{\sigma}_h - \boldsymbol{\sigma}^*, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h)),
\end{aligned}$$

which, employing the boundedness of  $A$ , yields

$$\begin{aligned}
& \alpha \|(\boldsymbol{\sigma} - \boldsymbol{\sigma}_h - \boldsymbol{\sigma}^*, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h)\|_{\mathbf{H}} \\
& \leq \sup_{\substack{(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}) \in \mathbf{H}\{\mathbf{0}\} \\ \mathbf{div}(\boldsymbol{\tau}) = \mathbf{0}}} \frac{A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}))}{\|(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})\|_{\mathbf{H}}} + M \|\boldsymbol{\sigma}^*\|_{H(\mathbf{div}; \Omega)}.
\end{aligned} \tag{5.5}$$

Hence, (5.3) follows straightforwardly from the triangle inequality, (5.4) and (5.5).  $\square$

It remains to bound the first term on the right hand side of (5.3). To this end, we will make use of the well known Clément interpolation operator,  $I_h : H^1(\Omega) \rightarrow X_h$  (cf. [12]), with  $X_h$  given by

$$X_h := \{v_h \in C(\bar{\Omega}) : v_h|_T \in \mathbb{P}_1(T) \quad \forall T \in \mathcal{T}_h\},$$

which satisfies the standard local approximation properties stated below in Lemma 5.3. It is important to remark that  $I_h$  is defined in [12] so that  $I_h(v) \in X_h \cap H_0^1(\Omega)$  for all  $v \in H_0^1(\Omega)$ .

LEMMA 5.3. *There exist constants  $C_1, C_2 > 0$ , independent of  $h$ , such that for all  $v \in H^1(\Omega)$  there holds*

$$\|v - I_h(v)\|_{L^2(T)} \leq C_1 h_T \|v\|_{H^1(\tilde{\omega}_T)} \quad \forall T \in \mathcal{T}_h,$$

and

$$\|v - I_h(v)\|_{L^2(e)} \leq C_2 h_e^{1/2} \|v\|_{H^1(\tilde{\omega}_e)} \quad \forall e \in \mathcal{E}_h,$$

where  $\tilde{\omega}_T$  and  $\tilde{\omega}_e$  are the union of all elements sharing at least one point with  $T$  and  $e$ , respectively.

*Proof.* See [12].  $\square$

We now let  $(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}) \in \mathbf{H}$ ,  $(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}) \neq \mathbf{0}$ , such that  $\operatorname{div}(\boldsymbol{\tau}) = \mathbf{0}$  in  $\Omega$ . Since  $\Omega$  is connected, there exists a stream function  $\boldsymbol{\varphi} := (\varphi_1, \varphi_2) \in [H^1(\Omega)]^2$  such that  $\int_{\Omega} \varphi_1 = \int_{\Omega} \varphi_2 = 0$  and  $\boldsymbol{\tau} = \underline{\operatorname{curl}}(\boldsymbol{\varphi})$ . Then, denoting  $\boldsymbol{\varphi}_h := (I_h(\varphi_1), I_h(\varphi_2))$ , we define  $\boldsymbol{\tau}_h := \underline{\operatorname{curl}}(\boldsymbol{\varphi}_h)$ .

It can be seen that, since  $\boldsymbol{\tau}_h$  has  $[H^1(T)]^{2 \times 2}$ -regularity on each triangle (in fact, it is piecewise constant), and its rows have continuous normal components across each interior edge,  $\boldsymbol{\tau}_h$  has a  $L^2(\Omega)$  divergence, which is zero. Thus,  $\boldsymbol{\tau}_h$  belongs to  $\tilde{H}_h^{\boldsymbol{\sigma}}$  (cf. (4.4)). The decomposition  $\boldsymbol{\tau}_h = \boldsymbol{\tau}_{h,0} + d_h \mathbf{I}$ , holds, where  $\boldsymbol{\tau}_{h,0} \in \tilde{H}_{0,h}^{\boldsymbol{\sigma}}$  (cf. (4.5)) and  $d_h = \frac{\int_{\Omega} \operatorname{tr}(\boldsymbol{\tau}_h)}{2|\Omega|} \in \mathbb{R}$ .

We also define  $\mathbf{v}_h := (I_h(v_1), I_h(v_2)) \in H_0^1$ , the vector Clément interpolant of  $\mathbf{v} := (v_1, v_2) \in [H_0^1(\Omega)]^2$ . From the Galerkin orthogonality, it follows that

$$\begin{aligned} A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})) &= \\ &= A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\boldsymbol{\tau} - \boldsymbol{\tau}_{h,0}, q, \mathbf{v} - \mathbf{v}_h, \boldsymbol{\eta})). \end{aligned} \quad (5.6)$$

Also, from (3.5), the orthogonality between symmetric and asymmetric tensors, and as a consequence, again, of the Galerkin orthogonality, it follows that

$$\begin{aligned} &A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (d_h \mathbf{I}, 0, \mathbf{0}, \mathbf{0})) \\ &= \frac{\kappa_0}{\mu} \int_{\Omega} \left( p - p_h + \frac{1}{2} \operatorname{tr}(\boldsymbol{\sigma} - \boldsymbol{\sigma}_h) \right) \frac{1}{2} \operatorname{tr}(d_h \mathbf{I}) \\ &= A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\mathbf{0}, d_h, \mathbf{0}, \mathbf{0})) \\ &= 0. \end{aligned} \quad (5.7)$$

Hence, (5.6), (5.7) and (4.1) give

$$\begin{aligned} &A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})) \\ &= A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\boldsymbol{\tau} - \boldsymbol{\tau}_h, q, \mathbf{v} - \mathbf{v}_h, \boldsymbol{\eta})) \\ &= F(\boldsymbol{\tau} - \boldsymbol{\tau}_h, q, \mathbf{v} - \mathbf{v}_h, \boldsymbol{\eta}) - A((\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h), (\boldsymbol{\tau} - \boldsymbol{\tau}_h, q, \mathbf{v} - \mathbf{v}_h, \boldsymbol{\eta})), \end{aligned} \quad (5.8)$$

which, after some algebraic manipulations, yields that

$$\begin{aligned}
& A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})) \\
&= \int_{\Omega} (\mathbf{f} + \mathbf{div}(\boldsymbol{\sigma}_h)) \cdot (\mathbf{v} - \mathbf{v}_h) + \int_{\Omega} \left( \frac{1}{2}(\boldsymbol{\sigma}_h - \boldsymbol{\sigma}_h^{\dagger}) - \kappa_3(\boldsymbol{\gamma}_h - \frac{1}{2}(\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\dagger})) \right) : \boldsymbol{\eta} \\
&- \int_{\Omega} \left\{ \kappa_1 \left( \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2}(\boldsymbol{\sigma}_h^{\text{d}} + (\boldsymbol{\sigma}_h^{\text{d}})^{\dagger}) \right) + \kappa_3(\boldsymbol{\gamma}_h - \frac{1}{2}(\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\dagger})) \right\} : \nabla(\mathbf{v} - \mathbf{v}_h) \\
&- \int_{\Omega} \left\{ \left( \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} - \nabla \mathbf{u}_h + \boldsymbol{\gamma}_h \right) + \frac{\kappa_0}{2\mu} \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) \mathbf{I} + \frac{\kappa_1}{2\mu} \left( \boldsymbol{\varepsilon}(\mathbf{u}_h)^{\text{d}} - \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} \right) \right\} : (\boldsymbol{\tau} - \boldsymbol{\tau}_h) \\
&- \frac{\kappa_0}{\mu} \int_{\Omega} \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) p.
\end{aligned} \tag{5.9}$$

The rest of reliability consists in deriving suitable upper bounds for each one of the terms appearing on the right hand side of (5.9). We begin by noticing that direct applications of the Cauchy-Schwarz inequality give

$$\left| \int_{\Omega} \frac{1}{2}(\boldsymbol{\sigma}_h - \boldsymbol{\sigma}_h^{\dagger}) : \boldsymbol{\eta} \right| \leq \| \boldsymbol{\sigma}_h - \boldsymbol{\sigma}_h^{\dagger} \|_{[L^2(\Omega)]^{2 \times 2}} \| \boldsymbol{\eta} \|_{[L^2(\Omega)]^{2 \times 2}}, \tag{5.10}$$

$$\left| \int_{\Omega} \left( \boldsymbol{\gamma}_h - \frac{1}{2}(\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\dagger}) \right) : \boldsymbol{\eta} \right| \leq \left\| \boldsymbol{\gamma}_h - \frac{1}{2}(\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\dagger}) \right\|_{[L^2(\Omega)]^{2 \times 2}} \| \boldsymbol{\eta} \|_{[L^2(\Omega)]^{2 \times 2}}, \tag{5.11}$$

and

$$\left| \int_{\Omega} \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) p \right| \leq \left\| p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right\|_{L^2(\Omega)} \| p \|_{L^2(\Omega)}. \tag{5.12}$$

The decomposition  $\bar{\Omega} = \bigcup_{T \in \mathcal{T}_h} T$  and the use of integration by parts formulae on each element are employed next to handle the terms from the third and the fourth rows of (5.9). We first replace  $\boldsymbol{\tau} - \boldsymbol{\tau}_h$  by  $\underline{\mathbf{curl}}(\boldsymbol{\varphi} - \boldsymbol{\varphi}_h)$  and use that  $\text{curl}(\nabla \mathbf{u}_h) = \mathbf{0}$  on each triangle  $T \in \mathcal{T}_h$ , to obtain

$$\begin{aligned}
\int_{\Omega} \left( \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} - \nabla \mathbf{u}_h + \boldsymbol{\gamma}_h \right) : (\boldsymbol{\tau} - \boldsymbol{\tau}_h) &= \sum_{T \in \mathcal{T}_h} \int_T \left( \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} - \nabla \mathbf{u}_h + \boldsymbol{\gamma}_h \right) : \underline{\mathbf{curl}}(\boldsymbol{\varphi} - \boldsymbol{\varphi}_h) \\
&= \sum_{T \in \mathcal{T}_h} \int_T \text{curl} \left( \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} + \boldsymbol{\gamma}_h \right) \cdot (\boldsymbol{\varphi} - \boldsymbol{\varphi}_h) \\
&- \sum_{e \in \mathcal{E}_h} \left\langle J \left[ \left( \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} - \nabla \mathbf{u}_h + \boldsymbol{\gamma}_h \right) \mathbf{t}_T \right], \boldsymbol{\varphi} - \boldsymbol{\varphi}_h \right\rangle_{[L^2(e)]^2}, \tag{5.13}
\end{aligned}$$

$$\begin{aligned}
\int_{\Omega} \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) \mathbf{I} : (\boldsymbol{\tau} - \boldsymbol{\tau}_h) &= \sum_{T \in \mathcal{T}_h} \int_T \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) \mathbf{I} : \underline{\mathbf{curl}}(\boldsymbol{\varphi} - \boldsymbol{\varphi}_h) \\
&= \sum_{T \in \mathcal{T}_h} \int_T \text{curl} \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) \cdot (\boldsymbol{\varphi} - \boldsymbol{\varphi}_h) \\
&- \sum_{e \in \mathcal{E}_h} \left\langle J \left[ \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) \mathbf{t}_T \right], \boldsymbol{\varphi} - \boldsymbol{\varphi}_h \right\rangle_{[L^2(e)]^2}, \tag{5.14}
\end{aligned}$$

and

$$\begin{aligned}
\int_{\Omega} \left( \boldsymbol{\varepsilon}(\mathbf{u}_h)^{\text{d}} - \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} \right) : (\boldsymbol{\tau} - \boldsymbol{\tau}_h) &= \sum_{T \in \mathcal{T}_h} \int_T \left( \boldsymbol{\varepsilon}(\mathbf{u}_h)^{\text{d}} - \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} \right) : \underline{\text{curl}}(\boldsymbol{\varphi} - \boldsymbol{\varphi}_h) \\
&= \sum_{T \in \mathcal{T}_h} \int_T \text{curl} \left( \boldsymbol{\varepsilon}(\mathbf{u}_h)^{\text{d}} - \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} \right) \cdot (\boldsymbol{\varphi} - \boldsymbol{\varphi}_h) \\
&\quad - \sum_{e \in \mathcal{E}_h} \left\langle J \left[ \left( \boldsymbol{\varepsilon}(\mathbf{u}_h)^{\text{d}} - \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} \right) \mathbf{t}_T \right], \boldsymbol{\varphi} - \boldsymbol{\varphi}_h \right\rangle_{[L^2(e)]^2}. \quad (5.15)
\end{aligned}$$

On the other hand, using that  $\mathbf{v} - \mathbf{v}_h = \mathbf{0}$  on  $\Gamma$ , we get

$$\begin{aligned}
\int_{\Omega} \left( \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^{\text{t}})^{\text{d}} \right) : \nabla(\mathbf{v} - \mathbf{v}_h) \\
= - \sum_{T \in \mathcal{T}_h} \int_T \mathbf{div} \left( \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^{\text{t}})^{\text{d}} \right) \cdot (\mathbf{v} - \mathbf{v}_h) \\
+ \sum_{e \in \mathcal{E}_{h,\Omega}} \left\langle J \left[ \left( \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^{\text{t}})^{\text{d}} \right) \boldsymbol{\nu}_T \right], \mathbf{v} - \mathbf{v}_h \right\rangle_{[L^2(e)]^2}, \quad (5.16)
\end{aligned}$$

and

$$\begin{aligned}
\int_{\Omega} \left( \boldsymbol{\gamma}_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\text{t}}) \right) : \nabla(\mathbf{v} - \mathbf{v}_h) \\
= - \sum_{T \in \mathcal{T}_h} \int_T \mathbf{div} \left( \boldsymbol{\gamma}_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\text{t}}) \right) \cdot (\mathbf{v} - \mathbf{v}_h) \\
+ \sum_{e \in \mathcal{E}_{h,\Omega}} \left\langle J \left[ \left( \boldsymbol{\gamma}_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\text{t}}) \right) \boldsymbol{\nu}_T \right], \mathbf{v} - \mathbf{v}_h \right\rangle_{[L^2(e)]^2}. \quad (5.17)
\end{aligned}$$

In what follows we apply again the Cauchy-Schwarz inequality, Lemma 5.3 and the fact that the number of triangles is bounded independently of  $h$  in both  $\tilde{\omega}_T$  and  $\tilde{\omega}_e$  to derive the estimates for the expression  $\int_{\Omega} (\mathbf{f} + \mathbf{div}(\boldsymbol{\sigma}_h)) \cdot (\mathbf{v} - \mathbf{v}_h)$  in (5.9) and the right hand sides of (5.13), (5.14), (5.15), (5.16), and (5.17), with constants  $C$  independent of  $h$ . Indeed, we easily have

$$\left| \int_{\Omega} (\mathbf{f} + \mathbf{div}(\boldsymbol{\sigma}_h)) \cdot (\mathbf{v} - \mathbf{v}_h) \right| \leq C \left\{ \sum_{T \in \mathcal{T}_h} h_T^2 \|\mathbf{f} + \mathbf{div}(\boldsymbol{\sigma}_h)\|_{[L^2(T)]^2}^2 \right\}^{1/2} \|\mathbf{v}\|_{[H^1(\Omega)]^2}. \quad (5.18)$$

In addition, for the terms containing the stream function  $\boldsymbol{\varphi}$  (cf. (5.13), (5.14), (5.15)), we get

$$\begin{aligned}
\left| \sum_{T \in \mathcal{T}_h} \int_T \text{curl} \left( \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} + \boldsymbol{\gamma}_h \right) \cdot (\boldsymbol{\varphi} - \boldsymbol{\varphi}_h) \right| \\
\leq C \left\{ \sum_{T \in \mathcal{T}_h} h_T^2 \left\| \text{curl} \left( \frac{1}{2\mu} \boldsymbol{\sigma}_h^{\text{d}} + \boldsymbol{\gamma}_h \right) \right\|_{[L^2(T)]^2}^2 \right\}^{1/2} \|\boldsymbol{\varphi}\|_{[H^1(\Omega)]^2}, \quad (5.19)
\end{aligned}$$



$$\begin{aligned}
& \left| \sum_{T \in \mathcal{T}_h} \int_T \mathbf{curl} \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) \cdot (\boldsymbol{\varphi} - \boldsymbol{\varphi}_h) \right| \\
& \leq C \left\{ \sum_{T \in \mathcal{T}_h} h_T^2 \left\| \mathbf{curl} \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) \right\|_{[L^2(T)]^2}^2 \right\}^{1/2} \|\boldsymbol{\varphi}\|_{[H^1(\Omega)]^2}, \quad (5.20)
\end{aligned}$$

$$\begin{aligned}
& \left| \sum_{T \in \mathcal{T}_h} \int_T \mathbf{curl} \left( \boldsymbol{\varepsilon}(\mathbf{u}_h)^d - \frac{1}{2\mu} \boldsymbol{\sigma}_h^d \right) \cdot (\boldsymbol{\varphi} - \boldsymbol{\varphi}_h) \right| \\
& \leq C \left\{ \sum_{T \in \mathcal{T}_h} h_T^2 \left\| \mathbf{curl} \left( \boldsymbol{\varepsilon}(\mathbf{u}_h)^d - \frac{1}{2\mu} \boldsymbol{\sigma}_h^d \right) \right\|_{[L^2(T)]^2}^2 \right\}^{1/2} \|\boldsymbol{\varphi}\|_{[H^1(\Omega)]^2}, \quad (5.21)
\end{aligned}$$

$$\begin{aligned}
& \left| \sum_{e \in \mathcal{E}_h} \left\langle J \left[ \left( \frac{1}{2\mu} \boldsymbol{\sigma}_h^d - \nabla \mathbf{u}_h + \boldsymbol{\gamma}_h \right) \mathbf{t}_T \right], \boldsymbol{\varphi} - \boldsymbol{\varphi}_h \right\rangle_{[L^2(e)]^2} \right| \\
& \leq C \left\{ \sum_{e \in \mathcal{E}_h} h_e \left\| J \left[ \left( \frac{1}{2\mu} \boldsymbol{\sigma}_h^d - \nabla \mathbf{u}_h + \boldsymbol{\gamma}_h \right) \mathbf{t}_T \right] \right\|_{[L^2(e)]^2}^2 \right\}^{1/2} \|\boldsymbol{\varphi}\|_{[H^1(\Omega)]^2}, \quad (5.22)
\end{aligned}$$

$$\begin{aligned}
& \left| \sum_{e \in \mathcal{E}_h} \left\langle J \left[ \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) \mathbf{t}_T \right], \boldsymbol{\varphi} - \boldsymbol{\varphi}_h \right\rangle_{[L^2(e)]^2} \right| \\
& \leq C \left\{ \sum_{e \in \mathcal{E}_h} h_e \left\| J \left[ \left( p_h + \frac{1}{2} \text{tr}(\boldsymbol{\sigma}_h) \right) \mathbf{t}_T \right] \right\|_{[L^2(e)]^2}^2 \right\}^{1/2} \|\boldsymbol{\varphi}\|_{[H^1(\Omega)]^2}, \quad (5.23)
\end{aligned}$$

and

$$\begin{aligned}
& \left| \sum_{e \in \mathcal{E}_h} \left\langle J \left[ \left( \boldsymbol{\varepsilon}(\mathbf{u}_h)^d - \frac{1}{2\mu} \boldsymbol{\sigma}_h^d \right) \mathbf{t}_T \right], \boldsymbol{\varphi} - \boldsymbol{\varphi}_h \right\rangle_{[L^2(e)]^2} \right| \\
& \leq C \left\{ \sum_{e \in \mathcal{E}_h} h_e \left\| J \left[ \left( \boldsymbol{\varepsilon}(\mathbf{u}_h)^d - \frac{1}{2\mu} \boldsymbol{\sigma}_h^d \right) \mathbf{t}_T \right] \right\|_{[L^2(e)]^2}^2 \right\}^{1/2} \|\boldsymbol{\varphi}\|_{[H^1(\Omega)]^2}. \quad (5.24)
\end{aligned}$$

We observe here, due to the equivalence between  $\|\boldsymbol{\varphi}\|_{[H^1(\Omega)]^2}$  and  $\|\nabla \boldsymbol{\varphi}\|_{[L^2(\Omega)]^{2 \times 2}}$ , that

$$\|\boldsymbol{\varphi}\|_{[H^1(\Omega)]^2} \leq C \|\nabla \boldsymbol{\varphi}\|_{[L^2(\Omega)]^{2 \times 2}} = C \|\mathbf{curl}(\boldsymbol{\varphi})\|_{[L^2(\Omega)]^{2 \times 2}} = C \|\boldsymbol{\tau}\|_{H(\text{div}; \Omega)},$$

which allows to replace  $\|\boldsymbol{\varphi}\|_{[H^1(\Omega)]^2}$  by  $\|\boldsymbol{\tau}\|_{H(\text{div}; \Omega)}$  in the above estimates (5.19) - (5.24).

Similarly, for the terms on the right hand side of (5.16) and (5.17), we find that

$$\begin{aligned} & \left| \sum_{T \in \mathcal{T}_h} \int_T \mathbf{div} \left( \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^{\dagger})^{\mathbf{d}} \right) \cdot (\mathbf{v} - \mathbf{v}_h) \right| \\ & \leq C \left\{ \sum_{T \in \mathcal{T}_h} h_T^2 \left\| \mathbf{div} \left( \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^{\dagger})^{\mathbf{d}} \right) \right\|_{[L^2(T)]^2}^2 \right\}^{1/2} \|\mathbf{v}\|_{[H^1(\Omega)]^2}, \quad (5.25) \end{aligned}$$

$$\begin{aligned} & \left| \sum_{T \in \mathcal{T}_h} \int_T \mathbf{div} \left( \gamma_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\dagger}) \right) \cdot (\mathbf{v} - \mathbf{v}_h) \right| \\ & \leq C \left\{ \sum_{T \in \mathcal{T}_h} h_T^2 \left\| \mathbf{div} \left( \gamma_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\dagger}) \right) \right\|_{[L^2(T)]^2}^2 \right\}^{1/2} \|\mathbf{v}\|_{[H^1(\Omega)]^2}, \quad (5.26) \end{aligned}$$

$$\begin{aligned} & \left| \sum_{e \in \mathcal{E}_{h,\Omega}} \left\langle J \left[ \left( \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^{\dagger})^{\mathbf{d}} \right) \boldsymbol{\nu}_T \right], \mathbf{v} - \mathbf{v}_h \right\rangle_{[L^2(e)]^2} \right| \\ & \leq C \left\{ \sum_{e \in \mathcal{E}_{h,\Omega}} h_e \left\| J \left[ \left( \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^{\dagger})^{\mathbf{d}} \right) \boldsymbol{\nu}_T \right] \right\|_{[L^2(e)]^2}^2 \right\}^{1/2} \|\mathbf{v}\|_{[H^1(\Omega)]^2}, \quad (5.27) \end{aligned}$$

and

$$\begin{aligned} & \left| \sum_{e \in \mathcal{E}_{h,\Omega}} \left\langle J \left[ \left( \gamma_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\dagger}) \right) \boldsymbol{\nu}_T \right], \mathbf{v} - \mathbf{v}_h \right\rangle_{[L^2(e)]^2} \right| \\ & \leq C \left\{ \sum_{e \in \mathcal{E}_{h,\Omega}} h_e \left\| J \left[ \left( \gamma_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^{\dagger}) \right) \boldsymbol{\nu}_T \right] \right\|_{[L^2(e)]^2}^2 \right\}^{1/2} \|\mathbf{v}\|_{[H^1(\Omega)]^2}. \quad (5.28) \end{aligned}$$

Therefore, placing (5.19) - (5.24) (resp. (5.25) - (5.28)) back into (5.13) - (5.15) (resp. (5.16) and (5.17)), employing the estimates (5.10), (5.11), (5.12) and (5.18), and using the identities

$$\sum_{e \in \mathcal{E}_{h,\Omega}} \int_e = \frac{1}{2} \sum_{T \in \mathcal{T}_h} \sum_{e \in \mathcal{E}(T) \cap \mathcal{E}_{h,\Omega}} \int_e$$

and

$$\sum_{e \in \mathcal{E}_h} \int_e = \sum_{e \in \mathcal{E}_{h,\Omega}} \int_e + \sum_{T \in \mathcal{T}_h} \sum_{e \in \mathcal{E}(T) \cap \mathcal{E}_{h,\Gamma}} \int_e,$$

we conclude from (5.9) that

$$\sup_{\substack{(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}) \in \mathbf{H}\{\mathbf{0}\} \\ \mathbf{div}(\boldsymbol{\tau}) = \mathbf{0}}} \frac{A((\boldsymbol{\sigma} - \boldsymbol{\sigma}_h, p - p_h, \mathbf{u} - \mathbf{u}_h, \boldsymbol{\gamma} - \boldsymbol{\gamma}_h), (\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}))}{\|(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta})\|_{\mathbf{H}}} \leq C\boldsymbol{\theta}. \quad (5.29)$$

This inequality and Lemma 5.2 complete the proof of reliability of  $\boldsymbol{\theta}$ .

We remark that when the finite element subspace  $\mathbf{H}_h$  is given by (4.10), that is, when  $\boldsymbol{\sigma}_h|_T \in [\mathbb{RT}_0(T)]^2$ ,  $p_h|_T \in \mathbb{P}_0(T)$ ,  $\mathbf{u}_h|_T \in [\mathbb{P}_1(T)]^2$ , and  $\boldsymbol{\gamma}_h|_T \in [\mathbb{P}_0(T)]^{2 \times 2}$ , then the expression (5.1) for  $\theta_T^2$  simplifies to

$$\begin{aligned}
\theta_T^2 &:= \|\mathbf{f} + \mathbf{div}(\boldsymbol{\sigma}_h)\|_{[L^2(T)]^2}^2 + \|\boldsymbol{\sigma}_h - \boldsymbol{\sigma}_h^\dagger\|_{[L^2(T)]^{2 \times 2}}^2 + \|\boldsymbol{\gamma}_h - \frac{1}{2}(\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^\dagger)\|_{[L^2(T)]^{2 \times 2}}^2 \\
&+ h_T^2 \|\mathbf{curl}\left(\frac{1}{2\mu}\boldsymbol{\sigma}_h^d\right)\|_{[L^2(T)]^2}^2 + h_T^2 \|\mathbf{curl}\left(\frac{1}{2}\text{tr}(\boldsymbol{\sigma}_h)\right)\|_{[L^2(T)]^2}^2 \\
&+ \sum_{e \in \mathcal{E}(T)} h_e \|J\left[\left(\frac{1}{2\mu}\boldsymbol{\sigma}_h^d - \nabla \mathbf{u}_h + \boldsymbol{\gamma}_h\right)\mathbf{t}_T\right]\|_{[L^2(e)]^2}^2 \\
&+ \sum_{e \in \mathcal{E}(T)} h_e \|J\left[\left(p_h + \frac{1}{2}\text{tr}(\boldsymbol{\sigma}_h)\right)\mathbf{t}_T\right]\|_{[L^2(e)]^2}^2 + \sum_{e \in \mathcal{E}(T)} h_e \|J\left[\left(\boldsymbol{\varepsilon}(\mathbf{u}_h)^d - \frac{1}{2\mu}\boldsymbol{\sigma}_h^d\right)\mathbf{t}_T\right]\|_{[L^2(e)]^2}^2 \\
&+ h_T^2 \|\mathbf{div}\left(\frac{1}{2\mu}\frac{1}{2}(\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^\dagger)^d\right)\|_{[L^2(T)]^2}^2 \\
&+ \sum_{e \in \mathcal{E}(T) \cap \mathcal{E}_{h,\Omega}} h_e \|J\left[\left(\boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu}\frac{1}{2}(\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^\dagger)^d\right)\boldsymbol{\nu}_T\right]\|_{[L^2(e)]^2}^2 \\
&+ \sum_{e \in \mathcal{E}(T) \cap \mathcal{E}_{h,\Omega}} h_e \|J\left[\left(\boldsymbol{\gamma}_h - \frac{1}{2}(\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^\dagger)\right)\boldsymbol{\nu}_T\right]\|_{[L^2(e)]^2}^2.
\end{aligned} \tag{5.30}$$

**5.2. Efficiency of the a posteriori error estimator.** In this section we proceed as in [5] and apply results ultimately based on inverse inequalities (see [11]) and the localization technique introduced in [25], which is based on triangle-bubble and edge-bubble functions, to prove the efficiency of our a posteriori estimator  $\boldsymbol{\theta}$  (lower bound of the estimate (5.2)).

Our goal is to estimate the thirteen terms defining the error indicator  $\theta_T^2$  (cf. (5.1)). Using  $\mathbf{f} = -\mathbf{div}(\boldsymbol{\sigma})$ , the symmetry of  $\boldsymbol{\sigma}$ , and  $\boldsymbol{\gamma} = \frac{1}{2}(\nabla \mathbf{u} - (\nabla \mathbf{u})^\dagger)$ , we first observe that there hold

$$\|\mathbf{f} + \mathbf{div}(\boldsymbol{\sigma}_h)\|_{[L^2(T)]^2}^2 = \|\mathbf{div}(\boldsymbol{\sigma} - \boldsymbol{\sigma}_h)\|_{[L^2(T)]^2}^2, \tag{5.31}$$

$$\|\boldsymbol{\sigma}_h - \boldsymbol{\sigma}_h^\dagger\|_{[L^2(T)]^{2 \times 2}}^2 \leq 4 \|\boldsymbol{\sigma} - \boldsymbol{\sigma}_h\|_{[L^2(T)]^{2 \times 2}}^2, \tag{5.32}$$

and

$$\|\boldsymbol{\gamma}_h - \frac{1}{2}(\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^\dagger)\|_{[L^2(T)]^{2 \times 2}}^2 \leq 2 \left\{ \|\boldsymbol{\gamma} - \boldsymbol{\gamma}_h\|_{[L^2(T)]^{2 \times 2}}^2 + |\mathbf{u} - \mathbf{u}_h|_{[H^1(T)]^2}^2 \right\}. \tag{5.33}$$

The upper bounds of the remaining ten terms, which depend on the mesh parameters  $h_T$  and  $h_e$ , will be derived next. To this end we will make use of Lemmata 5.4 - 5.7 below. Lemma 5.4 is required for the terms involving the curl and  $\mathbf{curl}$  operators, Lemma 5.5 handles the terms involving tangential jumps across the edges of  $\mathcal{T}_h$ , Lemma 5.6 is required for the terms containing the  $\mathbf{div}$  operator, and Lemma 5.7 is used to take care of the terms encompassing normal jumps across the edges of  $\mathcal{T}_h$ . For their proofs we refer to [5] and references therein. In what follows, we let

$$w_e := \cup\{T' \in \mathcal{T}_h : e \in \mathcal{E}(T')\}.$$

**LEMMA 5.4.** *Let  $\boldsymbol{\rho}_h \in [L^2(\Omega)]^{2 \times 2}$  be a piecewise polynomial of degree  $k \geq 0$  on each  $T \in \mathcal{T}_h$ . In addition, let  $\boldsymbol{\rho} \in [L^2(\Omega)]^{2 \times 2}$  be such that  $\mathbf{curl}(\boldsymbol{\rho}) = \mathbf{0}$  on each  $T \in \mathcal{T}_h$ . Then, there exists  $c > 0$ , independent of  $h$ , such that for any  $T \in \mathcal{T}_h$*

$$\|\mathbf{curl}(\boldsymbol{\rho}_h)\|_{[L^2(T)]^2} \leq ch_T^{-1} \|\boldsymbol{\rho} - \boldsymbol{\rho}_h\|_{[L^2(T)]^{2 \times 2}}. \tag{5.34}$$

LEMMA 5.5. *Let  $\boldsymbol{\rho}_h \in [L^2(\Omega)]^{2 \times 2}$  be a piecewise polynomial of degree  $k \geq 0$  on each  $T \in \mathcal{T}_h$ . Then, there exists  $c > 0$ , independent of  $h$ , such that for any  $e \in \mathcal{E}_h$*

$$\|J[\boldsymbol{\rho}_h \mathbf{t}_T]\|_{[L^2(e)]^2} \leq ch_e^{-1/2} \|\boldsymbol{\rho}_h\|_{[L^2(\omega_e)]^{2 \times 2}}. \quad (5.35)$$

LEMMA 5.6. *Let  $\boldsymbol{\rho}_h \in [L^2(\Omega)]^{2 \times 2}$  be a piecewise polynomial of degree  $k \geq 0$  on each  $T \in \mathcal{T}_h$ . Then, there exists  $c > 0$ , independent of  $h$ , such that for any  $T \in \mathcal{T}_h$*

$$\|\mathbf{div}(\boldsymbol{\rho}_h)\|_{[L^2(T)]^2} \leq ch_T^{-1} \|\boldsymbol{\rho}_h\|_{[L^2(T)]^{2 \times 2}}. \quad (5.36)$$

LEMMA 5.7. *Let  $\boldsymbol{\rho}_h \in [L^2(\Omega)]^{2 \times 2}$  be a piecewise polynomial of degree  $k \geq 0$  on each  $T \in \mathcal{T}_h$ . Then, there exists  $c > 0$ , independent of  $h$ , such that for any  $e \in \mathcal{E}_h$*

$$\|J[\boldsymbol{\rho}_h \boldsymbol{\nu}_T]\|_{[L^2(e)]^2} \leq ch_e^{-1/2} \|\boldsymbol{\rho}_h\|_{[L^2(\omega_e)]^{2 \times 2}}. \quad (5.37)$$

We now complete the proof of efficiency of  $\boldsymbol{\theta}$  by conveniently applying Lemmata 5.4 - 5.7 to the corresponding terms defining  $\boldsymbol{\theta}_T^2$ .

LEMMA 5.8. *There exist  $C_1, C_2, C_3 > 0$ , independent of  $h$ , such that for any  $T \in \mathcal{T}_h$*

$$h_T^2 \|\mathbf{curl}\left(\frac{1}{2\mu}\boldsymbol{\sigma}_h^{\mathbf{d}} + \boldsymbol{\gamma}_h\right)\|_{[L^2(T)]^2}^2 \leq C_1 \left\{ \|\boldsymbol{\sigma} - \boldsymbol{\sigma}_h\|_{[L^2(T)]^{2 \times 2}}^2 + \|\boldsymbol{\gamma} - \boldsymbol{\gamma}_h\|_{[L^2(T)]^{2 \times 2}}^2 \right\}, \quad (5.38)$$

$$h_T^2 \|\mathbf{curl}\left(p_h + \frac{1}{2}\mathbf{tr}(\boldsymbol{\sigma}_h)\right)\|_{[L^2(T)]^2}^2 \leq C_2 \left\{ \|p - p_h\|_{L^2(T)}^2 + \|\boldsymbol{\sigma} - \boldsymbol{\sigma}_h\|_{[L^2(T)]^{2 \times 2}}^2 \right\}, \quad (5.39)$$

and

$$h_T^2 \|\mathbf{curl}\left(\boldsymbol{\varepsilon}(\mathbf{u}_h)^{\mathbf{d}} - \frac{1}{2\mu}\boldsymbol{\sigma}_h^{\mathbf{d}}\right)\|_{[L^2(T)]^2}^2 \leq C_3 \left\{ \|\mathbf{u} - \mathbf{u}_h\|_{[H^1(T)]^2}^2 + \|\boldsymbol{\sigma} - \boldsymbol{\sigma}_h\|_{[L^2(T)]^{2 \times 2}}^2 \right\}. \quad (5.40)$$

*Proof.* Applying Lemma 5.4 with  $\boldsymbol{\rho}_h := \frac{1}{2\mu}\boldsymbol{\sigma}_h^{\mathbf{d}} + \boldsymbol{\gamma}_h$  and  $\boldsymbol{\rho} := \nabla \mathbf{u} = \frac{1}{2\mu}\boldsymbol{\sigma}^{\mathbf{d}} + \boldsymbol{\gamma}$ , and then using the triangle inequality and the continuity of  $\boldsymbol{\tau} \rightarrow \boldsymbol{\tau}^{\mathbf{d}}$  we obtain (5.38). Similarly, (5.39) and (5.40) follow from Lemma 5.4 with  $\boldsymbol{\rho}_h := p_h \mathbf{I} + \frac{1}{2}\mathbf{tr}(\boldsymbol{\sigma}_h)\mathbf{I}$  and  $\boldsymbol{\rho} := p \mathbf{I} + \frac{1}{2}\mathbf{tr}(\boldsymbol{\sigma})\mathbf{I} = \mathbf{0}$  (cf. (2.2)), and  $\boldsymbol{\rho}_h := \boldsymbol{\varepsilon}(\mathbf{u}_h)^{\mathbf{d}} - \frac{1}{2\mu}\boldsymbol{\sigma}_h^{\mathbf{d}}$  and  $\boldsymbol{\rho} := \boldsymbol{\varepsilon}(\mathbf{u})^{\mathbf{d}} - \frac{1}{2\mu}\boldsymbol{\sigma}^{\mathbf{d}} = \mathbf{0}$  (cf. (2.4)), respectively.  $\square$

LEMMA 5.9. *There exist  $C_4, C_5, C_6 > 0$ , independent of  $h$ , such that for any  $e \in \mathcal{E}_h$*

$$\begin{aligned} h_e \|J\left[\left(\frac{1}{2\mu}\boldsymbol{\sigma}_h^{\mathbf{d}} - \nabla \mathbf{u}_h + \boldsymbol{\gamma}_h\right)\mathbf{t}_T\right]\|_{[L^2(e)]^2}^2 \\ \leq C_4 \left\{ \|\boldsymbol{\sigma} - \boldsymbol{\sigma}_h\|_{[L^2(\omega_e)]^{2 \times 2}}^2 + \|\mathbf{u} - \mathbf{u}_h\|_{[H^1(\omega_e)]^2}^2 + \|\boldsymbol{\gamma} - \boldsymbol{\gamma}_h\|_{[L^2(\omega_e)]^{2 \times 2}}^2 \right\}, \end{aligned} \quad (5.41)$$

$$h_e \|J\left[\left(p_h + \frac{1}{2}\mathbf{tr}(\boldsymbol{\sigma}_h)\right)\mathbf{t}_T\right]\|_{[L^2(e)]^2}^2 \leq C_5 \left\{ \|p - p_h\|_{L^2(\omega_e)}^2 + \|\boldsymbol{\sigma} - \boldsymbol{\sigma}_h\|_{[L^2(\omega_e)]^{2 \times 2}}^2 \right\}, \quad (5.42)$$

and

$$h_e \|J\left[\left(\boldsymbol{\varepsilon}(\mathbf{u}_h)^{\mathbf{d}} - \frac{1}{2\mu}\boldsymbol{\sigma}_h^{\mathbf{d}}\right)\mathbf{t}_T\right]\|_{[L^2(e)]^2}^2 \leq C_6 \left\{ \|\mathbf{u} - \mathbf{u}_h\|_{[H^1(\omega_e)]^2}^2 + \|\boldsymbol{\sigma} - \boldsymbol{\sigma}_h\|_{[L^2(\omega_e)]^{2 \times 2}}^2 \right\}. \quad (5.43)$$

*Proof.* The estimate (5.41) follows from Lemma 5.5 with  $\boldsymbol{\rho}_h := \frac{1}{2\mu}\boldsymbol{\sigma}_h^d - \nabla \mathbf{u}_h + \boldsymbol{\gamma}_h$ , introducing  $\mathbf{0} = \frac{1}{2\mu}\boldsymbol{\sigma}^d - \nabla \mathbf{u} + \boldsymbol{\gamma}$  (cf. (2.2) - (2.4)) in the resulting estimate and applying the triangle inequality and the continuity of  $\boldsymbol{\tau} \longrightarrow \boldsymbol{\tau}^d$ . Analogously, estimate (5.42) (resp. (5.43)) is obtained from Lemma 5.5 defining  $\boldsymbol{\rho}_h$  as  $(p_h + \frac{1}{2}\text{tr}(\boldsymbol{\sigma}_h))\mathbf{I}$  (resp.  $\boldsymbol{\varepsilon}(\mathbf{u}_h)^d - \frac{1}{2\mu}\boldsymbol{\sigma}_h^d$ ) and then introducing  $\mathbf{0} = (p + \frac{1}{2}\text{tr}(\boldsymbol{\sigma}))\mathbf{I}$  (resp.  $\mathbf{0} = \boldsymbol{\varepsilon}(\mathbf{u}_h)^d - \frac{1}{2\mu}\boldsymbol{\sigma}_h^d$ ) (cf. (2.2) (resp. (2.4))).  $\square$

LEMMA 5.10. *There exist  $C_7, C_8 \geq 0$ , independent of  $h$ , such that for any  $T \in \mathcal{T}_h$*

$$h_T^2 \|\text{div} \left( \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^t)^d \right)\|_{[L^2(T)]^2}^2 \leq C_7 \left\{ |\mathbf{u} - \mathbf{u}_h|_{[H^1(T)]^2}^2 + \|\boldsymbol{\sigma} - \boldsymbol{\sigma}_h\|_{[L^2(T)]^{2 \times 2}}^2 \right\} \quad (5.44)$$

and

$$h_T^2 \|\text{div} \left( \boldsymbol{\gamma}_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^t) \right)\|_{[L^2(T)]^2}^2 \leq C_8 \left\{ \|\boldsymbol{\gamma} - \boldsymbol{\gamma}_h\|_{[L^2(T)]^{2 \times 2}}^2 + |\mathbf{u} - \mathbf{u}_h|_{[H^1(T)]^2}^2 \right\}. \quad (5.45)$$

*Proof.* The estimate (5.44) follows from Lemma 5.6 defining  $\boldsymbol{\rho}_h := \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^t)^d$ , introducing  $\mathbf{0} = \boldsymbol{\varepsilon}(\mathbf{u}) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma} + \boldsymbol{\sigma}^t)^d$  (cf. (2.4)), and then using the triangle inequality and the continuity of the operators  $\boldsymbol{\varepsilon}$  and  $\boldsymbol{\tau} \longrightarrow \boldsymbol{\tau}^d$ . Similarly, applying Lemma 5.6 with  $\boldsymbol{\rho}_h := \boldsymbol{\gamma}_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^t)$  and introducing  $\mathbf{0} = \boldsymbol{\gamma} - \frac{1}{2} (\nabla \mathbf{u} - (\nabla \mathbf{u})^t)$  yields (5.45).  $\square$

LEMMA 5.11. *There exist  $C_9, C_{10} > 0$ , independent of  $h$ , such that for any  $e \in \mathcal{E}_h$*

$$h_e \|J \left[ \left( \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^t)^d \right) \boldsymbol{\nu}_T \right]\|_{[L^2(e)]^2}^2 \leq C_9 \left\{ |\mathbf{u} - \mathbf{u}_h|_{[H^1(T)]^2}^2 + \|\boldsymbol{\sigma} - \boldsymbol{\sigma}_h\|_{[L^2(T)]^{2 \times 2}}^2 \right\} \quad (5.46)$$

and

$$h_e \|J \left[ \left( \boldsymbol{\gamma}_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^t) \right) \boldsymbol{\nu}_T \right]\|_{[L^2(e)]^2}^2 \leq C_{10} \left\{ \|\boldsymbol{\gamma} - \boldsymbol{\gamma}_h\|_{[L^2(T)]^{2 \times 2}}^2 + |\mathbf{u} - \mathbf{u}_h|_{[H^1(T)]^2}^2 \right\}. \quad (5.47)$$

*Proof.* The estimate (5.46) follows from Lemma 5.7 with  $\boldsymbol{\rho}_h := \boldsymbol{\varepsilon}(\mathbf{u}_h) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma}_h + \boldsymbol{\sigma}_h^t)^d$ , introducing  $\mathbf{0} = \boldsymbol{\varepsilon}(\mathbf{u}) - \frac{1}{2\mu} \frac{1}{2} (\boldsymbol{\sigma} + \boldsymbol{\sigma}^t)^d$  (cf. (2.4)) and then employing again the triangle inequality and the continuity of the operators  $\boldsymbol{\varepsilon}$  and  $\boldsymbol{\tau} \longrightarrow \boldsymbol{\tau}^d$ . Analogously, the estimate (5.47) follows from Lemma 5.7 defining  $\boldsymbol{\rho}_h := \boldsymbol{\gamma}_h - \frac{1}{2} (\nabla \mathbf{u}_h - (\nabla \mathbf{u}_h)^t)$  and then introducing  $\mathbf{0} = \boldsymbol{\gamma} - \frac{1}{2} (\nabla \mathbf{u} - (\nabla \mathbf{u})^t)$ .  $\square$

Thus, the efficiency of  $\boldsymbol{\theta}$  follows straightforwardly from the estimates (5.31) - (5.47) after summing over all  $T \in \mathcal{T}_h$  and using that the number of triangles on each domain  $\omega_e$  is bounded by two.

## 6. NUMERICAL RESULTS

In this section we present several numerical results illustrating the performance of the augmented finite element scheme (4.1) and the a posteriori error estimator  $\boldsymbol{\theta}$  analyzed in Section 5, using the specific finite element subspace  $\tilde{\mathbf{H}}_h$  (cf. (4.10)). We recall that in this case the local indicator  $\theta_T^2$  reduces to (5.30). Now, in order to implement the zero integral mean condition for functions of the space  $\tilde{H}_{0,h}^\sigma$  (cf. (4.5)), we introduce, as described in [18],

a Lagrange multiplier  $\varphi_h \in \mathbb{R}$ . That is, instead of (4.1) with  $\mathbf{H}_h = \tilde{\mathbf{H}}_h$ , we consider the equivalent problem: Find  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h, \varphi_h) \in \tilde{H}_h^\boldsymbol{\sigma} \times \tilde{H}_h^p \times \tilde{H}_{0,h}^\mathbf{u} \times \tilde{H}_h^\boldsymbol{\gamma} \times \mathbb{R}$  such that

$$\begin{aligned} A((\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h), (\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h)) + \varphi_h \int_{\Omega} \operatorname{tr}(\boldsymbol{\tau}_h) &= F((\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h)), \\ \psi_h \int_{\Omega} \operatorname{tr}(\boldsymbol{\sigma}_h) &= 0, \end{aligned} \quad (6.1)$$

for all  $(\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h, \psi_h) \in \tilde{H}_h^\boldsymbol{\sigma} \times \tilde{H}_h^p \times \tilde{H}_{0,h}^\mathbf{u} \times \tilde{H}_h^\boldsymbol{\gamma} \times \mathbb{R}$ . We state the equivalence between (4.1) and (6.1) through the application of the following Theorem, adapted from Theorem 4.3 in [18].

**THEOREM 6.1.**

- (1) Let  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h) \in \tilde{\mathbf{H}}_h$  be the solution of (4.1). Then  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h, 0)$  is a solution of (6.1).
- (2) Let  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h, \varphi_h) \in \tilde{H}_h^\boldsymbol{\sigma} \times \tilde{H}_h^p \times \tilde{H}_{0,h}^\mathbf{u} \times \tilde{H}_h^\boldsymbol{\gamma} \times \mathbb{R}$  be a solution of (6.1). Then  $\varphi_h = 0$  and  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h)$  is the solution of (4.1).

*Proof.* We first observe, according to the definition of  $A$  (cf. (3.5)), that for each  $(\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}) \in H(\mathbf{div}; \Omega) \times L^2(\Omega) \times [H_0^1(\Omega)]^2 \times [L^2(\Omega)]_{\text{asym}}^{2 \times 2}$  there holds

$$A((\boldsymbol{\tau}, q, \mathbf{v}, \boldsymbol{\eta}), (\mathbf{I}, -1, \mathbf{0}, \mathbf{0})) = 0. \quad (6.2)$$

Now, let  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h)$  be the solution of (4.1), and let  $(\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \in \tilde{H}_h^\boldsymbol{\sigma} \times \tilde{H}_h^p \times \tilde{H}_{0,h}^\mathbf{u} \times \tilde{H}_h^\boldsymbol{\gamma}$ . We write  $\boldsymbol{\tau}_h = \boldsymbol{\tau}_{0,h} + d_h \mathbf{I}$ , with  $\boldsymbol{\tau}_{0,h} \in \tilde{H}_{0,h}^\boldsymbol{\sigma}$  and  $d_h \in \mathbb{R}$  and observe that  $(\boldsymbol{\tau}_{0,h}, q_h + d_h, \mathbf{v}_h, \boldsymbol{\eta}_h) \in \tilde{\mathbf{H}}_h$ , whence (3.6), (4.1) and (6.2) yield

$$\begin{aligned} F(\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h) &= F(\boldsymbol{\tau}_{0,h}, q_h + d_h, \mathbf{v}_h, \boldsymbol{\eta}_h) = A((\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h), (\boldsymbol{\tau}_{0,h}, q_h + d_h, \mathbf{v}_h, \boldsymbol{\eta}_h)) \\ &= A((\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h), (\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h)). \end{aligned}$$

This identity and the fact that  $\boldsymbol{\sigma}_h$  clearly satisfies the second equation of (6.1), show that  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h, 0)$  is indeed a solution of (6.1).

Conversely, let  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h, \varphi_h) \in \tilde{H}_h^\boldsymbol{\sigma} \times \tilde{H}_h^p \times \tilde{H}_{0,h}^\mathbf{u} \times \tilde{H}_h^\boldsymbol{\gamma} \times \mathbb{R}$  be a solution of (6.1). Then, taking  $(\boldsymbol{\tau}_h, q_h, \mathbf{v}_h, \boldsymbol{\eta}_h) = (\mathbf{I}, -1, \mathbf{0}, \mathbf{0})$  in the first equation of (6.1) and using (3.6) and (6.2), we find that  $\varphi_h = 0$ , whence  $(\boldsymbol{\sigma}_h, p_h, \mathbf{u}_h, \boldsymbol{\gamma}_h)$  becomes the solution of (4.1).  $\square$

In what follows,  $N$  stands for the total number of degrees of freedom (unknowns) of (6.1), which, at least for uniform refinements, behaves asymptotically as six times the numbers of elements of each triangulation. Also, the individual and total errors are denoted by

$$\begin{aligned} e(\boldsymbol{\sigma}) &:= \|\boldsymbol{\sigma} - \boldsymbol{\sigma}_h\|_{H(\mathbf{div}; \Omega)}, & e(p) &:= \|p - p_h\|_{L^2(\Omega)}, \\ e(\mathbf{u}) &:= \|\mathbf{u} - \mathbf{u}_h\|_{[H^1(\Omega)]^2}, & e(\boldsymbol{\gamma}) &:= \|\boldsymbol{\gamma} - \boldsymbol{\gamma}_h\|_{[L^2(\Omega)]^{2 \times 2}}, \end{aligned}$$

and

$$e := \{[e(\boldsymbol{\sigma})]^2 + [e(p)]^2 + [e(\mathbf{u})]^2 + [e(\boldsymbol{\gamma})]^2\}^{1/2},$$

respectively, whereas the effectivity index with respect to  $\boldsymbol{\theta}$  is defined by  $e/\boldsymbol{\theta}$ .

Since the augmented method (for the compressible case) was shown in [18] to be robust with respect to the parameters  $\kappa_1$ ,  $\kappa_2$ , and  $\kappa_3$ , we simply consider for all the examples  $(\kappa_1, \kappa_2, \kappa_3) = (\mu, \frac{1}{2\mu}, \frac{\mu}{2})$ , which satisfy the assumptions of Theorem 4.4. In addition, since the choice of  $\kappa_0$  in (4.1) depends on the unknown constant  $c_1$  from Lemma 2.4, we simply take here  $\kappa_0 = \mu$ . As we will see below, this choice works out well in all the examples

We now specify the data of the three examples to be presented here. We take  $\Omega$  as either the square  $]0, 1[^2$  or the triangle  $\hat{T} := \{(x_1, x_2) : x_1, x_2 > 0 \text{ and } x_1 + x_2 < 1\}$ , and choose the

datum  $\mathbf{f}$  so that the exact solution  $\mathbf{u}(x_1, x_2) := (u_1(x_1, x_2), u_2(x_1, x_2))^t$  and  $p(x_1, x_2)$  are given in the table below. Actually, according to (2.1) we have  $\boldsymbol{\sigma} = 2\mu\boldsymbol{\varepsilon}(\mathbf{u}) - p\mathbf{I}$ , and hence simple computations show that  $\mathbf{f} := -\mathbf{div}(\boldsymbol{\sigma}) = -\mu\Delta\mathbf{u} - \mu\nabla(\mathbf{div}\mathbf{u}) + \nabla p = -\mu\Delta\mathbf{u} + \nabla p$ . We also recall that the rotation  $\gamma$  is defined by  $\frac{1}{2}(\nabla\mathbf{u} - (\nabla\mathbf{u})^t)$ . In all the examples we take  $\mu = 1.0$ .

We emphasize that from (2.1) an admissible solution  $\mathbf{u}$  must satisfy both  $\mathbf{u} = \mathbf{0}$  on  $\Gamma$  and  $\mathbf{div}(\mathbf{u}) = \mathbf{0}$  in  $\Omega$ , and from (2.2) and the fact that  $\boldsymbol{\sigma} \in H_0$  (cf. (2.6)) an admissible solution  $p$  must satisfy  $\int_{\Omega} p = 0$ .

EXAMPLE	$\Omega$	$\mathbf{u}(x_1, x_2)$	$p(x_1, x_2)$
1	$]0, 1[^2$	$\mathbf{curl}(x_1^2 x_2^2 (x_1 - 1)^2 (x_2 - 1)^2)$	$x_1^2 + x_2^2 - \frac{2}{3}$
2	$\hat{T}$	$10^2 \mathbf{curl}(x_1^2 x_2^2 (1 - x_1 - x_2)^2 (x_1^2 + x_2^2)^{-3/4})$	$x_1^2 + x_2^2 - \frac{1}{3}$
3	$]0, 1[^2$	$\mathbf{curl}\left(\frac{9x_1^2 x_2^2 (1 - x_1)^2 (1 - x_2)^2}{(300x_1 - 100)^2 + (300x_2 - 100)^2 + 90}\right)$	$\left(\frac{x_1}{100}\right)^2 + \left(\frac{x_2}{100}\right)^2 - \frac{2}{3} \times 10^{-4}$

We observe that the solution of Example 2 is singular at the boundary point  $(0, 0)$ . Thus, according to Theorem 4.5 we expect a rate of convergence lower than 1 for the uniform refinement. On the other hand, the solution of Example 3 shows a large stress region in the vicinity of the interior point  $(1/3, 1/3)$ .

The numerical results shown below were obtained in a *Pentium Xeon computer with dual processors* using a Fortran Code and the TRIANGLE mesh generator. The linear system arising from (6.1) is solved with the sequential LU package. Individual errors are computed on each triangle using a Gaussian quadrature rule.

We first utilize the Example 1 to illustrate the good behaviour of the a posteriori error estimator  $\boldsymbol{\theta}$  in a sequence of quasi-uniform meshes. In Table 1 we present the individual and total errors, the a posteriori estimators, and the effectivity indexes for this example with this sequence of quasi-uniform meshes. The index always remains in a neighborhood of 0.600 in this example, which confirms the reliability and efficiency of  $\boldsymbol{\theta}$ .

Next we consider Examples 2 and 3 to illustrate the performance of the following adaptive algorithm based on  $\boldsymbol{\theta}$  for the computation of solutions of (6.1):

1. Start with a coarse mesh  $\mathcal{T}_h$ .
2. Solve the Galerkin scheme (6.1) for the current mesh  $\mathcal{T}_h$ .
3. Compute  $\theta_T$  for each triangle  $T \in \mathcal{T}_h$ .
4. Consider stopping criterion and decide to finish or go to next step.
5. Instruct the mesh generator to ensure that in the next mesh the region enclosed by each element  $T' \in \mathcal{T}_h$  of the current mesh whose local indicator  $\theta_{T'}$  satisfies  $\theta_{T'} \geq \frac{1}{2} \max\{\theta_T : T \in \mathcal{T}_h\}$  encompasses no triangle with area larger than  $\frac{|T'|}{4}$ .
6. Generate the next mesh, store it as  $\mathcal{T}_h$  and go to step 2.

At this point we introduce the experimental rate of convergence, which, given two consecutive triangulations with degrees of freedom  $N$  and  $N'$  and corresponding errors  $e$  and  $e'$ , is defined by

$$r(e) := -2 \frac{\log(e/e')}{\log(N/N')}.$$

In Tables 2 through 5 we provide the individual and total errors, the experimental rates of convergence, the a posteriori error estimators and the effectivity indexes for the uniform and adaptive refinements as applied to Examples 2 and 3. In this case the quasi-uniform sequences of meshes are generated by instructing the mesh generator to provide only triangles with area

TABLE 1. Mesh sizes, individual and total errors, a posteriori error estimators, and effectivity indexes for a sequence of quasi-uniform meshes (Example 1).

$N$	$h$	$e(\boldsymbol{\sigma})$	$e(p)$	$e(\mathbf{u})$	$e(\boldsymbol{\gamma})$	$e$	$\boldsymbol{\theta}$	$e/\boldsymbol{\theta}$
99	0.500	0.681E-00	0.151E-00	0.130E-00	0.587E-01	0.712E-00	0.923E-00	0.772
165	0.500	0.557E-00	0.126E-00	0.844E-01	0.453E-01	0.579E-00	0.794E-00	0.729
207	0.500	0.528E-00	0.115E-00	0.818E-01	0.428E-01	0.548E-00	0.738E-00	0.743
363	0.288	0.374E-00	0.791E-01	0.609E-01	0.367E-01	0.389E-00	0.589E-00	0.660
435	0.271	0.345E-00	0.756E-01	0.584E-01	0.315E-01	0.359E-00	0.523E-00	0.687
627	0.257	0.282E-00	0.601E-01	0.463E-01	0.271E-01	0.293E-00	0.452E-00	0.648
849	0.250	0.253E-00	0.555E-01	0.420E-01	0.273E-01	0.264E-00	0.412E-00	0.639
1245	0.250	0.204E-00	0.485E-01	0.358E-01	0.220E-01	0.214E-00	0.335E-00	0.638
1707	0.147	0.181E-00	0.388E-01	0.305E-01	0.198E-01	0.188E-00	0.303E-00	0.622
2433	0.125	0.148E-00	0.323E-01	0.254E-01	0.154E-01	0.155E-00	0.243E-00	0.635
3369	0.125	0.128E-00	0.287E-01	0.218E-01	0.135E-01	0.133E-00	0.211E-00	0.632
4833	0.125	0.103E-00	0.229E-01	0.185E-01	0.120E-01	0.108E-00	0.180E-00	0.603
6927	0.077	0.880E-01	0.188E-01	0.154E-01	0.961E-02	0.918E-01	0.149E-00	0.615
9681	0.065	0.743E-01	0.159E-01	0.131E-01	0.851E-02	0.776E-01	0.129E-00	0.601
13563	0.062	0.632E-01	0.137E-01	0.112E-01	0.736E-02	0.661E-01	0.111E-00	0.595

TABLE 2. Individual and total errors, experimental rates of convergence, a posteriori error estimators, and effectivity indexes for a sequence of quasi-uniform meshes (Example 2).

$N$	$e(\boldsymbol{\sigma})$	$e(p)$	$e(\mathbf{u})$	$e(\boldsymbol{\gamma})$	$e$	$r(e)$	$\boldsymbol{\theta}$	$e/\boldsymbol{\theta}$
159	0.159E+03	0.626E+01	0.103E+02	0.527E+01	0.160E+03	—	0.166E+03	0.965
633	0.122E+03	0.363E+01	0.677E+01	0.366E+01	0.123E+03	0.383	0.127E+03	0.965
2367	0.941E+02	0.202E+01	0.365E+01	0.198E+01	0.942E+02	0.403	0.961E+02	0.979
9591	0.725E+02	0.107E+01	0.188E+01	0.106E+01	0.725E+02	0.373	0.733E+02	0.989

below a decreasing threshold, subject to a minimum angle constraint. We observe from these tables that the errors of the adaptive procedure decrease much faster than those obtained by the quasi-uniform one, which is confirmed by the experimental rates of convergence provided there. This fact can also be seen in Figures 1 and 2 where we display the total error  $e$  vs. the degrees of freedom  $N$  for both refinements. As shown by the values of  $r(e)$ , particularly in Example 2 (where  $r(e)$  approaches 0.38 for the quasi-uniform refinement), the adaptive method is able to recover, at least approximately, the quasi-optimal rate of convergence  $\mathcal{O}(h)$  for the total error. Furthermore, the effectivity indexes remain again bounded from above and below, which confirms the reliability and efficiency of  $\boldsymbol{\theta}$  for the adaptive algorithm. On the other hand, some intermediate meshes obtained with the adaptive refinement are displayed in Figures 3 and 4. Note that the method is able to recognize the singularities and large stress regions of the solutions. In particular, this fact is observed in Example 2 (see Figure 3) where adapted meshes are highly refined around the singular point  $(0,0)$ . Similarly, the adapted meshes obtained in Example 3 (see Figure 4) concentrate the refinement around the interior point  $(1/3, 1/3)$ , where the largest stress occur.

Summarizing, the numerical results presented in this section exhibit, on one hand, the expected  $\mathcal{O}(h)$  behaviour of this augmented method for smooth problems and, on the other hand, underline the reliability and efficiency of  $\boldsymbol{\theta}$ . In addition, they strongly demonstrate that the associated adaptive algorithm is much more suitable than a uniform discretization procedure when solving problems with non-smooth solutions.



TABLE 3. Individual and total errors, experimental rates of convergence, a posteriori error estimators, and effectivity indexes for the adaptive refinement (Example 2).

$N$	$e(\boldsymbol{\sigma})$	$e(p)$	$e(\mathbf{u})$	$e(\boldsymbol{\gamma})$	$e$	$r(e)$	$\boldsymbol{\theta}$	$e/\boldsymbol{\theta}$
159	0.159E+03	0.626E+01	0.103E+02	0.527E+01	0.160E+03	—	0.166E+03	0.965
249	0.119E+03	0.569E+01	0.892E+01	0.489E+01	0.119E+03	1.301	0.128E+03	0.929
345	0.109E+03	0.542E+01	0.837E+01	0.454E+01	0.110E+03	0.508	0.118E+03	0.926
417	0.993E+02	0.541E+01	0.836E+01	0.455E+01	0.999E+02	1.021	0.109E+03	0.912
531	0.910E+02	0.542E+01	0.825E+01	0.451E+01	0.916E+02	0.716	0.101E+03	0.899
627	0.841E+02	0.485E+01	0.809E+01	0.422E+01	0.847E+02	0.944	0.943E+02	0.898
981	0.711E+02	0.379E+01	0.622E+01	0.314E+01	0.715E+02	0.758	0.781E+02	0.915
1545	0.578E+02	0.313E+01	0.534E+01	0.273E+01	0.582E+02	0.906	0.642E+02	0.906
1899	0.560E+02	0.267E+01	0.441E+01	0.233E+01	0.563E+02	0.324	0.610E+02	0.922
2571	0.499E+02	0.255E+01	0.410E+01	0.210E+01	0.501E+02	0.759	0.544E+02	0.922
3651	0.413E+02	0.224E+01	0.353E+01	0.187E+01	0.416E+02	1.068	0.456E+02	0.912
5187	0.355E+02	0.202E+01	0.325E+01	0.162E+01	0.357E+02	0.867	0.390E+02	0.915
6957	0.310E+02	0.184E+01	0.297E+01	0.149E+01	0.312E+02	0.910	0.344E+02	0.906
9843	0.253E+02	0.133E+01	0.216E+01	0.115E+01	0.254E+02	1.179	0.280E+02	0.909
13707	0.214E+02	0.114E+01	0.194E+01	0.102E+01	0.215E+02	1.014	0.238E+02	0.904

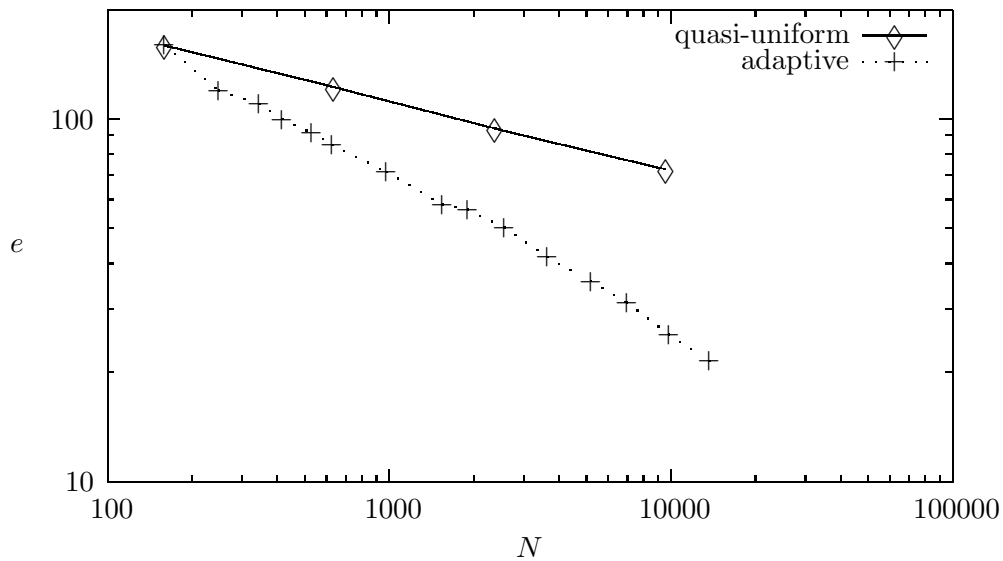


FIGURE 1. Total errors  $e$  vs. degrees of freedom  $N$  for the quasi-uniform and adaptive refinements (Example 2).

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TABLE 4. Individual and total errors, experimental rates of convergence, a posteriori error estimators, and effectivity indexes for a sequence of quasi-uniform meshes (Example 3).

$N$	$e(\boldsymbol{\sigma})$	$e(p)$	$e(\mathbf{u})$	$e(\boldsymbol{\gamma})$	$e$	$r(e)$	$\boldsymbol{\theta}$	$e/\boldsymbol{\theta}$
435	0.228E+03	0.691E+00	0.210E+01	0.106E+01	0.228E+03	—	0.228E+03	1.000
1245	0.195E+03	0.738E+01	0.562E+01	0.248E+01	0.195E+03	0.294	0.196E+03	0.999
3369	0.204E+03	0.613E+01	0.354E+01	0.177E+01	0.204E+03	—	0.204E+03	1.000
9681	0.133E+03	0.126E+01	0.183E+01	0.102E+01	0.133E+03	0.815	0.133E+03	0.998

TABLE 5. Individual and total errors, experimental rates of convergence, a posteriori error estimators, and effectivity indexes for the adaptive refinement (Example 3).

$N$	$e(\boldsymbol{\sigma})$	$e(p)$	$e(\mathbf{u})$	$e(\boldsymbol{\gamma})$	$e$	$r(e)$	$\boldsymbol{\theta}$	$e/\boldsymbol{\theta}$
435	0.228E+03	0.691E+00	0.210E+01	0.106E+01	0.228E+03	—	0.228E+03	1.000
555	0.221E+03	0.681E+01	0.542E+01	0.299E+01	0.221E+03	0.270	0.222E+03	0.995
651	0.178E+03	0.369E+01	0.300E+01	0.167E+01	0.178E+03	2.702	0.178E+03	0.998
819	0.119E+03	0.119E+01	0.155E+01	0.778E+00	0.119E+03	3.473	0.119E+03	0.997
1083	0.752E+02	0.919E+00	0.115E+01	0.591E+00	0.752E+02	3.320	0.756E+02	0.995
1431	0.500E+02	0.733E+00	0.914E+00	0.544E+00	0.500E+02	2.933	0.504E+02	0.992
2139	0.366E+02	0.453E+00	0.584E+00	0.391E+00	0.366E+02	1.552	0.368E+02	0.992
2775	0.303E+02	0.417E+00	0.488E+00	0.302E+00	0.303E+02	1.453	0.305E+02	0.993
3471	0.261E+02	0.369E+00	0.430E+00	0.262E+00	0.261E+02	1.333	0.262E+02	0.993
4707	0.211E+02	0.314E+00	0.367E+00	0.226E+00	0.211E+02	1.383	0.213E+02	0.992
6399	0.177E+02	0.284E+00	0.323E+00	0.197E+00	0.177E+02	1.123	0.179E+02	0.991
8667	0.151E+02	0.252E+00	0.283E+00	0.172E+00	0.151E+02	1.050	0.153E+02	0.991
12147	0.122E+02	0.220E+00	0.241E+00	0.146E+00	0.122E+02	1.271	0.123E+02	0.990

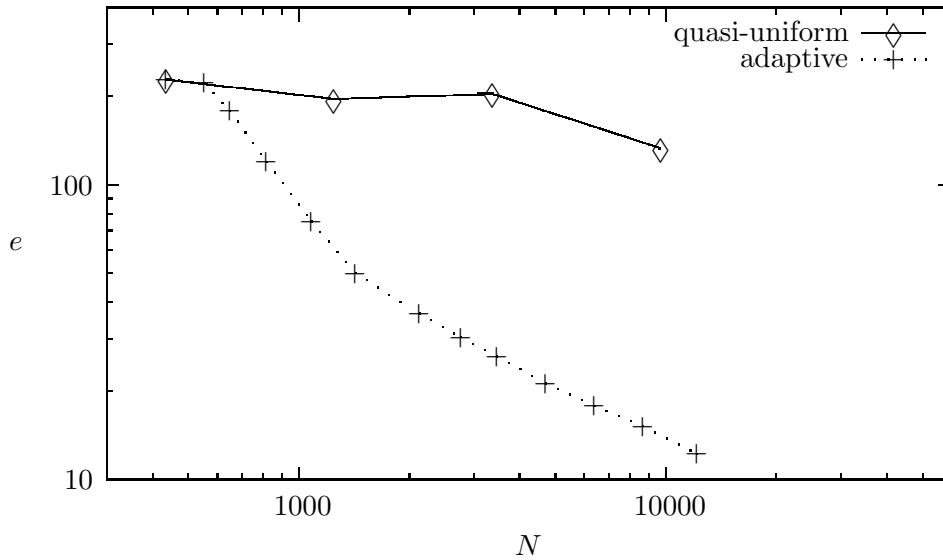


FIGURE 2. Total errors  $e$  vs. degrees of freedom  $N$  for the quasi-uniform and adaptive refinements (Example 3).

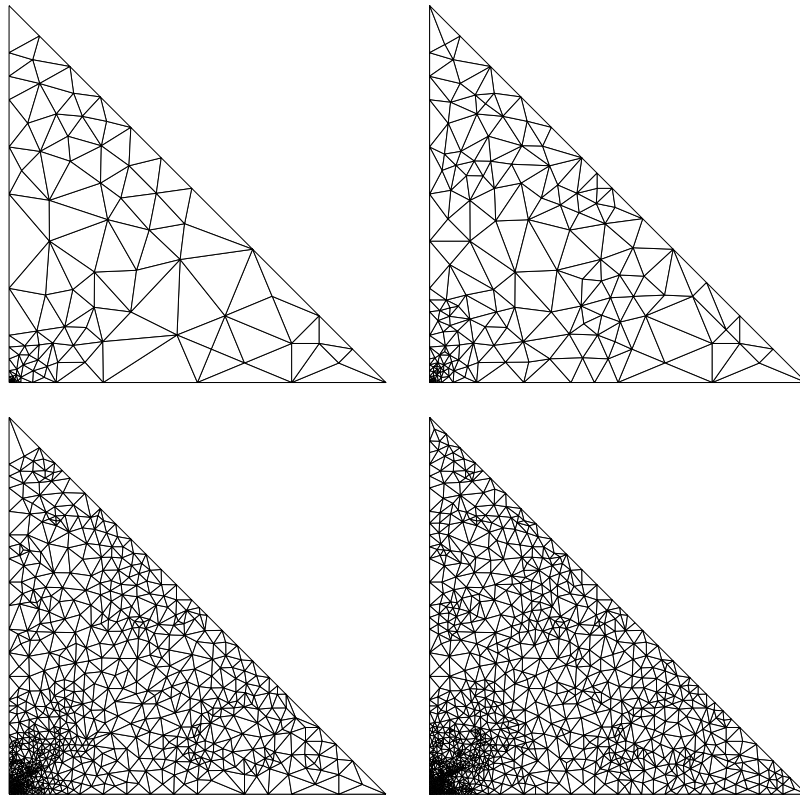


FIGURE 3. Adapted intermediate meshes with 981, 1899, 9843, and 13707 degrees of freedom (Example 2).

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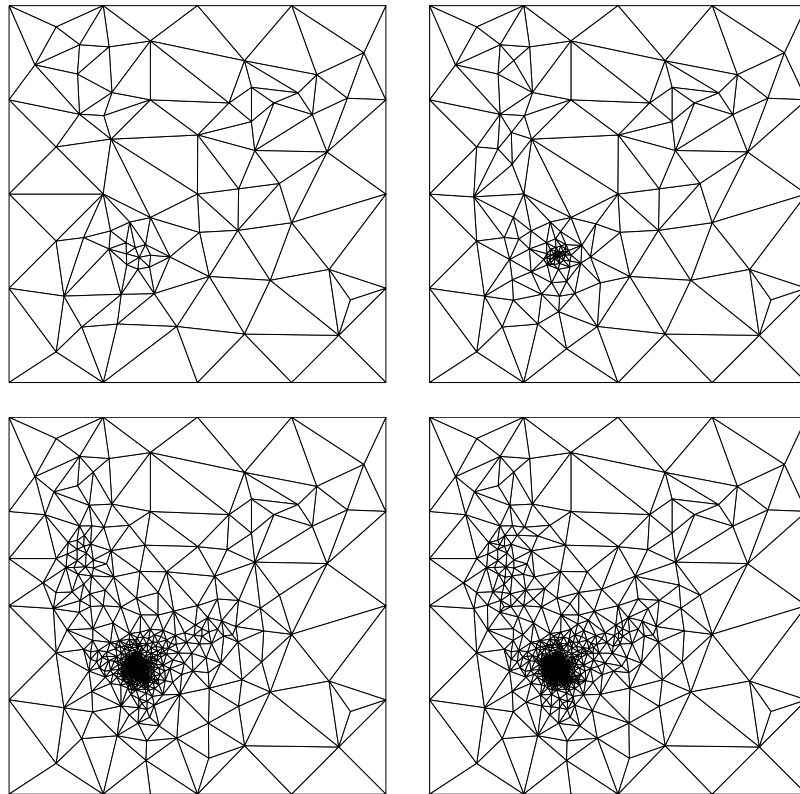


FIGURE 4. Adapted intermediate meshes with 819, 1431, 8667, and 12147 degrees of freedom (Example 3).

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