# EXISTENCE, A PRIORI AND A POSTERIORI ERROR ESTIMATES FOR A NONLINEAR THREE-FIELD PROBLEM ARISING FROM OLDROYD-B VISCOELASTIC FLOWS* 

Marco Picasso ${ }^{1}$ and Jacques Rappaz ${ }^{1}$


#### Abstract

In this paper, a nonlinear problem corresponding to a simplified Oldroyd-B model without convective terms is considered. Assuming the domain to be a convex polygon, existence of a solution is proved for small relaxation times. Continuous piecewise linear finite elements together with a Galerkin Least Square (GLS) method are studied for solving this problem. Existence and a priori error estimates are established using a Newton-chord fixed point theorem, a posteriori error estimates are also derived. An Elastic Viscous Split Stress (EVSS) scheme related to the GLS method is introduced. Numerical results confirm the theoretical predictions.


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## 1. Introduction

Numerical simulation of viscoelastic flows is of great importance for industrial processes involving plastics, paints or food. The modelling of viscoelastic flows generally consists in supplementing the mass and momentum equations with a rheological constitutive equation relating the velocity and the non Newtonian part of the stress.

When solving viscoelastic flows with finite element methods, the following points should be addressed, see for instance [1] for a review. Firstly, the finite element spaces used to approximate the velocity, pressure and extra-stress fields cannot be chosen arbitrarily, an inf-sup condition has to be satisfied [12, 13, 15, 24, 25]. Secondly, due to the presence of convective terms in both momentum and constitutive equations, adequate discretizations procedure must be used such as discontinuous finite elements, GLS stabilization procedures [5], or the characteristics method. Thirdly, the presence of nonlinear terms prevents numerical methods to converge at high Deborah numbers, this being consistent with theoretical [4, 19, 21, 23, 26] and experimental [22] studies.

In this paper, we focus on the first and last of these three points. In [6], a GLS method with continuous, piecewise linear finite elements was proposed for solving a three fields Stokes' problem. The method was stable even when the solvent viscosity was small compared to the polymer viscosity. The link with the EVSS method of [12] was proposed. The aim of this paper is to extend the work of [6] to a nonlinear model problem. Existence, a priori and a posteriori error estimates are derived. Numerical results confirm the theoretical predictions. Even though the model problem studied in this paper is simpler than those considered in [4, 19, 21, 23, 26], we believe that our results are interesting for the following reasons. Firstly, assuming the calculation domain to be a convex

[^0]polygon, we propose an original, but natural, variational setting in order to prove existence of the continuous problem. More precisely, the velocity is in $W^{2, r}$, for some $r>2$, whereas the pressure and the extra-stress are in $W^{1, r}$. Secondly, we consider an original stabilized finite element method and we prove existence, a priori and a posteriori error estimates. Finally, we present some numerical computations that i) confirm the optimality of our theoretical predictions ii) show that the method fails when the computational domain is not convex, thus suggesting that the problem is ill-posed.

The reader should note that the theoretical results presented in this paper can be seen as a first step towards the justification of some numerical methods used to perform mesoscopic calculations $[7,8]$.

## 2. The model problem

Let $\Omega$ be a bounded domain of $\mathbb{R}^{2}$ with Lipschitz boundary $\partial \Omega$. We consider the following problem. Given a force term $\mathbf{f}$, constant solvent and polymer viscosities $\eta_{s} \geqslant 0$ and $\eta_{p}>0$, a constant relaxation time $\lambda$, find the velocity $\mathbf{u}$, pressure $p$ and extra-stress $\boldsymbol{\sigma}$ such that

$$
\begin{align*}
-2 \eta_{s} \operatorname{div} \boldsymbol{\epsilon}(\mathbf{u})+\nabla p-\operatorname{div} \boldsymbol{\sigma} & =\mathbf{f} \\
\operatorname{div} \mathbf{u} & =0  \tag{1}\\
\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}-\frac{\lambda}{2 \eta_{p}}\left((\nabla \mathbf{u}) \boldsymbol{\sigma}+\boldsymbol{\sigma}\left(\nabla \mathbf{u}^{T}\right)\right)-\boldsymbol{\epsilon}(\mathbf{u}) & =0
\end{align*}
$$

in $\Omega$, where the velocity $\mathbf{u}$ vanishes on $\partial \Omega$. Here $\boldsymbol{\epsilon}(\mathbf{u})=\frac{1}{2}\left(\nabla \mathbf{u}+\nabla \mathbf{u}^{T}\right)$ is the rate of deformation tensor and $(\nabla \mathbf{u}) \boldsymbol{\sigma}$ is the the matrix product between $\nabla \mathbf{u}$ and $\boldsymbol{\sigma}$.

This model problem is a simplification of viscoelastic models for polymeric liquids. The first equation corresponds to momentum conservation. The total stress is split into three contributions: the pressure $-p \mathbf{I}$, the stress due to the (Newtonian) solvent $2 \eta_{s} \boldsymbol{\epsilon}(\mathbf{u})$, and the extra-stress $\boldsymbol{\sigma}$ due to the non Newtonian part of the fluid (for instance polymer chains). The third equation is a simplification of the Oldroyd-B constitutive relationship between the extra-stress and the velocity field. For the sake of simplicity, the convective terms in the first and last equations are removed. A future work should take these terms into account.

A theoretical result for Problem (1) with $\eta_{s}=0$ and with convective terms in the first and last equation has been obtained in [23]. Using an iterative procedure, the solution ( $\mathbf{u}, p, \boldsymbol{\sigma}$ ) was proved to be in $H^{3} \times H^{2} \times H^{2}$ provided $\mathbf{f} \in H^{1}$ was small enough and $\partial \Omega$ was sufficiently smooth. In this section, we shall prove that Problem (1) has a solution in $W^{2, r} \times W^{1, r} \times W^{1, r}$ for some $r>2$, provided $\lambda$ is small enough, when $\Omega$ is a convex polygonal domain.

Let $L_{0}^{2}(\Omega)$ be the space of $L^{2}(\Omega)$ functions having zero mean, let $L_{s}^{2}(\Omega)^{4}$ be the space of symmetric tensors having $L^{2}(\Omega)$ components. Given $\mathbf{f} \in H^{-1}(\Omega)^{2}$ and $\mathbf{g} \in L_{s}^{2}(\Omega)^{4}$, we first consider the following problem (we set $\lambda=0$ and add a source term in the third equation of (1)): find $(\mathbf{u}, p, \boldsymbol{\sigma}) \in H_{0}^{1}(\Omega)^{2} \times L_{0}^{2}(\Omega) \times L_{s}^{2}(\Omega)^{4}$ such that

$$
\begin{align*}
-2 \eta_{s} \operatorname{div} \boldsymbol{\epsilon}(\mathbf{u})+\nabla p-\operatorname{div} \boldsymbol{\sigma} & =\mathbf{f} \\
\operatorname{div} \mathbf{u} & =0  \tag{2}\\
\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}-\boldsymbol{\epsilon}(\mathbf{u}) & =\mathbf{g}
\end{align*}
$$

Eliminating $\boldsymbol{\sigma}$ we obtain

$$
\begin{align*}
-2\left(\eta_{s}+\eta_{p}\right) \operatorname{div} \boldsymbol{\epsilon}(\mathbf{u})+\nabla p & =\mathbf{f}+2 \eta_{p} \operatorname{div} \mathbf{g},  \tag{3}\\
\operatorname{div} \mathbf{u} & =0
\end{align*}
$$

Since $\mathbf{g} \in L_{s}^{2}(\Omega)^{4}$, div $\mathbf{g} \in H^{-1}(\Omega)$, and Problem (3) is a classical Stokes problem with solution $(\mathbf{u}, p) \in$ $H_{0}^{1}(\Omega)^{2} \times L_{0}^{2}(\Omega)$. Setting $\boldsymbol{\sigma}=2 \eta_{p}(\boldsymbol{\epsilon}(\mathbf{u})+\mathbf{g})$, then $(\mathbf{u}, p, \boldsymbol{\sigma})$ is solution of Problem (2). Thus we define the
operator

$$
\begin{aligned}
T: H^{-1}(\Omega)^{2} \times L_{s}^{2}(\Omega)^{4} & \longrightarrow H_{0}^{1}(\Omega)^{2} \times L_{s}^{2}(\Omega)^{4} \\
(\mathbf{f}, \mathbf{g}) & \longrightarrow T(\mathbf{f}, \mathbf{g}) \underset{\text { def. }}{=}(\mathbf{u}, \boldsymbol{\sigma}),
\end{aligned}
$$

where $(\mathbf{u}, \boldsymbol{\sigma})$ is the solution of (2). Problem (1) can then be formally written as follows: find $(\mathbf{u}, \boldsymbol{\sigma})$ such that

$$
\begin{equation*}
(\mathbf{u}, \boldsymbol{\sigma})=T\left(\mathbf{f}, \frac{\lambda}{2 \eta_{p}}\left((\nabla \mathbf{u}) \boldsymbol{\sigma}+\boldsymbol{\sigma}(\nabla \mathbf{u})^{T}\right)\right) . \tag{4}
\end{equation*}
$$

The question is now to find the correct spaces for the above problem. Indeed, since $(\mathbf{u}, \boldsymbol{\sigma})$ are a priori only in $H_{0}^{1}(\Omega)^{2} \times L_{s}^{2}(\Omega)^{4}$, then $(\nabla \mathbf{u}) \boldsymbol{\sigma}+\boldsymbol{\sigma}(\nabla \mathbf{u})^{T}$ is only in $L^{1}(\Omega)$, thus $\operatorname{div}\left((\nabla \mathbf{u}) \boldsymbol{\sigma}+\boldsymbol{\sigma}(\nabla \mathbf{u})^{T}\right) \notin H^{-1}(\Omega)!$

Let us recall some regularity properties of Stokes' problem. Let $\mu>0$ and consider the solution $(\mathbf{w}, q) \in$ $H_{0}^{1}(\Omega)^{2} \times L_{0}^{2}(\Omega)$ of

$$
\begin{aligned}
-2 \mu \operatorname{div} \boldsymbol{\epsilon}(\mathbf{w})+\nabla q & =\mathbf{f} \\
\operatorname{div} \mathbf{w} & =0
\end{aligned}
$$

According to Proposition 2.3 of [27], if the boundary $\partial \Omega$ is $\mathcal{C}^{2}$ and if $\mathbf{f} \in L^{r}(\Omega)^{2}$, then $\mathbf{w} \in W^{2, r}(\Omega)^{2}$, for all $1<r<\infty$. According to Theorem 7.3.3.1 of [18] and Proposition 5.3 of [17], if $\Omega$ is a convex polygon, then there exists $r>2\left(r\right.$ depends on the largest angle of the polygon) such that, if $\mathbf{f} \in L^{r}(\Omega)^{2}$, then $\mathbf{w} \in W^{2, r}(\Omega)^{2}$ and $\|\mathbf{w}\|_{W^{2, r}} \leq C\|\mathbf{f}\|_{L^{r}}$. In the sequel, this being crucial for our analysis, we will assume that $\Omega$ is a convex polygon and we will consider $r$ as above.

We set

$$
X=\left(W^{2, r}(\Omega) \cap H_{0}^{1}(\Omega)\right)^{2} \times W_{s}^{1, r}(\Omega)^{4}
$$

where $W_{s}^{1, r}(\Omega)^{4}$ is the space of symmetric tensors having components in $W^{1, r}(\Omega)$. We know that $W^{1, r}(\Omega)$, $r>2$, is an algebra, namely if $\varphi, \psi \in W^{1, r}(\Omega)$ then the product $\varphi \psi \in W^{1, r}(\Omega)$ and there exists $C$ such that

$$
\|\varphi \psi\|_{W^{1, r}(\Omega)} \leq C\|\varphi\|_{W^{1, r}(\Omega)}\|\psi\|_{W^{1, r}(\Omega)} \quad \forall \varphi, \psi \in W^{1, r}(\Omega)
$$

Therefore, if $(\mathbf{u}, \boldsymbol{\sigma}) \in X$ then

$$
\mathbf{g}=\frac{\lambda}{2 \eta_{p}}\left((\nabla \mathbf{u}) \boldsymbol{\sigma}+\boldsymbol{\sigma}(\nabla \mathbf{u})^{T}\right) \in W_{s}^{1, r}(\Omega)^{4}
$$

and $\operatorname{div} \mathbf{g} \in L^{r}(\Omega)^{2}$. Thus, if $\mathbf{f} \in L^{r}(\Omega)^{2}$, then

$$
T\left(\mathbf{f}, \frac{\lambda}{2 \eta_{p}}\left((\nabla \mathbf{u}) \boldsymbol{\sigma}+\boldsymbol{\sigma}(\nabla \mathbf{u})^{T}\right)\right) \in X
$$

From now on, we will restrict $T$ to $L^{r}(\Omega)^{2} \times W_{s}^{1, r}(\Omega)^{4}$ that is

$$
\begin{align*}
T: L^{r}(\Omega)^{2} & \times W_{s}^{1, r}(\Omega)^{4} \\
(\mathbf{f}, \mathbf{g}) & \longrightarrow X  \tag{5}\\
& \longrightarrow T(\mathbf{f}, \mathbf{g}) \underset{\text { def. }}{=}(\mathbf{u}, \boldsymbol{\sigma})
\end{align*}
$$

where $(\mathbf{u}, \boldsymbol{\sigma})$ is the solution of (2). The operator $T$ is bounded since

$$
\|T(\mathbf{f}, \mathbf{g})\|_{X} \leq C\left(\|\mathbf{f}\|_{L^{r}}+\|\mathbf{g}\|_{W^{1, r}}\right) \quad \forall(\mathbf{f}, \mathbf{g}) \in L^{r}(\Omega)^{2} \times W_{s}^{1, r}(\Omega)^{4}
$$

Finally, going back to (4), we are looking for $\mathbf{U}=(\mathbf{u}, \boldsymbol{\sigma}) \in X$ such that

$$
\begin{equation*}
F(\lambda, \mathbf{U})=0 \tag{6}
\end{equation*}
$$

where $F: \mathbf{R} \times X \rightarrow X$ is defined by

$$
F(\lambda, \mathbf{U})=\mathbf{U}-T(\mathbf{f}, \lambda S(\mathbf{U}))
$$

and $S: X \rightarrow W_{s}^{1, r}(\Omega)^{4}$ by

$$
\begin{equation*}
S(\mathbf{U})=\frac{1}{2 \eta_{p}}\left((\nabla \mathbf{u}) \boldsymbol{\sigma}+\boldsymbol{\sigma}(\nabla \mathbf{u})^{T}\right) \tag{7}
\end{equation*}
$$

We have the following result.
Lemma 2.1. For all $\lambda \geq 0$, the operator $F(\lambda, \cdot): X \rightarrow X$ is $\mathcal{C}^{1}$, with Frechet derivative given by

$$
D_{\mathbf{U}} F(\lambda, \mathbf{U}) \mathbf{V}=(\mathbf{v}, \boldsymbol{\tau})-T\left(\mathbf{0}, \frac{\lambda}{2 \eta_{p}}\left((\nabla \mathbf{u}) \boldsymbol{\tau}+\boldsymbol{\tau}(\nabla \mathbf{u})^{T}+(\nabla \mathbf{v}) \boldsymbol{\sigma}+\boldsymbol{\sigma}(\nabla \mathbf{v})^{T}\right)\right)
$$

for all $\mathbf{U}=(\mathbf{u}, \boldsymbol{\sigma}) \in X, \mathbf{V}=(\mathbf{v}, \boldsymbol{\tau})$ in $X$.
Proof. It suffices to note that the operator $S: X \rightarrow W_{s}^{1, r}(\Omega)^{4}$ is $\mathcal{C}^{1}$ with Frechet derivative

$$
D S(\mathbf{U}) \mathbf{V}=(\nabla \mathbf{u}) \boldsymbol{\tau}+\boldsymbol{\tau}(\nabla \mathbf{u})^{T}+(\nabla \mathbf{v}) \boldsymbol{\sigma}+\boldsymbol{\sigma}(\nabla \mathbf{v})^{T}
$$

We now observe that, when $\lambda=0,(6)$ reduces to a classical three fields Stokes' problem. Thus, there is a unique $\mathbf{U}_{0}=\left(\mathbf{u}_{0}, \sigma_{0}\right) \in X$ such that $F\left(0, \mathbf{U}_{0}\right)=0$. Using the implicit function theorem, we can then prove the following result.
Theorem 2.2. There exists $\lambda_{0}>0$ and $\delta>0$ such that, for all $\lambda \leq \lambda_{0}$, there exists a unique $\mathbf{U}(\lambda)=$ $(\mathbf{u}(\lambda), \boldsymbol{\sigma}(\lambda)) \in X$ such that $F(\lambda, \mathbf{U}(\lambda))=0$ and $\left\|\mathbf{U}(\lambda)-\mathbf{U}_{0}\right\|_{X} \leq \delta$. Moreover, the mapping $\lambda \in\left[0, \lambda_{0}\right] \rightarrow$ $\mathbf{U}(\lambda) \in X$ is continuous.
Proof. When $\lambda=0$ there is a unique $\mathbf{U}_{0}=\left(\mathbf{u}_{0}, \boldsymbol{\sigma}_{0}\right) \in X$ such that $F\left(0, \mathbf{U}_{0}\right)=0$. Moreover, using Lemma 2.1, we have $D_{\mathbf{U}} F\left(0, \mathbf{U}_{0}\right)=I$, which is an isomorphism onto $X$. The result is then an immediate consequence of the implicit function theorem.

## 3. A Galerkin Least Square method

In [6], a GLS method with continuous, piecewise linear finite elements was studied for solving (1) in the linear case, i.e. when $\lambda=0$. The method was proved to be stable and convergent even when the solvent viscosity $\eta_{s}$ was small compared to the polymer viscosity $\eta_{p}$. A priori error estimates where derived in the space $W=H_{0}^{1}(\Omega)^{2} \times L_{0}^{2}(\Omega) \times L_{s}^{2}(\Omega)^{4}$. Finally, this GLS method was shown to be equivalent to some EVSS method. The aim of the paper is to extend these results to the nonlinear case, that is when $\lambda \neq 0$.

For any $h>0$, let $\mathcal{T}_{h}$ be a mesh of $\bar{\Omega}$ into triangles $K$ with diameters $h_{K}$ less than $h$. We assume that the mesh satisfies the regularity and inverse assumptions in the sense of [10]. We consider $W_{h}$ the finite dimensional subspace of $W$ consisting in continuous, piecewise linear velocities, pressures, and stresses on the mesh $\mathcal{T}_{h}$. More precisely $W_{h} \subset W$ is defined by $W_{h}=V_{h} \times Q_{h} \times M_{h}$ where

$$
\begin{align*}
V_{h} & =\left\{\mathbf{v} \in \mathcal{C}^{0}(\bar{\Omega})^{2} ; \mathbf{v}_{\mid K} \in\left(\mathbb{P}_{1}\right)^{2}, \forall K \in \mathcal{T}_{h}\right\} \cap H_{0}^{1}(\Omega)^{2}, \\
Q_{h} & =\left\{q \in \mathcal{C}^{0}(\bar{\Omega}) ; q_{\mid K} \in \mathbb{P}_{1}, \forall K \in \mathcal{T}_{h}\right\} \cap L_{0}^{2}(\Omega),  \tag{8}\\
M_{h} & =\left\{\boldsymbol{\tau} \in \mathcal{C}^{0}(\bar{\Omega})^{4} ; \boldsymbol{\tau}_{\mid K} \in\left(\mathbb{P}_{1}\right)^{4}, \forall K \in \mathcal{T}_{h}\right\} \cap L_{s}^{2}(\Omega)^{4}
\end{align*}
$$

In order to approach the solution of (1) we consider the following problem: find $\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h}\right) \in W_{h}$ such that

$$
\begin{align*}
& 2 \eta_{s}\left(\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), \boldsymbol{\epsilon}(\mathbf{v})\right)-\left(p_{h}, \operatorname{div} \mathbf{v}\right)+\left(\boldsymbol{\sigma}_{h}, \boldsymbol{\epsilon}(\mathbf{v})\right)-(\mathbf{f}, \mathbf{v}) \\
& -\left(\operatorname{div} \mathbf{u}_{h}, q\right)-\sum_{K \in \mathcal{T}_{h}} \frac{\alpha h_{K}^{2}}{2 \eta_{p}}\left(-2 \eta_{s} \operatorname{div} \boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right)+\nabla p_{h}-\operatorname{div} \boldsymbol{\sigma}_{h}-\mathbf{f}, \nabla q\right)_{K}  \tag{9}\\
& -\left(\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}-\frac{\lambda}{2 \eta_{p}}\left(\nabla \mathbf{u}_{h} \boldsymbol{\sigma}_{h}+\boldsymbol{\sigma}_{h} \nabla \mathbf{u}_{h}^{T}\right)-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), \boldsymbol{\tau}+2 \eta_{p} \beta \boldsymbol{\epsilon}(\mathbf{v})\right)=0,
\end{align*}
$$

for all $(\mathbf{v}, q, \boldsymbol{\tau}) \in W_{h}$. Here $(\cdot, \cdot)$ denotes the $L^{2}(\Omega)$ scalar product for scalars, vectors and tensors, for instance

$$
(\boldsymbol{\sigma}, \boldsymbol{\tau})=\sum_{i, j=1}^{2} \int_{\Omega} \sigma_{i j} \tau_{i j} \mathrm{~d} x \quad \forall \boldsymbol{\sigma}, \boldsymbol{\tau} \in L^{2}(\Omega)^{4}
$$

Also, $\alpha>0$ and $0<\beta<2$ are dimensionless stabilization parameters. Note that, since the velocity is piecewise linear, then $\operatorname{div} \boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right)$ is zero in each of the mesh triangles.

In the next section, existence of a solution to (9) is proved, as well as a priori error estimates. In Section 5, a posteriori error estimates based on the equation residual are derived. In Section 6, an EVSS method related to the GLS method is presented. Numerical results on both GLS and EVSS schemes are reported in Section 7.

## 4. Existence and a priori Error estimates

Let us consider the discrete counterpart of the operator $T$, namely the operator

$$
\begin{align*}
T_{h}: L^{2}(\Omega)^{2} \times L_{s}^{2}(\Omega)^{4} & \longrightarrow H_{0}^{1}(\Omega)^{2} \times L_{s}^{2}(\Omega)^{4} \\
(\mathbf{f}, \mathbf{g}) & \longrightarrow T_{h}(\mathbf{f}, \mathbf{g}) \underset{\text { def. }}{=}\left(\mathbf{u}_{h}, \boldsymbol{\sigma}_{h}\right) \in V_{h} \times M_{h} \tag{10}
\end{align*}
$$

where $\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h}\right) \in W_{h}$ satisfies

$$
\begin{align*}
& 2 \eta_{s}\left(\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), \boldsymbol{\epsilon}(\mathbf{v})\right)-\left(p_{h}, \operatorname{div} \mathbf{v}\right)+\left(\boldsymbol{\sigma}_{h}, \boldsymbol{\epsilon}(\mathbf{v})\right)-(\mathbf{f}, \mathbf{v}) \\
& -\left(\operatorname{div} \mathbf{u}_{h}, q\right)-\sum_{K \in \mathcal{T}_{h}} \frac{\alpha h_{K}^{2}}{2 \eta_{p}}\left(-2 \eta_{s} \operatorname{div} \boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right)+\nabla p_{h}-\operatorname{div} \boldsymbol{\sigma}_{h}-\mathbf{f}, \nabla q\right)_{K}  \tag{11}\\
& -\left(\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}-\mathbf{g}-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), \boldsymbol{\tau}+2 \eta_{p} \beta \boldsymbol{\epsilon}(\mathbf{v})\right)=0
\end{align*}
$$

for all $(\mathbf{v}, q, \boldsymbol{\tau}) \in W_{h}$. We then have the following result.
Lemma 4.1. The operator $T_{h}$ is well defined and is uniformly bounded with respect to $h$. Moreover, there exists $h_{0}>0$ and $C$ such that, for all $(\mathbf{f}, \mathbf{g}) \in L^{r}(\Omega)^{2} \times W_{s}^{1, r}(\Omega)^{4}$ we have

$$
\left\|T(\mathbf{f}, \mathbf{g})-T_{h}(\mathbf{f}, \mathbf{g})\right\|_{H^{1} \times L^{2}} \leq C h\left(\|\mathbf{f}\|_{L^{r}}+\|\mathbf{g}\|_{W^{1, r}}\right) \quad \forall h \leq h_{0}
$$

Proof. We introduce, as in [6], the bilinear form $B_{h}\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h} ; \mathbf{v}, q, \boldsymbol{\tau}\right)$ corresponding to the left hand side of the weak formulation (11) when $\mathbf{f}$ and $\mathbf{g}$ are zero. The operator $T_{h}$ is then defined by $T_{h}(\mathbf{f}, \mathbf{g})=\left(\mathbf{u}_{h}, \boldsymbol{\sigma}_{h}\right)$ where $\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h}\right) \in W_{h}$ satisfies

$$
\begin{aligned}
B_{h}\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h} ; \mathbf{v}, q, \boldsymbol{\tau}\right)=(\mathbf{f}, \mathbf{v}) & -\left(\mathbf{g}, \boldsymbol{\tau}+2 \eta_{p} \beta \boldsymbol{\epsilon}(\mathbf{v})\right) \\
& -\sum_{K \in \mathcal{T}_{h}} \frac{\alpha h_{K}^{2}}{2 \eta_{p}}(\mathbf{f}, \nabla q)_{K}
\end{aligned}
$$

for all $(\mathbf{v}, q, \boldsymbol{\tau}) \in W_{h}$. It is proved in [6] that $B_{h}: W_{h} \times W_{h} \rightarrow \mathbb{R}$ satisfies the uniform (with respect to $h$ ) inf-sup condition in the norm $\|\cdot\|_{H^{1} \times L^{2} \times L^{2}}$. Thus $T_{h}$ is well defined and we have, for all $h>0$

$$
\left\|\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h}\right\|_{H^{1} \times L^{2} \times L^{2}} \leq C\|\mathbf{f}, \mathbf{g}\|_{L^{2} \times L^{2}}
$$

and thus

$$
\left\|T_{h}\right\|_{\mathcal{L}\left(L^{2} \times L^{2}, H^{1} \times L^{2}\right)} \leq C
$$

where $C$ does not depend on $h$. Moreover, if $(\mathbf{f}, \mathbf{g}) \in L^{r}(\Omega)^{2} \times W_{s}^{1, r}(\Omega)^{4}$, then we can introduce $T(\mathbf{f}, \mathbf{g})=$ $(\mathbf{u}, \boldsymbol{\sigma}) \in X$, where $(\mathbf{u}, \boldsymbol{\sigma})$ is the solution of (2) and where $T$ is defined in (5). We have

$$
\begin{aligned}
&\left\|T(\mathbf{f}, \mathbf{g})-T_{h}(\mathbf{f}, \mathbf{g})\right\|_{H^{1} \times L^{2}}=\left\|\mathbf{u}-\mathbf{u}_{h}, \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h}\right\|_{H^{1} \times L^{2}} \\
& \quad \leq\left\|\mathbf{u}-r_{h} \mathbf{u}, \boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma}\right\|_{H^{1} \times L^{2}}+\left\|r_{h} \mathbf{u}-\mathbf{u}_{h}, r_{h} \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h}\right\|_{H^{1} \times L^{2}}
\end{aligned}
$$

where $r_{h}$ denotes Lagrange interpolant on scalars, vectors or tensors. Note that, considering $r_{h} \mathbf{u}, r_{h} p$ and $r_{h} \boldsymbol{\sigma}$ has a meaning since $(\mathbf{u}, p, \boldsymbol{\sigma}) \in W^{2, r}(\Omega)^{2} \times W^{1, r}(\Omega) \times W_{s}^{1, r}(\Omega)^{4}$, and thus $\mathbf{u}, p$ and $\boldsymbol{\sigma}$ are continuous on $\bar{\Omega}$. Using classical interpolation results [10] together with the well posedness of the operator $T$, we have

$$
\begin{align*}
\left\|\mathbf{u}-r_{h} \mathbf{u}, \boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma}\right\|_{H^{1} \times L^{2}} & \leq C h\left(|\mathbf{u}|_{H^{2}}+|\boldsymbol{\sigma}|_{H^{1} \cap \mathcal{C}^{0}}\right) \\
& \leq C h\left(\|\mathbf{u}\|_{W^{2, r}}+\|\boldsymbol{\sigma}\|_{W^{1, r}}\right)  \tag{12}\\
& \leq C h\left(\|\mathbf{f}\|_{L^{r}}+\|\mathbf{g}\|_{W^{1, r}}\right)
\end{align*}
$$

On the other side, using the fact that $B_{h}$ satisfies the uniform inf-sup condition on $W_{h}$, there is a constant $C$ independent of $h$ such that

$$
\begin{equation*}
\left\|r_{h} \mathbf{u}-\mathbf{u}_{h}, r_{h} \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h}\right\|_{H^{1} \times L^{2}} \leq C \sup _{\mathbf{0} \neq(\mathbf{v}, q, \boldsymbol{\tau}) \in W_{h}} \frac{B_{h}\left(r_{h} \mathbf{u}-\mathbf{u}_{h}, r_{h} p-p_{h}, r_{h} \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h} ; \mathbf{v}, q, \boldsymbol{\tau}\right)}{\|\mathbf{v}, q, \boldsymbol{\tau}\|_{H^{1} \times L^{2} \times L^{2}}} \tag{13}
\end{equation*}
$$

The scheme (11) is consistent in the sense of [14], that is

$$
B_{h}\left(\mathbf{u}-\mathbf{u}_{h}, p-p_{h}, \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h} ; \mathbf{v}, q, \boldsymbol{\tau}\right)=0 \quad \forall(\mathbf{v}, q, \boldsymbol{\tau}) \in W_{h}
$$

and thus

$$
\left\|r_{h} \mathbf{u}-\mathbf{u}_{h}, r_{h} \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h}\right\|_{H^{1} \times L^{2}} \leq C \sup _{\mathbf{0} \neq(\mathbf{v}, q, \boldsymbol{\tau}) \in W_{h}} \frac{B_{h}\left(\mathbf{u}-r_{h} \mathbf{u}, p-r_{h} p, \boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma} ; \mathbf{v}, q, \boldsymbol{\tau}\right)}{\|\mathbf{v}, q, \boldsymbol{\tau}\|_{H^{1} \times L^{2} \times L^{2}}}
$$

From the definition of $B_{h}$ we have, for all $(\mathbf{v}, q, \boldsymbol{\tau}) \in W_{h}$,

$$
\begin{aligned}
B_{h}(\mathbf{u}- & \left.r_{h} \mathbf{u}, p-r_{h} p, \boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma} ; \mathbf{v}, q, \boldsymbol{\tau}\right) \\
= & 2\left(\eta_{s}+\eta_{p} \beta\right)\left(\boldsymbol{\epsilon}\left(\mathbf{u}-r_{h} \mathbf{u}\right), \boldsymbol{\epsilon}(\mathbf{v})\right) \\
& -\left(p-r_{h} p, \operatorname{div} \mathbf{v}\right)+(1-\beta)\left(\boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma}, \boldsymbol{\epsilon}(\mathbf{v})\right) \\
& -\left(\operatorname{div}\left(\mathbf{u}-r_{h} \mathbf{u}\right), q\right)-\sum_{K \in \mathcal{T}_{h}} \frac{\alpha h_{K}^{2}}{2 \eta_{p}}\left(\mathbf{f}-\nabla r_{h} p+\operatorname{div} r_{h} \boldsymbol{\sigma}, \nabla q\right)_{K} \\
& -\frac{1}{2 \eta_{p}}\left(\boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma}, \boldsymbol{\tau}\right)+\left(\boldsymbol{\epsilon}\left(\mathbf{u}-r_{h} \mathbf{u}\right), \boldsymbol{\tau}\right)
\end{aligned}
$$

Using Cauchy-Schwarz and Young inequalities, there is a constant $C$ independent of $h$ such that

$$
\begin{aligned}
&\left|B_{h}\left(\mathbf{u}-r_{h} \mathbf{u}, p-r_{h} p, \boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma} ; \mathbf{v}, q, \boldsymbol{\tau}\right)\right| \\
& \leq C\left(\left\|\mathbf{u}-r_{h} \mathbf{u}, p-r_{h} p, \boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma}\right\|_{H^{1} \times L^{2} \times L^{2}}\|\mathbf{v}, q, \boldsymbol{\tau}\|_{H^{1} \times L^{2} \times L^{2}}\right. \\
&\left.+\sum_{K \in \mathcal{T}_{h}} h_{K}^{2}\left(\|\mathbf{f}\|_{L^{2}(K)}+\left\|r_{h} p\right\|_{H^{1}(K)}+\left\|r_{h} \boldsymbol{\sigma}\right\|_{H^{1}(K)}\right)\|\nabla q\|_{L^{2}(K)}\right) .
\end{aligned}
$$

We now use the inverse inequality $h_{K}\|\nabla q\|_{L^{2}(K)} \leq C\|q\|_{L^{2}(K)}$, the fact that

$$
\begin{aligned}
& \left\|r_{h} p\right\|_{H^{1}} \leq C\left\|r_{h} p\right\|_{W^{1, r}} \leq \tilde{C}\|p\|_{W^{1, r}} \\
& \left\|r_{h} \boldsymbol{\sigma}\right\|_{H^{1}} \leq C\left\|r_{h} \boldsymbol{\sigma}\right\|_{W^{1, r}} \leq \tilde{C}\|\boldsymbol{\sigma}\|_{W^{1, r}}
\end{aligned}
$$

for $h$ sufficiently small, and standard interpolation estimates [10] to obtain

$$
\begin{aligned}
& \left|B_{h}\left(\mathbf{u}-r_{h} \mathbf{u}, p-r_{h} p, \boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma} ; \mathbf{v}, q, \boldsymbol{\tau}\right)\right| \\
& \leq \\
& \quad C h\left(\|\mathbf{u}, p, \boldsymbol{\sigma}\|_{W^{2, r} \times W^{1, r} \times W^{1, r}}\|\mathbf{v}, q, \boldsymbol{\tau}\|_{H^{1} \times L^{2} \times L^{2}}\right. \\
& \left.\quad+\left(\|\mathbf{f}\|_{L^{r}}+\|p\|_{W^{1, r}}+\|\boldsymbol{\sigma}\|_{W^{1, r}}\right)\|q\|_{L^{2}}\right)
\end{aligned}
$$

Using the well posedness of the operator $T$ we obtain

$$
\left|B_{h}\left(\mathbf{u}-r_{h} \mathbf{u}, p-r_{h} p, \boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma} ; \mathbf{v}, q, \boldsymbol{\tau}\right)\right| \leq C h\left(\|\mathbf{f}\|_{L^{r}}+\|\mathbf{g}\|_{W^{1, r}}\right)\|\mathbf{v}, q, \boldsymbol{\tau}\|_{H^{1} \times L^{2} \times L^{2}}
$$

for all $(\mathbf{v}, q, \boldsymbol{\tau}) \in W_{h}$. The above inequality in (13), together with (12) finally yields

$$
\left\|\mathbf{u}-\mathbf{u}_{h}, \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h}\right\|_{H^{1} \times L^{2}} \leq C h\left(\|\mathbf{f}\|_{L^{r}}+\|\mathbf{g}\|_{W^{1, r}}\right)
$$

which is the desired result.
We now would like to rewrite the nonlinear GLS scheme (9) in an abstract framework, as we did in (6) for the continuous Problem (1). We now introduce $X_{h}=V_{h} \times M_{h}$ equipped with the $\|\cdot\|_{H^{1} \times L^{2}}$ norm and we want to rewrite the nonlinear GLS scheme (9) as

$$
\begin{equation*}
F_{h}\left(\lambda, \mathbf{U}_{h}\right)=0 \quad \text { with } \quad \mathbf{U}_{h}=\left(\mathbf{u}_{h}, \boldsymbol{\sigma}_{h}\right) \tag{14}
\end{equation*}
$$

Here $F_{h}: \mathbf{R} \times X_{h} \rightarrow X_{h}$ is defined by

$$
F_{h}\left(\lambda, \mathbf{U}_{h}\right)=\mathbf{U}_{h}-T_{h}\left(\mathbf{f}, \lambda S\left(\mathbf{U}_{h}\right)\right)
$$

where $S$ is still formally defined by

$$
S\left(\mathbf{U}_{h}\right)=S\left(\mathbf{u}_{h}, \boldsymbol{\sigma}_{h}\right)=\frac{1}{2 \eta_{p}}\left(\left(\nabla \mathbf{u}_{h}\right) \boldsymbol{\sigma}_{h}+\boldsymbol{\sigma}_{h}\left(\nabla \mathbf{u}_{h}\right)^{T}\right)
$$

However, since $X_{h} \not \subset X$, we have to extend $S$ on $W_{0}^{1, r}(\Omega)^{2} \times W_{s}^{1, r}(\Omega)^{4}$, and since $W_{s}^{1, r}(\Omega)^{4} \subset \mathcal{C}^{0}(\bar{\Omega})$, we can consider $S: W_{0}^{1, r}(\Omega)^{2} \times W_{s}^{1, r}(\Omega)^{4} \rightarrow L_{s}^{r}(\Omega)^{4}$. Clearly, this operator $S$ is $\mathcal{C}^{1}$.

We have the following result.

Lemma 4.2. Let $\mathbf{U}(\lambda), 0 \leq \lambda \leq \lambda_{0}$, be given by Theorem 2.2. There exists $C$ such that, for all $h>0$, for all $\lambda \leq \lambda_{0}$ we have

$$
\begin{align*}
& \left\|F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)\right\|_{H^{1} \times L^{2}} \leq C h  \tag{15}\\
& \left\|D_{\mathbf{U}} F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)-D_{\mathbf{U}} F_{h}(\lambda, \mathbf{V})\right\|_{\mathcal{L}\left(X_{h}\right)} \leq C \frac{\lambda}{h}\left\|r_{h} \mathbf{U}(\lambda)-\mathbf{V}\right\|_{H^{1} \times L^{2}}, \tag{16}
\end{align*}
$$

for all $\mathbf{V} \in X_{h}$.
Proof. From the definition of $F_{h}$ and since $F(\lambda, \mathbf{U}(\lambda))=0$, we have

$$
F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)=r_{h} \mathbf{U}(\lambda)-\mathbf{U}(\lambda)-T_{h}\left(\mathbf{f}, \lambda S\left(r_{h} \mathbf{U}(\lambda)\right)\right)+T(\mathbf{f}, \lambda S(\mathbf{U}(\lambda)))
$$

so that

$$
\begin{aligned}
&\left\|F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)\right\|_{H^{1} \times L^{2}} \\
& \leq\left\|\mathbf{U}(\lambda)-r_{h} \mathbf{U}(\lambda)\right\|_{H^{1} \times L^{2}}+\left\|\left(T-T_{h}\right)(\mathbf{f}, \lambda S(\mathbf{U}(\lambda)))\right\|_{H^{1} \times L^{2}} \\
&+\left\|T_{h}\left(\mathbf{0}, \lambda\left[S(\mathbf{U}(\lambda))-S\left(r_{h} \mathbf{U}(\lambda)\right)\right]\right)\right\|_{H^{1} \times L^{2}}
\end{aligned}
$$

Using Lemma 4.1 for the second and third terms in the right hand side of the above inequality together with interpolation results, we obtain

$$
\begin{align*}
& \left\|F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)\right\|_{H^{1} \times L^{2}} \\
& \quad \leq \quad C h\|\mathbf{U}(\lambda)\|_{X}+C h\left(\|\mathbf{f}\|_{L^{r}}+\lambda\|S(\mathbf{U}(\lambda))\|_{W^{1, r}}\right)  \tag{17}\\
& \quad+C \lambda\left\|S(\mathbf{U}(\lambda))-S\left(r_{h} \mathbf{U}(\lambda)\right)\right\|_{L^{2}}
\end{align*}
$$

$C$ being independent of $h$ and $\lambda \in\left[0, \lambda_{0}\right]$. A simple calculation shows that

$$
\begin{equation*}
\|S(\mathbf{U}(\lambda))\|_{W^{1, r}} \leq C\|\mathbf{U}(\lambda)\|_{X}^{2} \tag{18}
\end{equation*}
$$

$C$ being independent of $h$ and $\lambda \in\left[0, \lambda_{0}\right]$. On the other hand, we also have

$$
\begin{aligned}
& 2 \eta_{p}\left(S(\mathbf{U})-S\left(r_{h} \mathbf{U}\right)\right) \\
& \quad=\nabla \mathbf{u} \boldsymbol{\sigma}+\boldsymbol{\sigma} \nabla \mathbf{u}^{T}-\left(\nabla r_{h} \mathbf{u}\right) r_{h} \boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma}\left(\nabla r_{h} \mathbf{u}\right)^{T} \\
& \quad=\nabla\left(\mathbf{u}-r_{h} \mathbf{u}\right) \boldsymbol{\sigma}+\left(\nabla r_{h} \mathbf{u}\right)\left(\boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma}\right)+\boldsymbol{\sigma} \nabla\left(\mathbf{u}-r_{h} \mathbf{u}\right)^{T}+\left(\boldsymbol{\sigma}-r_{h} \boldsymbol{\sigma}\right)\left(\nabla r_{h} \mathbf{u}\right)^{T}
\end{aligned}
$$

so that

$$
\begin{align*}
\left\|S(\mathbf{U}(\lambda))-S\left(r_{h} \mathbf{U}(\lambda)\right)\right\|_{L^{2}} & \leq C\left\|\mathbf{U}(\lambda)-r_{h} \mathbf{U}(\lambda)\right\|_{H^{1} \times L^{2}}\|\mathbf{U}(\lambda)\|_{X} \\
& \leq C h\|\mathbf{U}(\lambda)\|_{X}^{2} \tag{19}
\end{align*}
$$

$C$ being independent of $h$ and $\lambda$. Finally (19) and (18) in (17) yields (15) for $\lambda \in\left[0, \lambda_{0}\right]$.
Let us now prove (16). For all $\mathbf{V}=(\mathbf{v}, \boldsymbol{\tau}) \in X_{h}$, for all $\mathbf{W}=(\mathbf{w}, \boldsymbol{\gamma}) \in X_{h}$, we have

$$
\left(D_{\mathbf{U}} F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)-D_{\mathbf{U}} F_{h}(\lambda, \mathbf{V})\right) \mathbf{W}=-T_{h}\left(\mathbf{0}, \lambda\left[D S\left(r_{h} \mathbf{U}(\lambda)\right) \mathbf{W}-D S(\mathbf{V}) \mathbf{W}\right]\right)
$$

Using Lemma 4.1 we obtain

$$
\left\|\left(D_{\mathbf{U}} F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)-D_{\mathbf{U}} F_{h}(\lambda, \mathbf{V})\right) \mathbf{W}\right\|_{H^{1} \times L^{2}} \leq C \lambda\left\|\left(D S\left(r_{h} \mathbf{U}(\lambda)\right)-D S(\mathbf{V})\right) \mathbf{W}\right\|_{L^{2}}
$$

We have

$$
\begin{aligned}
& \left\|\left(D S\left(r_{h} \mathbf{U}\right)-D S(\mathbf{V})\right) \mathbf{W}\right\|_{L^{2}} \\
& \quad=\frac{1}{2 \eta_{p}}\left\|\nabla\left(r_{h} \mathbf{u}-\mathbf{v}\right) \boldsymbol{\gamma}+\gamma \nabla\left(r_{h} \mathbf{u}-\mathbf{v}\right)^{T}+\nabla \mathbf{w}\left(r_{h} \boldsymbol{\sigma}-\boldsymbol{\tau}\right)+\left(r_{h} \boldsymbol{\sigma}-\boldsymbol{\tau}\right) \nabla \mathbf{w}^{T}\right\|_{L^{2}} \\
& \quad \leq \frac{1}{\eta_{p}}\left(\left\|\nabla\left(r_{h} \mathbf{u}-\mathbf{v}\right)\right\|_{L^{r}}\|\gamma\|_{L^{q}}+\|\nabla \mathbf{w}\|_{L^{r}}\left\|r_{h} \boldsymbol{\sigma}-\boldsymbol{\tau}\right\|_{L^{q}}\right)
\end{aligned}
$$

with $q=\frac{2 r}{r-2}$. For $p>2$, the following inverse inequalities hold

$$
\|v\|_{L^{p}} \leq \frac{C}{h^{\frac{p-2}{p}}}\|v\|_{L^{2}}, \quad\|\nabla v\|_{L^{p}} \leq \frac{C}{h^{\frac{p-2}{p}}}\|\nabla v\|_{L^{2}}
$$

for all continuous, piecewise linear function $v$, see [10]. This yields

$$
\left\|\left(D S\left(r_{h} \mathbf{U}\right)-D S(\mathbf{V})\right) \mathbf{W}\right\|_{L^{2}} \leq C \frac{1}{h}\left\|r_{h} \mathbf{U}-\mathbf{V}\right\|_{H^{1} \times L^{2}}\|\mathbf{W}\|_{H^{1} \times L^{2}}
$$

$C$ being independent of $\lambda$ and $h$. This last inequality yields (16).
Before proving existence of a solution to (14) we still need to check that $D_{\mathbf{U}} F_{h}$ is invertible is the neighbourhood of $\mathbf{U}(\lambda)$.
Lemma 4.3. Let $\mathbf{U}(\lambda), 0 \leq \lambda \leq \lambda_{0}$, be given by Theorem 2.2. There exists $0<\lambda_{1} \leq \lambda_{0}$ such that, for all $\lambda \leq \lambda_{1}$ and for all $h$ we have

$$
\left\|D_{\mathbf{U}} F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)^{-1}\right\|_{\mathcal{L}\left(X_{h}\right)} \leq 2
$$

Proof. By definition of $F_{h}$, we have

$$
D_{\mathbf{U}} F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)=I-T_{h}\left(\mathbf{0}, \lambda D S\left(r_{h} \mathbf{U}(\lambda)\right)\right)
$$

where $D S\left(r_{h} \mathbf{U}(\lambda)\right)$ is such that

$$
D S\left(r_{h} \mathbf{U}\right) \mathbf{V}=\frac{1}{2 \eta_{p}}\left(\nabla\left(r_{h} \mathbf{u}\right) \boldsymbol{\tau}+\boldsymbol{\tau} \nabla\left(r_{h} \mathbf{u}\right)^{T}+\nabla \mathbf{v}\left(r_{h} \boldsymbol{\sigma}\right)+\left(r_{h} \boldsymbol{\sigma}\right) \nabla \mathbf{v}^{T}\right)
$$

for all $\mathbf{V}=(\mathbf{v}, \boldsymbol{\tau}) \in X_{h}$. Thus we can write

$$
D_{\mathbf{U}} F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)=I-G_{h} \quad \text { with } \quad G_{h}=T_{h}\left(\mathbf{0}, \lambda D S\left(r_{h} \mathbf{U}(\lambda)\right)\right)
$$

If we prove that $\left\|G_{h}\right\|_{\mathcal{L}\left(X_{h}\right)} \leq 1 / 2$ for sufficiently small values of $\lambda$, then $D_{\mathbf{U}} F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)$ is invertible and $\left\|D_{\mathbf{U}} F_{h}\left(\lambda, r_{h} \mathbf{U}(\lambda)\right)^{-1}\right\|_{\mathcal{L}\left(X_{h}\right)} \leq 2$. Using Lemma 4.1, it suffices to prove that $D S\left(r_{h} \mathbf{U}(\lambda)\right): X_{h} \rightarrow L^{2}(\Omega)^{4}$ is uniformly bounded with respect to $h$. Proceeding as in the proof of the previous Lemma we have

$$
\begin{aligned}
\left\|D S\left(r_{h} \mathbf{U}\right) \mathbf{V}\right\|_{L^{2}} & =\frac{1}{2 \eta_{p}}\left\|\nabla\left(r_{h} \mathbf{u}\right) \boldsymbol{\tau}+\boldsymbol{\tau} \nabla\left(r_{h} \mathbf{u}\right)^{T}+\nabla \mathbf{v}\left(r_{h} \boldsymbol{\sigma}\right)+\left(r_{h} \boldsymbol{\sigma}\right) \nabla \mathbf{v}^{T}\right\|_{L^{2}} \\
& \leq C\left(\|\nabla \mathbf{u}\|_{L^{\infty}}\|\boldsymbol{\tau}\|_{L^{2}}+\|\nabla \mathbf{v}\|_{L^{2}}\|\boldsymbol{\sigma}\|_{L^{\infty}}\right) \\
& \leq \tilde{C}\left(\|\mathbf{u}\|_{W^{2, r}}\|\boldsymbol{\tau}\|_{L^{2}}+\|\nabla \mathbf{v}\|_{L^{2}}\|\boldsymbol{\sigma}\|_{W^{1, r}}\right)
\end{aligned}
$$

$C, \tilde{C}$ being independent of $\lambda$ and $h$. In other words we have

$$
\left\|D S\left(r_{h} \mathbf{U}(\lambda)\right) \mathbf{V}\right\|_{L^{2}} \leq C\|\mathbf{U}(\lambda)\|_{X}\|\mathbf{V}\|_{H^{1} \times L^{2}}
$$

and, going back to $G_{h}$ we obtain

$$
\left\|G_{h}(\mathbf{V})\right\|_{H^{1} \times L^{2}} \leq \tilde{C} \lambda\|\mathbf{U}(\lambda)\|_{X}\|\mathbf{V}\|_{H^{1} \times L^{2}}
$$

Finally, for $\lambda$ sufficiently small we obtain $\left\|G_{h}\right\|_{\mathcal{L}\left(X_{h}\right)} \leq 1 / 2$ and the result follows.
We are now in position to state the main result of the section
Theorem 4.4. Let $\mathbf{U}(\lambda), 0 \leq \lambda \leq \lambda_{0}$, be given by Theorem 2.2. There exists $0<\bar{\lambda} \leq \lambda_{0}, \bar{h}$ and $\delta>0$ such that, for all $0 \leq \lambda \leq \bar{\lambda}$, for all $0<h \leq \bar{h}$, there exists a unique $\mathbf{U}_{h}(\lambda)$ in the ball of $X_{h}$ centered at $r_{h} \mathbf{U}(\lambda)$ with radius $\delta h$ in the norm $H^{1} \times L^{2}$, satisfying

$$
F_{h}\left(\lambda, \mathbf{U}_{h}(\lambda)\right)=0 .
$$

Moreover, the mapping $\lambda \in[0, \bar{\lambda}] \rightarrow \mathbf{U}_{h}(\lambda) \in X_{h}$ is continuous and there exists $C>0$ such that the following a priori error estimate holds

$$
\left\|\mathbf{U}(\lambda)-\mathbf{U}_{h}(\lambda)\right\|_{H^{1} \times L^{2}} \leq C h \quad \forall \lambda \leq \bar{\lambda} \quad \forall h \leq \bar{h}
$$

Remark 4.5. The statement of the above existence result is similar (although not the same) to those of [4,21], in which the convective term in the extra-stress constitutive equation are considered (see also [26] for analogous results on a second-grade fluid). However, a different technique is used in this paper, allowing a priori and a posteriori error estimates to be obtained more easily, with other assumptions.

In order to prove this theorem, we use the following abstract result.
Theorem 4.6 (Th. 2.1 of [9]). Let $Y$ and $Z$ be two real Banach spaces with norms $\|\cdot\|_{Y}$ and $\|\cdot\|_{Z}$ respectively. Let $G: Y \rightarrow Z$ be a $\mathcal{C}^{1}$ mapping and $v \in Y$ be such that $D G(v) \in \mathcal{L}(Y ; Z)$ is an isomorphism. We introduce the notations

$$
\begin{aligned}
\epsilon & =\|G(v)\|_{Z} \\
\gamma & =\left\|D G(v)^{-1}\right\|_{\mathcal{L}(Z ; Y)} \\
L(\alpha) & =\sup _{x \in \bar{B}(v, \alpha)}\|D G(v)-D G(x)\|_{\mathcal{L}(Y ; Z)} \\
& \text { with } \bar{B}(v, \alpha)=\left\{y \in Y ;\|v-y\|_{Y} \leq \alpha\right\}
\end{aligned}
$$

and we are interested in finding $u \in Y$ such that

$$
\begin{equation*}
G(u)=0 . \tag{20}
\end{equation*}
$$

We assume that $2 \gamma L(2 \gamma \epsilon) \leq 1$. Then Problem (20) has a unique solution $u$ in the ball $\bar{B}(v, 2 \gamma \epsilon)$ and, for all $x \in \bar{B}(v, 2 \gamma \epsilon)$, we have

$$
\begin{equation*}
\|x-u\|_{Y} \leq 2 \gamma\|G(x)\|_{Z} \tag{21}
\end{equation*}
$$

Proof of Theorem 4.4. We apply Theorem 4.6 with $Y=X_{h}, Z=X_{h}, G=F_{h}, v=r_{h} \mathbf{U}(\lambda)$ and the norm $\|\cdot\|_{H^{1} \times L^{2}}$ in $X_{h}$. The mapping $G: Y \rightarrow Z$ is $\mathcal{C}^{1}$ and, according to Lemma 4.2, for $\lambda$ sufficiently small there is a constant $C_{1}$ independent of $\lambda$ and $h$ such that $\epsilon \leq C_{1} h$. According to Lemma 4.3, for $\lambda$ sufficiently small $\gamma \leq 2$.

According to Lemma 4.2, there is a constant $C_{2}$ independent of $\lambda$ and $h$ such that $L(\alpha) \leq C_{2} \alpha \lambda / h$. Thus, we have

$$
2 \gamma L(2 \gamma \epsilon) \leq 2.2 C_{2}\left(2.2 . C_{1} h\right) \frac{\lambda}{h}=16 C_{1} C_{2} \lambda
$$

Thus, for $\lambda$ sufficiently small $2 \gamma L(2 \gamma \epsilon) \leq 1$ and Theorem 4.6 applies. There exists a unique $\mathbf{U}_{h}(\lambda)$ in the ball $\bar{B}(v, 2 \gamma \epsilon)$ such that $F_{h}\left(\lambda, \mathbf{U}_{h}(\lambda)\right)=0$ and we have

$$
\left\|r_{h} \mathbf{U}(\lambda)-\mathbf{U}_{h}(\lambda)\right\|_{H^{1} \times L^{2}} \leq 4 C_{1} h
$$

It suffices to use the triangle inequality

$$
\left\|\mathbf{U}(\lambda)-\mathbf{U}_{h}(\lambda)\right\|_{H^{1} \times L^{2}} \leq\left\|\mathbf{U}(\lambda)-r_{h} \mathbf{U}(\lambda)\right\|_{H^{1} \times L^{2}}+\left\|r_{h} \mathbf{U}(\lambda)-\mathbf{U}_{h}(\lambda)\right\|_{H^{1} \times L^{2}}
$$

and standard interpolation [10] results to obtain the a priori estimates. The fact that the mapping $\lambda \rightarrow \mathbf{U}_{h}(\lambda)$ is continuous is a direct consequence of the implicit function theorem.

## 5. A POSTERIORI ERROR ESTIMATES

Let us consider again the operator

$$
\begin{aligned}
T_{h}: L^{2}(\Omega)^{2} \times L_{s}^{2}(\Omega)^{4} & \longrightarrow H_{0}^{1}(\Omega)^{2} \times L_{s}^{2}(\Omega)^{4} \\
(\mathbf{f}, \mathbf{g}) & \longrightarrow T_{h}(\mathbf{f}, \mathbf{g}) \underset{\text { def. }}{=}\left(\mathbf{u}_{h}, \boldsymbol{\sigma}_{h}\right) \in V_{h} \times M_{h}
\end{aligned}
$$

where $\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h}\right) \in W_{h}$ satisfies (11). We now introduce a residual based error estimator for $T_{h}$. The notations are those of [2]. For any triangle $K$ of the triangulation $\mathcal{T}_{h}$, let $E_{K}$ be the set of its three edges. For each interior edge $\ell$ of $\mathcal{T}_{h}$, let us choose an arbitrary normal direction $\mathbf{n}$, let $[\cdot]_{\ell}$ denote the jump across edge $\ell$. For each edge $\ell$ of $\mathcal{T}_{h}$ lying on the boundary $\partial \Omega$, we set $[\cdot]_{\ell}=0$. The local error estimator corresponding to (11) is then defined by

$$
\begin{aligned}
\mu_{K}^{2}(\mathbf{f}, \mathbf{g})= & \frac{1}{\eta_{s}+\eta_{p}}\left(h_{K}^{2}\left\|-2 \eta_{s} \operatorname{div} \boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right)+\nabla p_{h}-\operatorname{div} \boldsymbol{\sigma}_{h}-\mathbf{f}\right\|_{L^{2}(K)}^{2}\right. \\
& \left.+\frac{1}{2} \sum_{\ell \in E_{K}}|\ell|\left\|\left[2 \eta_{s} \boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right) \mathbf{n}\right]_{\ell}\right\|_{L^{2}(\ell)}^{2}\right)+\left(\eta_{s}+\eta_{p}\right)\left\|\operatorname{div} \mathbf{u}_{h}\right\|_{L^{2}(K)}^{2} \\
& +\eta_{p}\left\|\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}-\mathbf{g}-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right)\right\|_{L^{2}(K)}^{2}
\end{aligned}
$$

We have the following a posteriori error estimate for operator $T-T_{h}$.
Lemma 5.1. There exists $C$ and $h_{0}>0$ such that, for all $(\mathbf{f}, \mathbf{g}) \in L^{r}(\Omega)^{2} \times W_{s}^{1, r}(\Omega)^{4}$ we have

$$
\left\|T(\mathbf{f}, \mathbf{g})-T_{h}(\mathbf{f}, \mathbf{g})\right\|_{H^{1} \times L^{2}} \leq C\left(\sum_{K \in \mathcal{T}_{h}} \mu_{K}^{2}(\mathbf{f}, \mathbf{g})\right)^{1 / 2} \quad \forall h \leq h_{0}
$$

Proof. We set $T(\mathbf{f}, \mathbf{g})=(\mathbf{u}, \boldsymbol{\sigma})$ where $(\mathbf{u}, p, \boldsymbol{\sigma})$ is the solution of $(2)$, we also set $T_{h}(\mathbf{f}, \mathbf{g})=\left(\mathbf{u}_{h}, \boldsymbol{\sigma}_{h}\right)$ where $\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h}\right)$ is the solution of (11). We recall that $W=H_{0}^{1}(\Omega)^{2} \times L_{0}^{2}(\Omega) \times L_{s}^{2}(\Omega)^{4}$. Let $B: W \times W \rightarrow \mathbb{R}$ be the bilinear form corresponding to the weak formulation of (2), namely

$$
\begin{aligned}
B(\mathbf{u}, p, \boldsymbol{\sigma} ; \mathbf{v}, q, \boldsymbol{\tau})= & 2 \eta_{s}(\boldsymbol{\epsilon}(\mathbf{u}), \boldsymbol{\epsilon}(\mathbf{v}))-(p, \operatorname{div} \mathbf{v})+(\boldsymbol{\sigma}, \boldsymbol{\epsilon}(\mathbf{v})) \\
& -(\operatorname{div} \mathbf{u}, q)-\frac{1}{2 \eta_{p}}(\boldsymbol{\sigma}, \boldsymbol{\tau})+(\boldsymbol{\epsilon}(\mathbf{u}), \boldsymbol{\tau})
\end{aligned}
$$

for all $(\mathbf{u}, p, \boldsymbol{\sigma})$ and $(\mathbf{v}, q, \boldsymbol{\tau})$ in $W$. In Lemma 2 of $[6]$, it is proved that $B$ satisfies an inf-sup condition, therefore we have

$$
\left\|\mathbf{u}-\mathbf{u}_{h}, p-p_{h}, \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h}\right\|_{H^{1} \times L^{2} \times L^{2}} \leq C \sup _{\mathbf{0} \neq(\mathbf{v}, q, \boldsymbol{\tau}) \in W} \frac{B\left(\mathbf{u}-\mathbf{u}_{h}, p-p_{h}, \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h} ; \mathbf{v}, q, \boldsymbol{\tau}\right)}{\|\mathbf{v}, q, \boldsymbol{\tau}\|_{H^{1} \times L^{2} \times L^{2}}}
$$

$C$ being independent of $h$. For all $(\mathbf{v}, q, \boldsymbol{\tau}) \in W$ we have

$$
B\left(\mathbf{u}-\mathbf{u}_{h}, p-p_{h}, \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h} ; \mathbf{v}, q, \boldsymbol{\tau}\right)=(\mathbf{f}, \mathbf{v})-(\mathbf{g}, \boldsymbol{\tau})-B\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h} ; \mathbf{v}, q, \boldsymbol{\tau}\right)
$$

Introducing $B_{h}: W_{h} \times W_{h} \rightarrow \mathbb{R}$ as in the proof of Lemma 4.1, we have

$$
\begin{aligned}
B\left(\mathbf{u}-\mathbf{u}_{h}, p-p_{h}, \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h} ; \mathbf{v}, q, \boldsymbol{\tau}\right)= & \left(\mathbf{f}, \mathbf{v}-\mathbf{v}_{h}\right)-\left(\mathbf{g}, \boldsymbol{\tau}-\boldsymbol{\tau}_{h}\right)-B\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h} ; \mathbf{v}-\mathbf{v}_{h}, q-q_{h}, \boldsymbol{\tau}-\boldsymbol{\tau}_{h}\right) \\
& +\left(B_{h}-B\right)\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h} ; \mathbf{v}_{h}, q_{h}, \boldsymbol{\tau}_{h}\right)
\end{aligned}
$$

for all $\left(\mathbf{v}_{h}, q_{h}, \boldsymbol{\tau}_{h}\right) \in W_{h}$, that is:

$$
\begin{aligned}
& B(\mathbf{u}\left.-\mathbf{u}_{h}, p-p_{h}, \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h} ; \mathbf{v}, q, \boldsymbol{\tau}\right) \\
&=-2 \eta_{s}\left(\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), \boldsymbol{\epsilon}\left(\mathbf{v}-\mathbf{v}_{h}\right)\right)+\left(p_{h}, \operatorname{div}\left(\mathbf{v}-\mathbf{v}_{h}\right)\right)-\left(\boldsymbol{\sigma}_{h}, \boldsymbol{\epsilon}\left(\mathbf{v}-\mathbf{v}_{h}\right)\right)+\left(\mathbf{f}, \mathbf{v}-\mathbf{v}_{h}\right) \\
& \quad+\left(\operatorname{div} \mathbf{u}_{h}, q-q_{h}\right)-\sum_{K \in \mathcal{T}_{h}} \frac{\alpha h_{K}^{2}}{2 \eta_{p}}\left(-2 \eta_{s} \operatorname{div} \boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right)+\nabla p_{h}-\operatorname{div} \boldsymbol{\sigma}_{h}-\mathbf{f}, \nabla q_{h}\right)_{K} \\
&+\left(\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}-\mathbf{g}-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), \boldsymbol{\tau}-\boldsymbol{\tau}_{h}-2 \eta_{p} \beta \boldsymbol{\epsilon}\left(\mathbf{v}_{h}\right)\right) .
\end{aligned}
$$

We then proceed as in [3,28], integrate by parts on each triangle $K \in \mathcal{T}_{h}$ the first three terms in the right hand side of the above equation, and choose $\mathbf{v}_{h}=\mathbf{R}_{h} \mathbf{v}$ (where $R_{h}$ is Clément's interpolant [11]), $q_{h}=0, \boldsymbol{\tau}_{h}=0$ to conclude.

We are now in position to state a posteriori error estimates for the solution to (9). Let us briefly recall the notations. Problem (1) is written as $F(\lambda, \mathbf{U})=0$ with $\mathbf{U}=(\mathbf{u}, \boldsymbol{\sigma}) \in X$ and $F(\lambda, \mathbf{U})=\mathbf{U}-T(\mathbf{f}, \lambda S(\mathbf{U}))$, the operators $T$ and $S$ being defined in (5) and (7). Problem (9) is written as $F_{h}\left(\lambda, \mathbf{U}_{h}\right)=0$ with $\mathbf{U}_{h}=\left(\mathbf{u}_{h}, \boldsymbol{\sigma}_{h}\right) \in$ $X_{h} \not \subset X$ and $F_{h}\left(\lambda, \mathbf{U}_{h}\right)=\mathbf{U}_{h}-T_{h}\left(\mathbf{f}, \lambda S\left(\mathbf{U}_{h}\right)\right)$, the operator $T_{h}$ being defined in (10). According to Theorem 2.2, for $\lambda$ sufficiently small, there is a unique $\mathbf{U}(\lambda)$ such that $F(\lambda, \mathbf{U}(\lambda))=0$. According to Theorem 4.4, for $h$ and $\lambda$ sufficiently small, there exists a unique $\mathbf{U}_{h}(\lambda)$ in a neighbourhood of $r_{h} \mathbf{U}(\lambda)$ depending on $h$ (in the norm $\left.\|\cdot\|_{H^{1} \times L^{2}}\right)$ such that $F_{h}\left(\lambda, \mathbf{U}_{h}(\lambda)\right)=0$. Moreover, when $h$ goes to zero, $\mathbf{U}_{h}(\lambda)$ converges to $\mathbf{U}(\lambda)$ in the norm $\|\cdot\|_{H^{1} \times L^{2}}$.

Theorem 5.2. There exists $\lambda_{0}, h_{0}$ and $C>0$ such that, for all $\lambda \leq \lambda_{0}$, for all $h \leq h_{0}$,

$$
\begin{equation*}
\left\|\mathbf{U}(\lambda)-\mathbf{U}_{h}(\lambda)\right\|_{H^{1} \times L^{2}} \leq C\left(\sum_{K \in \mathcal{T}_{h}} \mu_{K}^{2}\left(\mathbf{f}, \lambda S\left(\mathbf{U}_{h}(\lambda)\right)\right)\right)^{1 / 2} \tag{22}
\end{equation*}
$$

Proof. Using the definition of $F$ and $F_{h}$ we have

$$
\begin{align*}
\mathbf{U}-\mathbf{U}_{h} & =T(\mathbf{f}, \lambda S(\mathbf{U}))-T_{h}\left(\mathbf{f}, \lambda S\left(\mathbf{U}_{h}\right)\right) \\
& =T\left(\mathbf{0}, \lambda\left(S(\mathbf{U})-S\left(\mathbf{U}_{h}\right)\right)\right)+\left(T-T_{h}\right)\left(\mathbf{f}, \lambda S\left(\mathbf{U}_{h}\right)\right) \tag{23}
\end{align*}
$$

We now bound the first term in the right hand side of (23). If $\mathbf{V}=(\mathbf{v}, \boldsymbol{\tau})$ is defined by

$$
\mathbf{V}=T\left(\mathbf{0}, \lambda\left(S(\mathbf{U})-S\left(\mathbf{U}_{h}\right)\right)\right)
$$

then there exists $q \in L_{0}^{2}(\Omega)$ such that $(\mathbf{v}, q) \in H_{0}^{1}(\Omega)^{2} \times L_{0}^{2}(\Omega)$ satisfies

$$
\begin{aligned}
-2\left(\eta_{s}+\eta_{p}\right) \operatorname{div} \boldsymbol{\epsilon}(\mathbf{v})+\nabla q & =\lambda \operatorname{div}\left(S(\mathbf{U})-S\left(\mathbf{U}_{h}\right)\right) \\
\operatorname{div} \mathbf{v} & =0
\end{aligned}
$$

We then have the following estimates

$$
\begin{aligned}
\|\mathbf{V}\|_{H^{1} \times L^{2}}=\|\mathbf{v}, q\|_{H^{1} \times L^{2}} & \leq C \lambda\left\|\operatorname{div}\left(S(\mathbf{U})-S\left(\mathbf{U}_{h}\right)\right)\right\|_{H^{-1}} \\
& \leq \tilde{C} \lambda\left\|S(\mathbf{U})-S\left(\mathbf{U}_{h}\right)\right\|_{L^{2}}
\end{aligned}
$$

Now we have

$$
\begin{aligned}
& 2 \eta_{p}\left(S(\mathbf{U})-S\left(\mathbf{U}_{h}\right)\right) \\
& \quad=\nabla \mathbf{u} \boldsymbol{\sigma}+\boldsymbol{\sigma} \nabla \mathbf{u}^{T}-\left(\nabla \mathbf{u}_{h}\right) \boldsymbol{\sigma}_{h}-\boldsymbol{\sigma}_{h}\left(\nabla \mathbf{u}_{h}\right)^{T} \\
& \quad=\nabla\left(\mathbf{u}-\mathbf{u}_{h}\right) \boldsymbol{\sigma}+\left(\nabla \mathbf{u}_{h}\right)\left(\boldsymbol{\sigma}-\boldsymbol{\sigma}_{h}\right)+\boldsymbol{\sigma} \nabla\left(\mathbf{u}-\mathbf{u}_{h}\right)^{T}+\left(\boldsymbol{\sigma}-\boldsymbol{\sigma}_{h}\right)\left(\nabla \mathbf{u}_{h}\right)^{T}
\end{aligned}
$$

so that

$$
\left\|S(\mathbf{U})-S\left(\mathbf{U}_{h}\right)\right\|_{L^{2}} \leq \frac{1}{\eta_{p}}\left(\left\|\nabla\left(\mathbf{u}-\mathbf{u}_{h}\right)\right\|_{L^{2}}\|\boldsymbol{\sigma}\|_{L^{\infty}}+\left\|\boldsymbol{\sigma}-\boldsymbol{\sigma}_{h}\right\|_{L^{2}}\left\|\nabla \mathbf{u}_{h}\right\|_{L^{\infty}}\right)
$$

We now bound the last term of the above inequality. We have

$$
\begin{aligned}
\left\|\nabla \mathbf{u}_{h}\right\|_{L^{\infty}} & \leq\left\|\nabla\left(\mathbf{u}-\mathbf{u}_{h}\right)\right\|_{L^{\infty}}+\|\nabla \mathbf{u}\|_{L^{\infty}} \\
& \leq\left\|\nabla\left(\mathbf{u}-r_{h} \mathbf{u}\right)\right\|_{L^{\infty}}+\left\|\nabla\left(r_{h} \mathbf{u}-\mathbf{u}_{h}\right)\right\|_{L^{\infty}}+\|\nabla \mathbf{u}\|_{L^{\infty}} .
\end{aligned}
$$

We then use standard interpolation estimates [10], an inverse estimate, and a Sobolev imbedding theorem to obtain

$$
\left\|\nabla \mathbf{u}_{h}\right\|_{L^{\infty}} \leq C\left(\|\mathbf{u}\|_{W^{2, r}}+\frac{1}{h}\left\|\nabla\left(r_{h} \mathbf{u}-\mathbf{u}_{h}\right)\right\|_{L^{2}}+\|\mathbf{u}\|_{W^{2, r}}\right)
$$

with $C$ independent of $\lambda$ and $h$. Finally, applying Theorem 4.4 to the above estimate we have, for $h \leq \bar{h}$ and $\lambda \leq \bar{\lambda},\left\|\nabla \mathbf{u}_{h}\right\|_{L^{\infty}} \leq C$ so that

$$
\left\|S(\mathbf{U})-S\left(\mathbf{U}_{h}\right)\right\|_{L^{2}} \leq C\left\|\mathbf{U}-\mathbf{U}_{h}\right\|_{H^{1} \times L^{2}}
$$

$C$ being independent of $h$ and $\lambda \in[0, \bar{\lambda}]$. Thus, we have shown that

$$
\begin{equation*}
\left\|T\left(\mathbf{0}, \lambda\left(S(\mathbf{U})-S\left(\mathbf{U}_{h}\right)\right)\right)\right\|_{H^{1} \times L^{2}} \leq C \lambda\left\|\mathbf{U}-\mathbf{U}_{h}\right\|_{H^{1} \times L^{2}} \tag{24}
\end{equation*}
$$

where $C$ does not depend on $h$ and $\lambda \in[0, \bar{\lambda}]$.

In order to bound the second term in the right hand side of (23), we use Lemma 5.1. There is a constant $C$ independent of $h$ and $\lambda$ such that we have, for $h$ sufficiently small

$$
\begin{equation*}
\left\|\left(T-T_{h}\right)\left(\mathbf{f}, \lambda S\left(\mathbf{U}_{h}\right)\right)\right\|_{H^{1} \times L^{2}} \leq C\left(\sum_{K \in \mathcal{T}_{h}} \mu_{K}^{2}\left(\mathbf{f}, \lambda S\left(\mathbf{U}_{h}\right)\right)\right)^{1 / 2} \tag{25}
\end{equation*}
$$

Finally, estimates (24) and (25) in (23) yield the result provided $\lambda$ is small enough.
Remark 5.3. Proceeding as in $[2,28]$, we can prove a lower bound similar to the upper bound of Lemma 5.1. Then, we can also prove a lower bound similar to the upper bound of Theorem 5.2. Thus, the error estimator is equivalent to the true error. We did not include such a result in this paper in order to shorten the presentation.

Remark 5.4. We can show a sharper estimate than (22). Indeed, let us introduce as in [6] the norm $\|\cdot\|_{W}$ defined, for all $(\mathbf{v}, q, \boldsymbol{\tau}) \in W$, by

$$
\begin{equation*}
\|\mathbf{v}, q, \boldsymbol{\tau}\|_{W}^{2}=2\left(\eta_{s}+\eta_{p}\right)\|\boldsymbol{\epsilon}(\mathbf{v})\|^{2}+\frac{1}{2\left(\eta_{s}+\eta_{p}\right)}\|q\|^{2}+\frac{1}{2 \eta_{p}}\|\boldsymbol{\tau}\|^{2} \tag{26}
\end{equation*}
$$

Then, the following estimate also holds

$$
\left\|\mathbf{u}-\mathbf{u}_{h}, p-p_{h}, \boldsymbol{\sigma}-\boldsymbol{\sigma}_{h}\right\|_{W} \leq C\left(\sum_{K \in \mathcal{T}_{h}} \mu_{K}^{2}\left(\mathbf{f}, \lambda S\left(\mathbf{u}_{h}, \boldsymbol{\sigma}_{h}\right)\right)\right)^{1 / 2}
$$

where $C$ is independent of $\lambda, h$ and $\eta_{s}, \eta_{p}$. Here $(\mathbf{u}, p, \boldsymbol{\sigma})$ is the solution of $(1),\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h}\right)$ is the solution of (9). Thus, our a posteriori error estimates also hold when the solvent viscosity is small.

## 6. Link with an EVSS formulation

The Elastic Viscous Split Stress (EVSS) formulation corresponding to (9) is obtained from the following differential problem: find the velocity $\mathbf{u}$, pressure $p$, extra-stresses $\boldsymbol{\sigma}$ and $\mathbf{D}$ such that

$$
\begin{align*}
-2\left(\eta_{s}+\eta_{p}\right) \operatorname{div} \boldsymbol{\epsilon}(\mathbf{u})+\nabla p-\operatorname{div}\left(\boldsymbol{\sigma}-2 \eta_{p} \mathbf{D}\right) & =\mathbf{f} \\
\operatorname{div} \mathbf{u} & =0 \\
\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}-\frac{\lambda}{2 \eta_{p}}\left((\nabla \mathbf{u}) \boldsymbol{\sigma}+\boldsymbol{\sigma}\left(\nabla \mathbf{u}^{T}\right)\right)-\boldsymbol{\epsilon}(\mathbf{u}) & =0  \tag{27}\\
\mathbf{D}-\boldsymbol{\epsilon}(\mathbf{u}) & =0
\end{align*}
$$

in $\Omega$, where the velocity $\mathbf{u}$ vanishes on $\partial \Omega$. Obviously, at the continuous level (27) is equivalent to (1).
In order to solve this problem, we consider the following EVSS scheme: find $\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h}, \mathbf{D}_{h}\right) \in V_{h} \times Q_{h} \times$ $M_{h} \times M_{h}$ such that

$$
\begin{align*}
& 2\left(\eta_{s}+\eta_{p}\right)\left(\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), \boldsymbol{\epsilon}(\mathbf{v})\right)-\left(p_{h}, \operatorname{div} \mathbf{v}\right)+\left(\boldsymbol{\sigma}_{h}-2 \eta_{p} \mathbf{D}_{h}, \boldsymbol{\epsilon}(\mathbf{v})\right)-(\mathbf{f}, \mathbf{v}) \\
& -\left(\operatorname{div} \mathbf{u}_{h}, q\right)-\sum_{K \in \mathcal{T}_{h}} \frac{\alpha h_{K}^{2}}{2 \eta_{p}}\left(-2 \eta_{s} \operatorname{div} \boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right)+\nabla p_{h}-\operatorname{div} \boldsymbol{\sigma}_{h}-\mathbf{f}, \nabla q\right)_{K}  \tag{28}\\
& -\left(\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}-\frac{\lambda}{2 \eta_{p}}\left(\nabla \mathbf{u}_{h} \boldsymbol{\sigma}_{h}+\boldsymbol{\sigma}_{h} \nabla \mathbf{u}_{h}^{T}\right)-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), \boldsymbol{\tau}\right) \\
& +\left(\mathbf{D}_{h}-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), \mathbf{E}\right)=0
\end{align*}
$$

for all $(\mathbf{v}, q, \boldsymbol{\tau}, \mathbf{E}) \in V_{h} \times Q_{h} \times M_{h} \times M_{h}$. In the above scheme, only the stabilization terms corresponding to the first equation of (27) have been added in order to avoid spurious oscillations for the pressure. However, no special attention is paid to the extra-stress $\boldsymbol{\sigma}_{h}$ and the stabilization term

$$
\left(\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}-\frac{\lambda}{2 \eta_{p}}\left(\nabla \mathbf{u}_{h} \boldsymbol{\sigma}_{h}+\boldsymbol{\sigma}_{h} \nabla \mathbf{u}_{h}^{T}\right)-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), 2 \eta_{p} \beta \boldsymbol{\epsilon}(\mathbf{v})\right)
$$

present in (9) is missing in (28). The extra-stress $\mathbf{D}_{h}$ is added only for stability purposes and thus (28) can be interpreted as an EVSS scheme for solving (1).

In the linear case (i.e. when $\lambda=0$ ), the GLS scheme (9) with $\beta=1$ is equivalent to the EVSS formulation (28), as explained in [6]. Indeed, $\mathbf{D}_{h} \in M_{h}$ is such that

$$
\left(\mathbf{D}_{h}-\boldsymbol{\epsilon}\left(u_{h}\right), \mathbf{E}\right)=0 \quad \forall \mathbf{E} \in M_{h},
$$

and, when $\lambda=0, \sigma_{h} \in M_{h}$ is such that

$$
\left(\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}\right), \boldsymbol{\tau}\right)=0 \quad \forall \boldsymbol{\tau} \in M_{h}
$$

Thus, $\mathbf{D}_{h}=\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}$ and (28) reduces to (9) with $\beta=1$.
In the nonlinear case, the two GLS and EVSS schemes are not equivalent anymore. Indeed, when $\lambda \neq 0$, we have

$$
\left(\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}+\frac{\lambda}{2 \eta_{p}}\left(\nabla \mathbf{u}_{h} \boldsymbol{\sigma}_{h}+\boldsymbol{\sigma}_{h} \nabla \mathbf{u}_{h}^{T}\right)-\mathbf{D}_{h}, \boldsymbol{\tau}\right)=0 \quad \forall \boldsymbol{\tau} \in M_{h},
$$

and since $\left(\nabla \mathbf{u}_{h} \boldsymbol{\sigma}_{h}+\boldsymbol{\sigma}_{h} \nabla \mathbf{u}_{h}^{T}\right) \notin M_{h}$ we cannot conclude.
Proceeding as we did for the GLS scheme (9), and under the same assumptions, we can prove that (28) has a solution converging to the solution of (27).

## 7. Numerical Results

We now discuss iterative decoupling schemes for solving (9) with $\beta=1$ and (28).
Our iterative procedures are an extension of those presented in [6] and [12] and allow velocity and pressure computations to be decoupled from extra-stresses computations.

Let us start with the GLS scheme (9) and $\beta=1$. Let $\left(\mathbf{u}_{h}^{n}, p_{h}^{n}, \boldsymbol{\sigma}_{h}^{n}\right)$ be the known approximation of $\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h}\right)$ after $n$ steps. Step $(n+1)$ of the algorithm consists in first computing $\left(\mathbf{u}_{h}^{n+1}, p_{h}^{n+1}\right)$ by using the mass and momentum equations, then under-relaxing with parameter $0 \leq \omega \leq 1$, and finally computing $\sigma_{h}^{n+1}$ with the constitutive relationship. Thus iteration $(n+1)$ consists in finding $\left(\tilde{\mathbf{u}}_{h}^{n+1}, \tilde{p}_{h}^{n+1}\right) \in V_{h} \times Q_{h}$ such that

$$
\begin{align*}
& 2 \eta_{s}\left(\boldsymbol{\epsilon}\left(\tilde{\mathbf{u}}_{h}^{n+1}\right), \boldsymbol{\epsilon}(\mathbf{v})\right)-\left(\tilde{p}_{h}^{n+1}, \operatorname{div} \mathbf{v}\right)+\left(\boldsymbol{\sigma}_{h}^{n}, \boldsymbol{\epsilon}(\mathbf{v})\right) \\
& -\left(\boldsymbol{\sigma}_{h}^{n}-\lambda\left(\nabla \mathbf{u}_{h}^{n} \boldsymbol{\sigma}_{h}^{n}+\boldsymbol{\sigma}_{h}^{n}\left(\nabla \mathbf{u}_{h}^{n}\right)^{T}\right)-2 \eta_{p} \boldsymbol{\epsilon}\left(\tilde{\mathbf{u}}_{h}^{n+1}\right), \boldsymbol{\epsilon}(\mathbf{v})\right)-(\mathbf{f}, \mathbf{v}) \\
& -\left(\operatorname{div} \tilde{\mathbf{u}}_{h}^{n+1}, q\right)  \tag{29}\\
& -\sum_{K \in \mathcal{T}_{h}} \frac{\alpha h_{K}^{2}}{2 \eta_{p}}\left(-2 \eta_{s} \operatorname{div} \boldsymbol{\epsilon}\left(\tilde{\mathbf{u}}_{h}^{n+1}\right)+\nabla \tilde{p}_{h}^{n+1}-\operatorname{div} \boldsymbol{\sigma}_{h}^{n}-\mathbf{f}, \nabla q\right)_{K}=0,
\end{align*}
$$

for all $(\mathbf{v}, q) \in V_{h} \times Q_{h}$, then updating $\mathbf{u}_{h}^{n+1}$ and $p_{h}^{n+1}$ as following

$$
\begin{aligned}
& \mathbf{u}_{h}^{n+1}=\omega \tilde{\mathbf{u}}_{h}^{n+1}+(1-\omega) \mathbf{u}_{h}^{n} \\
& p_{h}^{n+1}=\omega \tilde{p}_{h}^{n+1}+(1-\omega) p_{h}^{n}
\end{aligned}
$$

and finally finding $\boldsymbol{\sigma}_{h}^{n+1} \in M_{h}$ such that

$$
\begin{equation*}
\left(\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}^{n+1}-\frac{\lambda}{2 \eta_{p}}\left(\nabla \mathbf{u}_{h}^{n} \boldsymbol{\sigma}_{h}^{n}+\boldsymbol{\sigma}_{h}^{n}\left(\nabla \mathbf{u}_{h}^{n}\right)^{T}\right)-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}^{n+1}\right), \boldsymbol{\tau}\right)=0 \quad \forall \boldsymbol{\tau} \in M_{h} \tag{30}
\end{equation*}
$$

Let us now turn to the EVSS scheme (28). Let $\left(\mathbf{u}_{h}^{n}, p_{h}^{n}, \boldsymbol{\sigma}_{h}^{n}, \mathbf{D}_{h}^{n}\right)$ be the known approximation of $\left(\mathbf{u}_{h}, p_{h}, \boldsymbol{\sigma}_{h}, \mathbf{D}_{h}\right)$ after $n$ steps. Iteration $(n+1)$ consists in finding $\left(\tilde{\mathbf{u}}_{h}^{n+1}, \tilde{p}_{h}^{n+1}\right) \in V_{h} \times Q_{h}$ such that

$$
\begin{align*}
& 2\left(\eta_{s}+\eta_{p}\right)\left(\boldsymbol{\epsilon}\left(\tilde{\mathbf{u}}_{h}^{n+1}\right), \boldsymbol{\epsilon}(\mathbf{v})\right)-\left(\tilde{p}_{h}^{n+1}, \operatorname{div} \mathbf{v}\right)+\left(\boldsymbol{\sigma}_{h}^{n}-2 \eta_{p} \mathbf{D}_{h}^{n}, \boldsymbol{\epsilon}(\mathbf{v})\right)-(\mathbf{f}, \mathbf{v}) \\
& -\left(\operatorname{div} \tilde{\mathbf{u}}_{h}^{n+1}, q\right) \\
& -\sum_{K \in \mathcal{T}_{h}} \frac{\alpha h_{K}^{2}}{2 \eta_{p}}\left(-2 \eta_{s} \operatorname{div} \boldsymbol{\epsilon}\left(\tilde{\mathbf{u}}_{h}^{n+1}\right)+\nabla \tilde{p}_{h}^{n+1}-\operatorname{div} \boldsymbol{\sigma}_{h}^{n}-\mathbf{f}, \nabla q\right)_{K}=0 \tag{31}
\end{align*}
$$

for all $(\mathbf{v}, q) \in V_{h} \times Q_{h}$, then updating $\mathbf{u}_{h}^{n+1}$ and $p_{h}^{n+1}$ as following

$$
\begin{aligned}
\mathbf{u}_{h}^{n+1} & =\omega \tilde{\mathbf{u}}_{h}^{n+1}+(1-\omega) \mathbf{u}_{h}^{n} \\
p_{h}^{n+1} & =\omega \tilde{p}_{h}^{n+1}+(1-\omega) p_{h}^{n}
\end{aligned}
$$

and finally finding $\boldsymbol{\sigma}_{h}^{n+1} \in M_{h}$ such that

$$
\begin{equation*}
\left(\frac{1}{2 \eta_{p}} \boldsymbol{\sigma}_{h}^{n+1}-\frac{\lambda}{2 \eta_{p}}\left(\nabla \mathbf{u}_{h}^{n} \boldsymbol{\sigma}_{h}^{n}+\boldsymbol{\sigma}_{h}^{n}\left(\nabla \mathbf{u}_{h}^{n}\right)^{T}\right)-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}^{n+1}\right), \boldsymbol{\tau}\right)=0 \quad \forall \boldsymbol{\tau} \in M_{h} \tag{32}
\end{equation*}
$$

and $\mathbf{D}_{h}^{n+1} \in M_{h}$ such that

$$
\begin{equation*}
\left(\mathbf{D}_{h}^{n+1}-\boldsymbol{\epsilon}\left(\mathbf{u}_{h}^{n+1}\right), \mathbf{E}\right)=0, \quad \forall \mathbf{E} \in M_{h} \tag{33}
\end{equation*}
$$

The computational effort required to compute the velocity and pressure thus corresponds to solving a Stokes' problem with a conventional GLS method, whereas the computations of the extra-stresses are explicit provided the mass matrices are lumped. The decoupled EVSS procedure is more interesting than the GLS one from the implementation point of view. Indeed, if the constitutive relationship between the extra-stress and the velocity field is replaced by a more realistic equation (for instance the Phan-Thien-Tanner or the FENE-P models), then only the portion of the code corresponding to (32) should be updated. This is not the case of the GLS scheme since some of the terms present in (29) are due to the constitutive equation.

The EVSS scheme is particularly interesting for molecular models, the constitutive relationship being replaced by a stochastic differential equation for the dumbbells elongations, see for instance $[8,16,20]$.

### 7.1. A simple test case

In order to investigate numerically the rate of convergence of both GLS and EVSS schemes we have set

$$
\mathbf{u}\left(x_{1}, x_{2}\right)=\binom{u_{1}\left(x_{2}\right)}{u_{2}\left(x_{1}\right)}=\binom{\sin \left(\pi x_{2}\right) e^{x_{2}}}{\sin \left(\pi x_{1}\right) e^{x_{1}}}, \quad p\left(x_{1}, x_{2}\right)=0
$$

The constitutive relationship (the last equation of (1)) then leads to

$$
\sigma_{11}=2 \eta_{p} \lambda u_{1}^{\prime} \gamma, \quad \sigma_{12}=\eta_{p} \gamma, \quad \sigma_{22}=2 \eta_{p} \lambda u_{2}^{\prime} \gamma
$$

where $\gamma$ is defined by

$$
\gamma\left(x_{1}, x_{2}\right)=\frac{u_{1}^{\prime}\left(x_{2}\right)+u_{2}^{\prime}\left(x_{1}\right)}{1-4 \lambda^{2} u_{1}^{\prime}\left(x_{2}\right) u_{2}^{\prime}\left(x_{1}\right)}
$$

whenever $\lambda<1 / \sqrt{4 \max \left|u_{1}^{\prime} u_{2}^{\prime}\right|}=\frac{1}{2 \pi e} \simeq 0.058$. Then, the source term in the momentum equation is given by

$$
\mathbf{f}=-\eta_{s}\binom{u_{1}^{\prime \prime}}{u_{2}^{\prime \prime}}-\eta_{p}\binom{2 \lambda u_{1}^{\prime} \partial \gamma / \partial x_{1}+\partial \gamma / \partial x_{2}}{\partial \gamma / \partial x_{1}+2 \lambda u_{2}^{\prime} \partial \gamma / \partial x_{2}}
$$

The calculation domain was the unit square cut into several squares, each square being cut into two triangles along one of its diagonal. The viscosities were $\eta_{s}=0.01, \eta_{p}=1$, the GLS stabilization parameter was $\alpha=0.01$, the relaxation parameter was $\omega=0.5$, the elastic time scale was $\lambda=0.02$. In Table 1 we have reported the $L^{2}$ error of the scalar unknowns (velocity components $u_{1}, u_{2}$, pressure $p$, extra-stress components $\sigma_{11}, \ldots$ ) with several meshes. Clearly, the order of convergence with respect to $h$ is close to two for the velocity and one for the pressure and extra-stress.

TABLE 1. $L^{2}$ error and number of iterations to achieve convergence with several meshes, top: GLS scheme, bottom EVSS scheme.

| Mesh | $u_{1}$ | $u_{2}$ | $p$ | $\sigma_{11}$ | $\sigma_{12}$ | $\sigma_{22}$ | iterations |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $10 \times 10$ | 0.00041 | 0.00041 | 0.36 | 0.19 | 0.41 | 0.19 | 22 |
| $20 \times 20$ | 0.00012 | 0.00012 | 0.17 | 0.066 | 0.14 | 0.066 | 22 |
| $40 \times 40$ | 0.000037 | 0.000037 | 0.082 | 0.023 | 0.048 | 0.023 | 23 |
| $80 \times 80$ | 0.000010 | 0.000010 | 0.040 | 0.0079 | 0.017 | 0.0079 | 23 |
| Mesh | $u_{1}$ | $u_{2}$ | $p$ | $\sigma_{11}$ | $\sigma_{12}$ | $\sigma_{22}$ | iterations |
| $10 \times 10$ | 0.00049 | 0.00049 | 0.38 | 0.19 | 0.41 | 0.19 | 22 |
| $20 \times 20$ | 0.00016 | 0.00016 | 0.18 | 0.066 | 0.14 | 0.066 | 22 |
| $40 \times 40$ | 0.000052 | 0.000052 | 0.091 | 0.022 | 0.047 | 0.022 | 23 |
| $80 \times 80$ | 0.000015 | 0.000015 | 0.045 | 0.0078 | 0.016 | 0.0078 | 23 |

The stopping criterion for the iterative procedures (29-33) to reach convergence with respect to $n$ was $10^{-6}$ on the relative discrepancy. From Table 2, we can see that, for a fixed $\lambda$, the number of iterations does not depend on the mesh size nor on the solvent viscosity. However, when $\lambda$ increases, the number of iterations increases.

TABLE 2. Number of iterations to achieve convergence with respect to $n$ on a $20 \times 20$ mesh. Left table: with several solvent viscosities $(\lambda=0.02)$. Right table: with several relaxation times $\left(\eta_{s}=0.01\right)$. The xx symbol means that the scheme was divergent.

| $\eta_{s}$ | GLS | EVSS |
| :---: | :---: | :---: |
| 1 | 22 | 22 |
| 0.01 | 22 | 22 |
| 0 | 22 | 22 |


| $\lambda$ | GLS | EVSS |
| :--- | :---: | :---: |
| 0.02 | 22 | 22 |
| 0.03 | 22 | 22 |
| 0.04 | 22 | 22 |
| 0.05 | 36 | 35 |
| 0.055 | 88 | 62 |
| 0.06 | xx | 225 |
| 0.065 | xx | xx |



Figure 1. The coarsest mesh (mesh 1) used for the computations (201 vertices). Mesh 2 is obtained by cutting the mesh size by two ( 721 vertices), mesh 3 by four ( 2721 vertices) and mesh 4 by eight ( 10561 vertices).

### 7.2. The 4:1 planar contraction

The 4:1 planar contraction is a classical test case in the frame of non Newtonian flows [22]. This test case is not covered by our theory since the calculation domain is not convex, and the velocity gradient is not bounded at the reentrant corner. From the physical point of view, instabilities are observed at high Deborah numbers. From the numerical point of view, most viscoelastic codes fail to converge at high Deborah numbers. Moreover, as reported in several papers (see for instance Sect. 7.2 of [1] for details), the maximum attainable Deborah number seems to decrease with mesh size. The well posedness of the problem is still an open question.

We have performed computations on half of the contraction, with four different meshes, the coarsest mesh being the one of Figure 1. The inlet and outlet velocities were imposed to be parabolic, with maximum velocity one at the inlet. The viscosities were $\eta_{s}=0.01, \eta_{p}=1$, the GLS stabilization parameter was $\alpha=0.01$, the elastic time scale $\lambda$ was ranging from 0.01 to 0.04 . In Table 3 we have reported the number of iterations to reach convergence for several mesh sizes and several values of $\lambda$ when using the EVSS scheme (similar results were obtained with the GLS scheme). It seems that, the more the mesh is refined, the smaller the maximum attainable Deborah number. Moreover, increasing the solvent viscosity does not help convergence.

There are (at least) two possible reasons to explain failure of the numerical procedure (31-33).

- The decoupled procedure (31-33) used to obtain the solution of (28) is not appropriate. Newton's method should be implemented to check if convergence could be obtained with higher Deborah numbers, and several meshes.
- Problem (1) has many solutions or no solution at all for the contraction flow.

Table 3. The $4: 1$ planar contraction flow. Number of iterations to achieve convergence with four meshes (see caption of Fig. 1) and several values of $\lambda$. Left: $\lambda=0.04$, middle: $\lambda=0.02$, right: $\lambda=0.01$. The xx symbol means that the scheme was divergent.

| Mesh | iterations |
| :---: | :---: |
| 1 | 17 |
| 2 | xx |
| 3 | xx |
| 4 | xx |


| Mesh | iterations |
| :---: | :---: |
| 1 | 16 |
| 2 | 16 |
| 3 | 17 |
| 4 | xx |


| Mesh | iterations |
| :---: | :---: |
| 1 | 15 |
| 2 | 15 |
| 3 | 16 |
| 4 | 18 |

## 8. Conclusion and perspectives

In this paper, existence, a priori and a posteriori error estimates have been obtained for the GLS approximation of an Oldroyd-B problem without convection. An EVSS method is also introduced and the link with the GLS method is shown. Numerical results are presented.

An important point is now to extend these theoretical predictions to the case when the convective terms are added to the momentum equation and to the constitutive equation. Also, we are looking forward to extending this work to the case of mesoscopic models in order to justify the computations performed in $[7,8]$.

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    1 Département de Mathématiques, École Polytechnique Fédérale de Lausanne, 1015 Lausanne, Switzerland.

