HAL
open science

# A distributed resistance inverse method for flow obstacle identification from internal velocity measurements 

Jorge Aguayo, Cristobal Bertoglio, Axel Osses

## To cite this version:

Jorge Aguayo, Cristobal Bertoglio, Axel Osses. A distributed resistance inverse method for flow obstacle identification from internal velocity measurements. 2020. hal-02913958

HAL Id: hal-02913958

## https://hal.science/hal-02913958

Preprint submitted on 11 Aug 2020

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers.

L'archive ouverte pluridisciplinaire HAL, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

# A distributed resistance inverse method for flow obstacle identification from internal velocity measurements 

Jorge Aguayo ${ }^{\text {a }}$, Cristóbal Bertoglio ${ }^{\text {b }}$, Axel Osses ${ }^{\text {a,c }}$<br>${ }^{a}$ Mathematical Engineering Department, Faculty of Physical and Mathematical Sciences, Universidad de Chile, Santiago, Chile<br>${ }^{b}$ Bernoulli Institute, University of Groningen, Groningen, The Netherlands<br>${ }^{c}$ Center for Mathematical Modeling, Universidad de Chile, Santiago, Chile


#### Abstract

We present a penalization parameter method for obstacle identification in an incompressible fluid flow for a modified version of the Oseen equations. The proposed method consist in adding a high resistance potential to the system such that some subset of its boundary support represents the obstacle. This allows to work in a fixed domain and highly simplify the solution of the inverse problem via some suitable cost functional. Existence of minimizers and first and second order optimality conditions are derived through the differentiability of the solutions of the Oseen equation with respect to the potential. Finally, several numerical experiments using Navier-Stokes flow illustrate the applicability of the method, for the localization of a bi-dimensional cardiac valve from MRI and ultrasound flow type imaging data.


## 1. Introduction

The pathway of blood flow through the heart is regulated by four membrane structures or valves, opening and closing by differential blood pressures. Valvular heart diseases affect $20 \%$ of the population. Among them, aortic valve stenosis is the most prevalent one in developed countries [21].

Imaging the 3D shape of the valves is a challenging task. For instance, the aortic valve is very thin $(0.5 \mathrm{~mm})$, and therefore its shape can be imaged nowadays using only two modalities: computerized tomography (CT) and transesophageal echocardiography (TEE). Since CT is based on X-rays, it is only used in patients that are selected for valvular replacement in order to obtain the aortic root dimensions for sizing the prosthesis. Such CT images are usually obtained when the valve is closed. Obtaining the image at open valve position requires about 5 times larger radiation dose since the complete cardiac cycle has to be imaged. This is equivalent to the annual recommended total radiation dose (see [25]). TEE is a newer technique, but highly invasive: a wire is inserted through the esophagus of the patient involving a cumbersome procedure. TEE is most of the time applied to monitor valve surgeries (see [18]), and is therefore rarely applied in a diagnostic phase.

Magnetic resonance imaging (MRI) allows to image anatomical structures in a non-invasive and nonionising way. Unfortunately, the aortic valve geometry is not directly visible with MRI, since the valve thickness is smaller than the image voxel size.

However, visual inspection of 3D Flow MRI Imaging (4D Flow, see [17) data sets clearly shows that the valvular shape could be retrieved from the flow pattern in the valve surroundings. This fact motivates to formulate, analyze and numerically assess the inverse problem of inferring rigid obstacles from interior velocity measurements, with the ultimate goal of recovering the shape of cardiac valves at the opening position from velocity images.

[^0]Available approaches studied and reported in the mathematical literature are limited to the detection of obstacles in viscous fluid flows using boundary stress measurements, which limits the inverse problem to oversimplified shapes, usually of circular nature [3, 7, 14, 15]. Noteworthy, boundary stress measurements would mean in practice to introduce a catheter close to the valve.

The novelty of this work is the identification of flow obstacles by a distributed resistance inverse method. That is, we propose the incorporation of a large resistance term that allows to model the effect of an obstacle in the viscous fluid. This method reduces the obstacle identification to a simpler potential inverse problem. This idea is inspired by [13, 5, [19, who model cardiac valves using a resistive immerse surface (RIS) given by a Dirac distribution. The problem of using RIS for valve identification would be the need of modifying the discretization of the domain at every iteration of the identification procedure. In contrast, the distributed resistance term that we propose here allows us to work in a fixed domain to solve the valve shape identification problem. This distributed resistance method can also be useful to estimate porosity in porous media flows following Brinkman's law [12, 6.

The paper is structured as follows. In Section 2, a parameter identification problem is defined for the Oseen equations (as a linearization of Navier-Stokes equations) adding a resistive potential term with the form $\gamma \boldsymbol{u}$, where $\boldsymbol{u}$ is the velocity field and $\gamma \geq 0$ is a function that takes large values in zones where is the obstacle should be located, or 0. In Section 3, a proof of the existence of minimizers using appropriate techniques is presented. In Sections 4 and 5, first and second order optimality conditions for some suitable cost functions are established, motivating the choice of suitable spaces for the parameter and the feasibility of implementing optimization algorithms to numerically solve this problem. Section 6 presents numerical examples that illustrate the feasibility of the proposed penalty method using Navier-Stokes equations from a 2 D reference case representing the blood flow passing through the aortic valve. We reconstruct the position of the valve from global or local velocity field measurements. For this purpose, synthetic images based on MRI and ultrasound type measurements are used. Finally, conclusions and future work are presented in Section 7.

## 2. Model problem

Consider a non-empty bounded domain $\Omega \subseteq \mathbb{R}^{N}, N \in\{2,3\}$. The Lebesgue measure of $\Omega$ is denoted by $|\Omega|$, which extends to lesser dimension spaces. The norm and seminorms for Sobolev spaces $W^{m, p}(\Omega)$ is denoted by $\|\cdot\|_{m, p, \Omega}$ and $|\cdot|_{m, p, \Omega}$, respectively. For $p=2$, the norm, seminorms and inner product of the space $W^{m, 2}(\Omega)=H^{m}(\Omega)$ are denoted by $\|\cdot\|_{m, \Omega},|\cdot|_{m, \Omega}$ and $(\cdot, \cdot)_{m, \Omega}$, respectively. Also, $\|\cdot\|_{\infty, \Omega}$ denotes the norm of $L^{\infty}(\Omega)$. The spaces $\boldsymbol{H}^{m}(\Omega)$ and $\boldsymbol{W}^{m, p}(\Omega)$ are defined by $\boldsymbol{H}^{m}(\Omega)=\left[H^{m}(\Omega)\right]^{N}$ and $\boldsymbol{W}^{m, p}(\Omega)=\left[W^{m, p}(\Omega)\right]^{N}$. The notation for norms, seminorms and inner products will be extended from $W^{m, p}(\Omega)$ or $H^{m}(\Omega)$.

Consider $\alpha>0, \nu>0, s \geq 0, M>0, \Omega$ with a Lipschitz boundary $\partial \Omega=\Gamma_{D} \cup \Gamma_{N}$ such that $\operatorname{int}\left(\Gamma_{D}\right) \cap \operatorname{int}\left(\Gamma_{N}\right)=\emptyset$, an open subset $\omega \subseteq \Omega, \boldsymbol{u}_{R} \in \boldsymbol{H}^{1}(\omega)^{d}, \boldsymbol{u}_{D} \in \boldsymbol{H}^{1 / 2}\left(\Gamma_{D}\right)$ and $\boldsymbol{a} \in \boldsymbol{W}^{1, \infty}(\Omega)$ such that $\operatorname{div} \boldsymbol{a}=0$ and $\boldsymbol{a} \cdot \boldsymbol{n} \geq 0$ on $\Gamma_{N}$, where $\boldsymbol{n}$ is the outer normal vector. The model problem is defined by

$$
\begin{array}{rlrl}
\operatorname{minimize} & J(\gamma, \boldsymbol{u})=\frac{1}{2}\left\|\boldsymbol{u}-\boldsymbol{u}_{R}\right\|_{0, \omega}^{2}+\frac{\alpha}{2}\|\gamma\|_{s, \Omega}^{2} & & \\
\text { subject to } & -\nu \triangle \boldsymbol{u}+(\nabla \boldsymbol{u}) \boldsymbol{a}+\nabla p+\gamma \boldsymbol{u} & =\mathbf{0} &  \tag{2}\\
\text { in } \Omega \\
\operatorname{div} \boldsymbol{u} & =0 & & \text { in } \Omega \\
\boldsymbol{u} & =\boldsymbol{u}_{D} & & \text { on } \Gamma_{D} \\
-\nu \frac{\partial \boldsymbol{u}}{\partial \boldsymbol{n}}+p \boldsymbol{n} & =\mathbf{0} & & \text { on } \Gamma_{N} \\
& \boldsymbol{u} \in \boldsymbol{H}^{1}(\Omega), \gamma \in H^{s}(\Omega) & & \\
0 \leq \gamma \leq M \text { a.e. in } \Omega . & &
\end{array}
$$

The modified version of Oseen equations given by Equation (2) models the velocity $\boldsymbol{u}$ and pressure $p$ of a fluid that passes through the control volume $\Omega$. The term $\gamma \boldsymbol{u}$, with $\gamma \geq 0$, represents the distributed resistance that $\Omega$ imposes to the fluid. A zero value of $\gamma$ represents that the fluid follows the Oseen equations. As the $\gamma$ values increases in a certain area, the magnitude of $\boldsymbol{u}$ tends to decrease to 0 . According to Brinkman's law [12, 6, $1 /$ gamma is an indicator of permeability. The media is completly permeable when $\gamma=0$, while $\gamma$ tends to $+\infty$ in zones where there are obstacles. Bounded values of $\gamma$ model porous media, where porosity decreases as gamma increases.

In order to simplifying the analysis of this problem, an equivalent formulation with homogeneous Dirichlet boundary condition is proposed. Let $\boldsymbol{G} \in \boldsymbol{H}^{1}(\Omega)$ such that $\boldsymbol{G}=\boldsymbol{u}_{D}$ on $\Gamma_{D}$. Taking $\boldsymbol{v}=\boldsymbol{u}-\boldsymbol{G}, \boldsymbol{f}=$ $\nu \triangle \boldsymbol{G}-(\nabla \boldsymbol{G}) \boldsymbol{a}-\gamma \boldsymbol{G}$ in $\boldsymbol{H}^{-1}(\Omega), g=-\operatorname{div} \boldsymbol{G}$ in $L^{2}(\Omega)$ and $\boldsymbol{h}=\nu \frac{\partial \boldsymbol{G}}{\partial \boldsymbol{n}}$ in $\boldsymbol{H}^{1 / 2}\left(\Gamma_{N}\right)$, the model problem can be written in a new equivalent form.

$$
\begin{align*}
& \text { minimize } \quad J(\gamma, \boldsymbol{u})=\frac{1}{2}\left\|\boldsymbol{v}+\boldsymbol{G}-\boldsymbol{u}_{R}\right\|_{0, \omega}^{2}+\frac{\alpha}{2}\|\gamma\|_{s, \Omega}^{2}  \tag{3}\\
& \text { subject to } \quad-\nu \triangle \boldsymbol{v}+(\nabla \boldsymbol{v}) \boldsymbol{a}+\nabla p+\gamma \boldsymbol{v}=\boldsymbol{f} \quad \text { in } \boldsymbol{H}^{-1}(\Omega)  \tag{4}\\
& \operatorname{div} \boldsymbol{v}=g \quad \text { in } \Omega \\
& \boldsymbol{v}=\mathbf{0} \quad \text { on } \Gamma_{D} \\
& -\nu \frac{\partial \boldsymbol{v}}{\partial \boldsymbol{n}}+p \boldsymbol{n}=\boldsymbol{h} \quad \text { on } \Gamma_{N} \\
& \boldsymbol{v} \in \boldsymbol{H}^{1}(\Omega), \gamma \in H^{s}(\Omega) \\
& 0 \leq \gamma \leq M \text { a.e. in } \Omega \text {. }
\end{align*}
$$

Remark 1. The boundary condition imposed in $\Gamma_{N}$ strongly depends of $\boldsymbol{a} \cdot \boldsymbol{n} \geq 0$, but that hypothesis can be restrictive for some existing phenomena in the applied case. For example, if a vortex crosses the outlet, there may be a region where the flow returns to the domain. To overcome this problem, the so called "backflow" stabilization techniques [8] can be applied. For example, directional do-nothing condition can be used for this problem. This condition verifies

$$
-\nu \frac{\partial \boldsymbol{u}}{\partial \boldsymbol{n}}+p \boldsymbol{n}+\frac{1}{2}(\boldsymbol{a} \cdot \boldsymbol{n})_{-} \boldsymbol{u}=0
$$

where $(x)_{-}=\min \{x, 0\}$. For this case, the existence of a solution of the Oseen equations is verified [11.

## 3. Existence of solution

The demonstration of the existence of an optimal solution presented follows the same scheme as [4]. A first step is to introduce some helpful notations.

Definition 2. Let $s \geq 0$. The set of admissible parameter functions is defined by

$$
\mathcal{A}=\left\{\gamma \in H^{s}(\Omega) \mid 0 \leq \gamma \leq M\right\}
$$

Definition 3. The space $\boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega)$, subspace of $\boldsymbol{H}^{1}(\Omega)$, is defined by

$$
\boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega)=\left\{\boldsymbol{v} \in \boldsymbol{H}^{1}(\Omega) \mid \boldsymbol{v}=\mathbf{0} \text { on } \Gamma_{D}\right\}
$$

Definition 4. Let $\gamma \in L^{\infty}(\Omega)$, the map $A: L^{\infty}(\Omega) \rightarrow \boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega) \times L_{0}^{2}(\Omega)$ is defined by $A(\gamma)=\left[A_{1}(\gamma), A_{2}(\gamma)\right]=$ $[\boldsymbol{v}, p]$, where $\boldsymbol{v} \in \boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega)$ and $p \in L_{0}^{2}(\Omega)$ are the unique variational solutions of the Problem (4). In what follows, it is denoted $\mathcal{H}=\boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega) \times L_{0}^{2}(\Omega)$, which is a Banach space behind the norm

$$
\|[\boldsymbol{v}, p]\|_{\mathcal{H}}=\nu|\boldsymbol{v}|_{1, \Omega}+\|p\|_{0, \Omega}
$$

Remark 5. It is possible to ensure that $A$ is well defined, because (4) has an unique solution for every $\gamma \in \mathcal{A}$ (see Lemma 5.8 on [22]). Futhermore, there exists a constant $C>0$, independent of $\gamma, \nu, \boldsymbol{u}_{D}$ and $\boldsymbol{G}$, such that

$$
\begin{aligned}
\nu|\boldsymbol{v}|_{1, \Omega}^{2}+(\gamma \boldsymbol{v}, \boldsymbol{v})_{0, \Omega} & \leq C \max \left\{\nu, \frac{\|\boldsymbol{a}\|_{\infty}^{2}}{\nu},\|\gamma\|_{\infty, \Omega}\right\}\|\boldsymbol{G}\|_{1, \Omega}^{2} \\
\|p\|_{0, \Omega}^{2} & \leq C \max \left\{\nu,\|\boldsymbol{a}\|_{\infty, \Omega}, \frac{\|\boldsymbol{a}\|_{\infty, \Omega}^{2}}{\nu},\|\gamma\|_{\infty, \Omega}\right\}\|\boldsymbol{G}\|_{1, \Omega}
\end{aligned}
$$

The uniformly boundedness of $\|\boldsymbol{v}\|_{1, \Omega}$ and $\|p\|_{0, \Omega}$ are obtained because $\|\gamma\|_{\infty, \Omega} \leq M$.
Theorem 6. Problem (3) has at least one solution $\gamma^{*} \in \mathcal{A}$ with optimal states $A\left(\gamma^{*}\right)=\left[\boldsymbol{v}^{*}, p^{*}\right]$, i.e.

$$
(\forall \gamma \in \mathcal{A}) \quad J\left(\gamma, A_{1}(\gamma)\right) \geq J\left(\gamma^{*}, \boldsymbol{v}^{*}\right)
$$

Proof. First, $J(\gamma)=J(\gamma, A(\gamma))$ is bounded below by 0 and

$$
(\forall \gamma \in \mathcal{A}) \quad J(\gamma) \geq \frac{\alpha}{2}\|\gamma\|_{s, \Omega}^{2}
$$

Thus, $m=\inf _{\gamma \in \mathcal{A}} J(\gamma)$ is well defined. Let $\left\{\gamma_{n}\right\}_{n \in \mathbb{N}} \subseteq \mathcal{A}$ and $\left\{\left[\boldsymbol{v}_{n}, p_{n}\right]\right\}_{n \in \mathbb{N}} \subseteq \boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega) \times L_{0}^{2}(\Omega)$ such that $A\left(\gamma_{n}\right)=\left[\boldsymbol{v}_{n}, p_{n}\right],\left\{J\left(\gamma_{n}\right)\right\}$ is decreasing and $J\left(\gamma_{n}\right) \rightarrow m$. It is clear that

$$
(\forall n \in \mathbb{N}) \quad J\left(\gamma_{n}\right) \leq J\left(\gamma_{1}\right)
$$

Then, $\left\{\gamma_{n}\right\}_{n \in \mathbb{N}}$ is uniformly bounded in $H^{s}(\Omega)$. Since $\mathcal{A}$ is weakly closed, there exists a subsequence (denoted by $\left\{\gamma_{n}\right\}_{n \in \mathbb{N}}$ ) such that $\gamma_{n} \rightharpoonup \gamma^{*}$ in $H^{s}(\Omega)$, with $\gamma^{*} \in \mathcal{A}$. In particular, $\gamma_{n} \rightharpoonup \gamma^{*}$ in $L^{2}(\Omega)$. By the same way, $\left\{\boldsymbol{v}_{n}\right\}_{n \in \mathbb{N}}$ and $\left\{p_{n}\right\}_{n \in \mathbb{N}}$ are also uniformly bounded (see Remark 5 on $\boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega)$ and $L_{0}^{2}(\Omega)$, respectively. Applying the Sobolev embedding Theorem (see Section 6.6 in 16), there exists a subsequence (also denoted by $\left\{\boldsymbol{v}_{n}\right\}_{n \in \mathbb{N}}$ ) such that $\boldsymbol{v}_{n} \rightharpoonup \boldsymbol{v}^{*}$ weakly in $\boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega)$, with $\boldsymbol{v}^{*} \in \boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega)$, and $\boldsymbol{v}_{n} \rightarrow \boldsymbol{v}^{*}$ in $\boldsymbol{L}^{p}(\Omega)$ for $p \in[2,6)$ when $N \geq 3$, or for $p \geq 2$ when $N=2$. In particular,

$$
\boldsymbol{v}_{n} \rightarrow \boldsymbol{v}^{*} \text { in } \boldsymbol{L}^{4}(\Omega) \text { and } \boldsymbol{L}^{2}(\Omega)
$$

Let $\boldsymbol{w} \in \boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega)$. Since $d \in\{2,3\}$, then $\boldsymbol{w} \in \boldsymbol{L}^{4}(\Omega)$. Later, $\boldsymbol{v}_{n} \cdot \boldsymbol{w} \rightarrow \boldsymbol{v}^{*} \cdot \boldsymbol{w}$ in $L^{2}(\Omega)$ and $\left(\nabla \boldsymbol{v}_{n}\right) \boldsymbol{a} \rightharpoonup(\nabla \boldsymbol{v}) \boldsymbol{a}$ weakly in $L^{2}(\Omega)$. Thus, for every $\boldsymbol{w} \in \boldsymbol{H}_{\Gamma_{D}}^{1}(\Omega)$

$$
\left(\gamma_{n} \boldsymbol{v}_{n}, \boldsymbol{w}\right)_{0, \Omega} \rightarrow\left(\gamma^{*} \boldsymbol{v}^{*}, \boldsymbol{w}\right)_{0, \Omega} \quad\left(\left(\nabla \boldsymbol{v}_{n}\right) \boldsymbol{a}, \boldsymbol{w}\right)_{0, \Omega} \rightarrow\left(\left(\nabla \boldsymbol{v}^{*}\right) \boldsymbol{a}, \boldsymbol{w}\right)_{0, \Omega}
$$

Repeating this argument, there exists a subsequence also denoted by $\left\{p_{n}\right\}_{n \in \mathbb{N}}$ such that

$$
p_{n} \rightharpoonup p^{*} \text { in } L_{0}^{2}(\Omega)
$$

In conclusion, $\left(\boldsymbol{v}^{*}, p^{*}\right) \in \mathcal{H}$ is the variational solution of Equation (4). In other words, $A\left(\gamma^{*}\right)=\left[\boldsymbol{v}^{*}, p^{*}\right]$. Now, $\gamma^{*}$ is optimal. Indeed,

$$
\begin{aligned}
m & =\lim _{n \rightarrow \infty} J\left(\gamma_{n}\right) \\
& \geq \lim _{n \rightarrow \infty} \frac{1}{2}\left\|\boldsymbol{v}_{n}+\boldsymbol{G}-\boldsymbol{u}_{R}\right\|_{0, \omega}^{2}+\frac{\alpha}{2} \liminf _{n \rightarrow \infty}\left\|\gamma_{n}\right\|_{s, \Omega}^{2} \\
& \geq \frac{1}{2}\left\|\boldsymbol{v}^{*}+\boldsymbol{G}-\boldsymbol{u}_{R}\right\|_{0, \omega}^{2}+\frac{\alpha}{2}\left\|\gamma^{*}\right\|_{s, \Omega}^{2}=J\left(\gamma^{*}\right) \geq m
\end{aligned}
$$

Hence, $J\left(\gamma, A_{1}(\gamma)\right) \geq J\left(\gamma^{*}, \boldsymbol{v}^{*}\right)$ for every $\gamma \in \mathcal{A}$, proving this theorem.

## 4. First order necessary optimality condition

In order to implement a descent method to numerically solve this problem, it is necessary to establish the differentiability of functional $J$. However, this directly depends on the differentiability of map $A$. From this result, it is possible to establish necessary optimality conditions.

Theorem 7. The map $A$ is Fréchet-differentiable. For each $\gamma, \gamma_{1} \in \mathcal{A}$, if $A(\gamma)=[\boldsymbol{v}, p]$, then the Fréchet derivative $A^{\prime}(\gamma)\left[\gamma_{1}\right]=\left[\boldsymbol{v}_{1}^{\prime}, p_{1}^{\prime}\right]$ can be described by the weak solution of the following problem

$$
\begin{align*}
-\nu \triangle \boldsymbol{v}_{1}^{\prime}+\left(\nabla \boldsymbol{v}_{1}^{\prime}\right) \boldsymbol{a}+\nabla p_{1}^{\prime}+\gamma \boldsymbol{v}_{1}^{\prime} & =-\gamma_{1}(\boldsymbol{v}+\boldsymbol{G}) & & \text { in } \Omega  \tag{5}\\
\operatorname{div} \boldsymbol{v}_{1}^{\prime} & =0 & & \text { in } \Omega \\
\boldsymbol{v}_{1}^{\prime} & =\mathbf{0} & & \text { on } \Gamma_{D} \\
-\nu \frac{\partial \boldsymbol{v}_{1}^{\prime}}{\partial \boldsymbol{n}}+p_{1}^{\prime} \boldsymbol{n} & =\mathbf{0} & & \text { on } \Gamma_{N} .
\end{align*}
$$

Proof. See Appendix A
Defining $B(\boldsymbol{v})=\frac{1}{2}\left\|\boldsymbol{v}+\boldsymbol{G}-\boldsymbol{u}_{R}\right\|_{0, \omega}^{2}$ and $C(\gamma)=\frac{\alpha}{2}\|\gamma\|_{s, \Omega}^{2}$, an expression for Frechét derivatives of $B$ and $C$ is given by

$$
B^{\prime}(\boldsymbol{v})[w]=\left(\boldsymbol{v}+\boldsymbol{G}-\boldsymbol{u}_{R}, \boldsymbol{w}\right)_{0, \omega} \quad C^{\prime}(\gamma)[\beta]=\alpha(\gamma, \beta)_{s, \Omega}
$$

Applying chain rule, it is obtained that

$$
\begin{aligned}
J^{\prime}(\gamma)\left[\gamma_{1}\right] & =B^{\prime}(\boldsymbol{v})\left[A_{1}^{\prime}(\gamma)\left[\gamma_{1}\right]\right]+C^{\prime}(\gamma)\left[\gamma_{1}\right] \\
& =\left(\boldsymbol{v}+\boldsymbol{G}-\boldsymbol{u}_{R}, \boldsymbol{v}_{1}^{\prime}\right)_{0, \omega}+\alpha\left(\gamma, \gamma_{1}\right)_{s, \Omega}
\end{aligned}
$$

In order to reduce this expression, the following definition is introduced similarly to [1].
Definition 8. Let $\gamma \in \mathcal{A}$ and $A(\gamma)=[\boldsymbol{v}, p]$. The adjoint states $A^{*}(\gamma)=[\boldsymbol{\varphi}, \xi]$ are defined as the unique weak solution of the problem

$$
\begin{align*}
-\nu \triangle \boldsymbol{\varphi}-(\nabla \boldsymbol{\varphi}) \boldsymbol{a}+\nabla \xi+\gamma \boldsymbol{\varphi} & =\chi_{\omega}\left(\boldsymbol{v}+\boldsymbol{G}-\boldsymbol{u}_{R}\right) & & \text { in } \boldsymbol{\Omega}  \tag{6}\\
\operatorname{div} \boldsymbol{\varphi} & =0 & & \text { in } \Omega \\
\boldsymbol{\varphi} & =\mathbf{0} & & \text { on } \Gamma_{D} \\
-\nu \frac{\partial \varphi}{\partial \boldsymbol{n}}+\xi \boldsymbol{n}-(\boldsymbol{a} \cdot \boldsymbol{n}) \boldsymbol{\varphi} & =\mathbf{0} & & \text { on } \Gamma_{N} .
\end{align*}
$$

where $\chi_{\omega}$ is the indicator function of $\omega$.
Using this defintion, it is possible to rewrite $J^{\prime}\left(\gamma_{1}\right)$ depending of the adjoint state. That expression is simpler to analyze, since it depends on the adjoint state, allowing a simple form of a first order optimality condition using a variational inequality.

Theorem 9. Let $\gamma, \gamma_{1} \in \mathcal{A}$ and $s \geq 0$. Then,

$$
J^{\prime}(\gamma)\left[\gamma_{1}\right]=-\left(\gamma_{1}(\boldsymbol{v}+\boldsymbol{G}), \boldsymbol{\varphi}\right)_{0, \Omega}+\alpha\left(\gamma, \gamma_{1}\right)_{s, \Omega}
$$

where $A(\gamma)=[\boldsymbol{v}, p]$. If $\gamma^{*} \in \mathcal{A}$ is an optimal for Problem (3), then

$$
(\forall \gamma \in \mathcal{A}) \quad-\left(\left(\gamma-\gamma^{*}\right)\left(\boldsymbol{v}^{*}+\boldsymbol{G}\right), \boldsymbol{\varphi}\right)_{0, \Omega}+\alpha\left(\gamma^{*}, \gamma-\gamma^{*}\right)_{s, \Omega} \geq 0
$$

where $A\left(\gamma^{*}\right)=\left[\boldsymbol{v}^{*}, p^{*}\right]$ and $A^{*}\left(\gamma^{*}\right)=[\boldsymbol{\varphi}, \xi]$ are the states and adjont states of $\gamma^{*}$, respectively.

Proof. First, using integration by parts with the adjoint states $[\boldsymbol{\varphi}, \xi]$ as tests functions, it is obtained that

$$
\begin{aligned}
-\nu \int_{\Omega} \Delta \boldsymbol{v}_{1}^{\prime} \cdot \boldsymbol{\varphi} d \boldsymbol{x} & =-\nu \int_{\Omega} \Delta \boldsymbol{\varphi} \cdot \boldsymbol{v}_{1}^{\prime} d \boldsymbol{x}+\int_{\partial \Omega} \boldsymbol{\varphi} \cdot\left(-\nu \frac{\partial \boldsymbol{v}_{1}^{\prime}}{\partial \boldsymbol{n}}\right)-\boldsymbol{v}_{1}^{\prime} \cdot\left(-\nu \frac{\partial \boldsymbol{\varphi}}{\partial \boldsymbol{n}}\right) d S \\
\int_{\Omega}\left[\left(\nabla \boldsymbol{v}_{1}^{\prime}\right) \boldsymbol{a}\right] \cdot \boldsymbol{\varphi} d \boldsymbol{x} & =-\int_{\Omega}\left[(\nabla \boldsymbol{a}) \boldsymbol{v}_{1}^{\prime}\right] \cdot \boldsymbol{\varphi} d \boldsymbol{x}+\int_{\partial \Omega}(\boldsymbol{a} \cdot \boldsymbol{n})\left(\boldsymbol{\varphi} \cdot \boldsymbol{v}_{1}^{\prime}\right) d S \\
\int_{\Omega} \nabla p_{1}^{\prime} \cdot \boldsymbol{\varphi} d \boldsymbol{x} & =-\int_{\Omega} p_{1}^{\prime} \operatorname{div} \boldsymbol{\varphi} d \boldsymbol{x}+\int_{\partial \Omega} \boldsymbol{\varphi} \cdot\left(p_{1}^{\prime} \boldsymbol{n}\right) d S \\
-\int_{\Omega} \xi \operatorname{div} \boldsymbol{v}_{1}^{\prime} d \boldsymbol{x} & =\int_{\Omega} \nabla \xi \cdot \boldsymbol{v}_{1}^{\prime} d \boldsymbol{x}-\int_{\partial \Omega} \boldsymbol{v}_{1}^{\prime} \cdot(\xi \boldsymbol{n}) d S
\end{aligned}
$$

Then,

$$
\begin{aligned}
-\left(\gamma_{1}(\boldsymbol{v}+\boldsymbol{G}), \boldsymbol{\varphi}\right)_{0, \Omega} & =-\nu \int_{\Omega} \Delta \boldsymbol{v}_{1}^{\prime} \cdot \boldsymbol{\varphi} d \boldsymbol{x}+\int_{\Omega}\left[\left(\nabla \boldsymbol{v}_{1}^{\prime}\right) \boldsymbol{a}\right] \cdot \boldsymbol{\varphi} d \boldsymbol{x}+\int_{\Omega} \nabla p_{1}^{\prime} \cdot \boldsymbol{\varphi} d \boldsymbol{x}+\int_{\Omega} \gamma \boldsymbol{v}_{1}^{\prime} \cdot \boldsymbol{\varphi} d \boldsymbol{x}-\int_{\Omega} \xi \operatorname{div} \boldsymbol{v}_{1}^{\prime} d \boldsymbol{x} \\
& =-\nu \int_{\Omega} \Delta \boldsymbol{\varphi} \cdot \boldsymbol{v}_{1}^{\prime} d \boldsymbol{x}-\int_{\Omega}[(\nabla \boldsymbol{\varphi}) \boldsymbol{a}] \cdot \boldsymbol{v}_{1}^{\prime} d \boldsymbol{x}+\int_{\Omega} \nabla \xi \cdot \boldsymbol{v}_{1}^{\prime} d \boldsymbol{x}+\int_{\Omega} \gamma \boldsymbol{v}_{1}^{\prime} \cdot \boldsymbol{\varphi} d \boldsymbol{x}-\int_{\Omega} p_{1}^{\prime} \operatorname{div} \boldsymbol{\varphi} d \boldsymbol{x} \\
& +\int_{\Gamma_{N}} \boldsymbol{\varphi} \cdot\left(-\nu \frac{\partial \boldsymbol{v}_{1}^{\prime}}{\partial \boldsymbol{n}}+p_{1}^{\prime} \boldsymbol{n}\right)-\int_{\Gamma_{N}} \boldsymbol{v}_{1}^{\prime} \cdot\left(-\nu \frac{\partial \boldsymbol{\varphi}}{\partial \boldsymbol{n}}+\xi \boldsymbol{n}-(\boldsymbol{a} \cdot \boldsymbol{n}) \boldsymbol{\varphi}\right) d S \\
& =-\nu \int_{\Omega} \triangle \boldsymbol{\varphi} \cdot \boldsymbol{v}_{1}^{\prime} d \boldsymbol{x}-\int_{\Omega}\left[\left(\nabla \boldsymbol{v}_{1}^{\prime}\right) \boldsymbol{a}\right] \cdot \boldsymbol{\varphi} d \boldsymbol{x}+\int_{\Omega} \nabla \xi \cdot \boldsymbol{v}_{1}^{\prime} d \boldsymbol{x}+\int_{\Omega} \gamma \boldsymbol{v}_{1}^{\prime} \cdot \boldsymbol{\varphi} d \boldsymbol{x}-\int_{\Omega} p_{1}^{\prime} \operatorname{div} \boldsymbol{\varphi} d \boldsymbol{x} \\
& =\left(\left(\boldsymbol{v}+\boldsymbol{G}-\boldsymbol{u}_{D}\right), \boldsymbol{v}_{1}^{\prime}\right)_{0, \omega}
\end{aligned}
$$

Thus,

$$
J^{\prime}(\gamma)\left[\gamma_{1}\right]=-\left(\gamma_{1}(\boldsymbol{v}+\boldsymbol{G}), \boldsymbol{\varphi}\right)_{0, \Omega}+\alpha\left(\gamma, \gamma_{1}\right)_{s, \Omega}
$$

Later, if $\gamma^{*} \in \mathcal{A}$ is optimal for the problem, then

$$
(\forall \gamma \in \mathcal{A}) \quad J(\gamma) \geq J\left(\gamma^{*}\right)
$$

Finally, it is obtained that

$$
J^{\prime}\left(\gamma^{*}\right)\left[\gamma-\gamma^{*}\right]=\lim _{\varepsilon \rightarrow 0^{+}} \frac{J\left(\gamma^{*}+\varepsilon\left(\gamma-\gamma^{*}\right)\right)-J\left(\gamma^{*}\right)}{\varepsilon} \geq 0
$$

proving this theorem.

## 5. Second order sufficient optimality condition.

The stability results of the optimization algorithms depend, for the most part, on the existence of the second derivative of $J$. Likewise, it is possible to establish second order sufficient optimality conditions.

In first place, it is necesary to introduce new forms of emebedding inequalties.
Theorem 10. There exists a constant $q^{*}>2$, dependent only of $\Omega$, such that for each $p \in\left[2, q^{*}\right]$ there exists a constant $C>0$, dependent only of $\Omega, M, \nu$ and $p$ such that

$$
\|\boldsymbol{v}\|_{1, p, \Omega} \leq C\|G\|_{-1, p, \Omega}
$$

where $G \in W^{-1, p}(\Omega)$ (the dual of $\left.W_{0}^{1, q}(\Omega)\right)$ is defined by

$$
\left(\forall \boldsymbol{w} \in W_{0}^{1, q}(\Omega)\right) \quad\langle G(\boldsymbol{v}), \boldsymbol{w}\rangle=\langle\boldsymbol{f}, \boldsymbol{w}\rangle_{\boldsymbol{H}^{-1}(\Omega), \boldsymbol{H}^{1}(\Omega)}+(\boldsymbol{h}, \boldsymbol{w})_{0, \Gamma_{N}}
$$

with $q \geq 2$ such that $\frac{1}{p}+\frac{1}{q}=1$.

## Proof. See Theorem 1 in [20].

Remark 11. The hypotheses $\boldsymbol{u}_{R} \in L^{\infty}(\omega)$ and $\boldsymbol{u}_{D} \in L^{\infty}\left(\Gamma_{D}\right)$ appear frequently. Asumming $\boldsymbol{u}_{R} \in L^{\infty}(\omega)$ and $\boldsymbol{u}_{D} \in L^{\infty}\left(\Gamma_{D}\right)$, it is clear that $\boldsymbol{G} \in \boldsymbol{W}^{1, \infty}(\Omega)$. Then, there exists a constant $C_{p}>0$, depending only of $p$ and $\Omega$, such that

$$
\|G\|_{-1, p, \Omega} \leq C_{p} \max \left\{\nu,\|\boldsymbol{a}\|_{\infty, \Omega},\|\gamma\|_{\infty, \Omega}\right\}\|\boldsymbol{G}\|_{1, \infty, \Omega}
$$

Futhermore, using Theorem 10.

$$
\|\boldsymbol{v}\|_{1, p, \Omega} \leq C C_{p} \max \left\{\nu,\|\boldsymbol{a}\|_{\infty, \Omega}, M\right\}\|\boldsymbol{G}\|_{1, \infty, \Omega}
$$

So, $\|\boldsymbol{v}\|_{1, p, \Omega}$ is uniformly bounded for each $p \in\left[2, q^{*}\right]$. Also, Theorem 1 in [20] allows us to be more precise about the asymptotic behavior of $q^{*}$. Indeed, if the value of $\|\gamma\|_{\infty, \Omega}$ increases, $q^{*}$ decreases to 2 .
Lemma 12. Let $s>0$ and $p \in[1, N]$.

1. If $N>s p$, then the embedding from $W^{s, p}(\Omega)$ to $L^{r}(\Omega)$ is continuous for $r \in\left[p, \frac{d}{d-s p}\right]$.
2. If $N=s p$, then the embedding from $W^{s, p}(\Omega)$ to $L^{r}(\Omega)$ is continuous for $r \in[p,+\infty)$.

Proof. See Section 6.6 in 16
Remark 13. Using this embedding result with $s \geq \frac{N}{q^{*}}$ and $p=2$, Lemma 12 can be written by

1. If $q^{*}>2$, then the embedding from $H^{d / q^{*}}(\Omega)$ to $L^{r}(\Omega)$ is compact for $r \in\left[2, \frac{1}{\frac{1}{p}-\frac{1}{q^{*}}}\right]$.
2. If $q^{*}=2$, then the embedding from $H^{d / q^{*}}(\Omega)$ to $L^{r}(\Omega)$ is compact for $r \in[2,+\infty)$.

Then, for $2 \leq p \leq \frac{1}{\frac{1}{2}-\frac{1}{q^{*}}}$ there exists $C>0$ such that

$$
(\forall \gamma \in \mathcal{A}) \quad\|\gamma\|_{0, p, \Omega} \leq C\|\gamma\|_{s, \Omega}
$$

In order to establish and prove second order optimality conditions, a first step is to show that the map $A$ is twice Fréchet-differentiable.

Theorem 14. The map $\gamma \longmapsto A^{\prime}(\gamma)\left[\gamma_{1}\right]$ from $L^{\infty}(\Omega)$ to $\mathcal{H}$ for each $\gamma_{1} \in \mathcal{A}$ is Fréchet-differentiable. Let $\gamma, \gamma_{1}, \gamma_{2} \in \in \mathcal{A}$, the Fréchet derivative of $A^{\prime}(\gamma)\left[\gamma_{1}\right]$ on $\gamma_{2}$ direction is given by $A^{\prime \prime}(\gamma)\left[\gamma_{1}, \gamma_{2}\right]=\left[\boldsymbol{v}^{\prime \prime}, p^{\prime \prime}\right]$, where $\left[\boldsymbol{v}^{\prime \prime}, p^{\prime \prime}\right]$ is the unique weak solution of the problem

$$
\begin{align*}
-\nu \triangle \boldsymbol{v}^{\prime \prime}+\left(\nabla \boldsymbol{v}^{\prime \prime}\right) \boldsymbol{a}+\nabla p^{\prime \prime}+\gamma \boldsymbol{v}^{\prime \prime} & =-\left(\gamma_{2} \boldsymbol{v}_{1}^{\prime}+\gamma_{1} \boldsymbol{v}_{2}^{\prime}\right) & & \text { in } \Omega  \tag{7}\\
\operatorname{div} \boldsymbol{v}^{\prime \prime} & =0 & & \text { in } \Omega \\
\boldsymbol{v}^{\prime \prime} & =\mathbf{0} & & \text { on } \Gamma_{D} \\
-\nu \frac{\partial \boldsymbol{v}^{\prime \prime}}{\partial \boldsymbol{n}}+p^{\prime \prime} \boldsymbol{n} & =\mathbf{0} & & \text { on } \Gamma_{N} .
\end{align*}
$$

and $A^{\prime}(\gamma) \gamma_{j}=\left[\boldsymbol{v}_{j}^{\prime}, p_{j}^{\prime}\right]$, for $j \in\{1,2\}$.
Proof. See Appendix B

Let $\gamma, \gamma_{1}, \gamma_{2} \in \mathcal{A}$. An expresion for the Fréchet second derivative of $J(\gamma)$ on directions $\gamma_{1}$ and $\gamma_{2}$ is given by

$$
\begin{aligned}
J^{\prime \prime}(\gamma)\left[\gamma_{1}, \gamma_{2}\right] & =B^{\prime \prime}(A(\gamma))\left[A_{1}^{\prime}\left(\gamma_{1}\right), A^{\prime}\left(\gamma_{2}\right)\right]+B^{\prime}((A(\gamma)))\left[A_{1}^{\prime \prime}(\gamma)\left[\gamma_{1}, \gamma_{2}\right]\right]+C^{\prime \prime}(\gamma)\left[\gamma_{1}, \gamma_{2}\right] \\
& =\left(\boldsymbol{v}_{1}^{\prime}, \boldsymbol{v}_{2}^{\prime}\right)_{0, \Omega}+\left(\boldsymbol{v}+\boldsymbol{G}-\boldsymbol{u}_{D}, \boldsymbol{v}^{\prime \prime}\right)_{0, \omega}+\alpha\left(\gamma_{1}, \gamma_{2}\right)_{s, \Omega}
\end{aligned}
$$

where, reasoning as in the proof of Theorem 9

$$
\left(\boldsymbol{v}+\boldsymbol{G}-\boldsymbol{u}_{D}, \boldsymbol{v}^{\prime \prime}\right)_{0, \omega}=-\left(\gamma_{1} \boldsymbol{v}_{2}^{\prime}+\gamma_{2} \boldsymbol{v}_{1}^{\prime}, \boldsymbol{\varphi}\right)_{0, \Omega}
$$

In consequence,

$$
\begin{aligned}
J^{\prime \prime}(\gamma)\left[\gamma_{1}, \gamma_{2}\right] & =B^{\prime \prime}(A(\gamma))\left[A_{1}^{\prime}\left(\gamma_{1}\right), A^{\prime}\left(\gamma_{2}\right)\right]+B^{\prime}((A(\gamma)))\left[A_{1}^{\prime \prime}(\gamma)\left[\gamma_{1}, \gamma_{2}\right]\right]+C^{\prime \prime}(\gamma)\left[\gamma_{1}, \gamma_{2}\right] \\
& =\left(\boldsymbol{v}_{1}^{\prime}, \boldsymbol{v}_{2}^{\prime}\right)_{0, \Omega}-\left(\gamma_{1} \boldsymbol{v}_{2}^{\prime}+\gamma_{2} \boldsymbol{v}_{1}^{\prime}, \boldsymbol{\varphi}\right)_{0, \Omega}+\alpha\left(\gamma_{1}, \gamma_{2}\right)_{s, \Omega} .
\end{aligned}
$$

In what follows, a second order optimality condition is proved. For this, a series of technical results are required. Consider $r$ such that $\frac{1}{q^{*}}+\frac{1}{r}=\frac{1}{2}$. For $\gamma, \gamma_{1}, \gamma_{2}, h \in \mathcal{A}$, consider $A(\gamma)=[\boldsymbol{v}, p], A(\gamma+h)=\left[\boldsymbol{v}_{h}, p_{h}\right]$, with respective adjoint states $[\boldsymbol{\varphi}, \xi]$ and $\left[\boldsymbol{\varphi}_{h}, \xi_{h}\right], A^{\prime}(\gamma)\left[\gamma_{k}\right]=\left[\boldsymbol{v}_{k}^{\prime}, p_{k}^{\prime}\right]$ and $\left[\boldsymbol{v}_{h, k}^{\prime}, p_{h, k}^{\prime}\right]=A^{\prime}(\gamma+h)\left[\gamma_{k}\right]$ for $k \in\{1,2\}$. Using this, it is possible to obtain the following estimations.

Lemma 15. For each $q \in\left[2, q^{*}\right]$, there exists $C>0$, independent of $\boldsymbol{v}, \gamma, \gamma_{1}, \gamma_{2}$ and $h$ such that $\left\|\boldsymbol{v}_{h}-\boldsymbol{v}\right\|_{1, q, \Omega} \leq C\|h\|_{\infty, \Omega}$,
Proof. subtracting the equations of $A(\gamma+h)$ and $A(\gamma)$, it is obtained

$$
-\nu \triangle\left(\boldsymbol{v}_{h}-\boldsymbol{v}\right)+\left[\nabla\left(\boldsymbol{v}_{h}-\boldsymbol{v}\right)\right] \boldsymbol{a}+\nabla\left(p_{h}-p\right)+\gamma\left(\boldsymbol{v}_{h}-\boldsymbol{v}\right)=-h\left(\boldsymbol{G}+\boldsymbol{v}_{h}\right)
$$

Using Theorem 10 twice and triangular inequality, there exist $c_{1}, c_{2}>0$ such that

$$
\begin{aligned}
\|\boldsymbol{v}\|_{1, q, \Omega} & \leq c_{1}\left\|h\left(\boldsymbol{G}+\boldsymbol{v}_{h}\right)\right\|_{-1, q, \Omega} \\
& \leq c_{1}\|h\|_{\infty, \Omega}\left\|\left(\boldsymbol{G}+\boldsymbol{v}_{h}\right)\right\|_{1, q, \Omega} \\
& \leq c_{2}\|h\|_{\infty, \Omega}\|\boldsymbol{G}\|_{1, \infty, \Omega} .
\end{aligned}
$$

Taking $C=c_{2}\|\boldsymbol{G}\|_{1, \infty, \Omega}$, the estimation is proved.
Lemma 16. Let $k \in\{1,2\}$. There exists $C>0$ independent of $\boldsymbol{v}, \gamma, \gamma_{1}, \gamma_{2}$ and $h$ such that $\left|\boldsymbol{v}_{k, h}^{\prime}\right|_{1, \Omega} \leq$ $C\left\|\gamma_{k}\right\|_{s, \Omega}$.

Proof. First, testing the equation of $\boldsymbol{v}_{h, k}^{\prime}$ with $\boldsymbol{v}_{h, k}^{\prime}$ and appliying Friedichs-Poincaré and Hölder inequalities, there exists a constant $c_{1}>0$ such that

$$
\begin{aligned}
\nu\left|\boldsymbol{v}_{h, k}^{\prime}\right|_{1, \Omega} & \leq c_{1}\left\|\gamma_{k}\left(\boldsymbol{G}+\boldsymbol{v}_{k, h}\right)\right\|_{0, \Omega} \\
& \leq c_{1}\left\|\gamma_{k}\right\|_{0, r, \Omega}\left\|\left(\boldsymbol{G}+\boldsymbol{v}_{k, h}\right)\right\|_{0, q^{*}, \Omega}
\end{aligned}
$$

From Remark $11,\left\|\left(\boldsymbol{G}+\boldsymbol{v}_{k, h}\right)\right\|_{0, q^{*}, \Omega}$ is uniformly bounded by $c_{2}>0$. Also, due to $r=\frac{1}{\frac{1}{2}-\frac{1}{q^{*}}}$, Remark 13 points out that there exists $c_{3}>0$ such that $\left\|\gamma_{k}\right\|_{0, r, \Omega} \leq c_{3}\left\|\gamma_{k}\right\|_{s, \Omega}$. In conclusion,

$$
\nu\left|\boldsymbol{v}_{h, k}^{\prime}\right|_{1, \Omega} \leq c_{2} c_{3}\left\|\gamma_{k}\right\|_{s, \Omega}
$$

proving this estimation.

Lemma 17. There exists $C>0$ such that

$$
\begin{aligned}
\left|\left(\boldsymbol{v}_{h, 1}^{\prime}-\boldsymbol{v}_{1}^{\prime}, \boldsymbol{v}_{2}^{\prime}\right)_{0, \Omega}\right| & \leq C\left\|\gamma_{1}\right\|_{s, \Omega}\left\|\gamma_{2}\right\|_{s, \Omega}\|h\|_{\infty, \Omega} \\
\left|\left(\boldsymbol{v}_{h, 1}^{\prime}, \boldsymbol{v}_{h, 2}^{\prime}-\boldsymbol{v}_{2}^{\prime}\right)_{0, \Omega}\right| & \leq C\|h\|_{\infty, \Omega}\left\|\gamma_{1}\right\|_{s, \Omega}\left\|\gamma_{2}\right\|_{s, \Omega}
\end{aligned}
$$

Proof. Using Cauchy-Schwartz and Friedrichs-Poincaré inequalities, there exist $c_{1}, c_{2}>0$ such that

$$
\left|\left(\boldsymbol{v}_{h, 1}^{\prime}-\boldsymbol{v}_{1}^{\prime}, \boldsymbol{v}_{2}^{\prime}\right)_{0, \Omega}\right| \leq c_{1}\left|\boldsymbol{v}_{2}^{\prime}\right|_{1, \Omega}\left|\boldsymbol{v}_{h, 1}^{\prime}-\boldsymbol{v}_{1}^{\prime}\right|_{1, \Omega} \leq c_{2}\left\|\gamma_{2}\right\|_{s, \Omega}\left|\boldsymbol{v}_{h, 1}^{\prime}-\boldsymbol{v}_{1}^{\prime}\right|_{1, \Omega}
$$

Using previous lemmas, Friedrichs-Poincaré inequality and Remark 13, there exist $c_{3}, c_{4}, c_{5}, c_{6}>0$ such that

$$
\begin{aligned}
\nu\left|\boldsymbol{v}_{h, 1}^{\prime}-\boldsymbol{v}_{1}^{\prime}\right|_{1, \Omega} & \leq c_{3}\left(\|h\|_{\infty, \Omega}\left|\boldsymbol{v}_{h, 1}^{\prime}\right|_{1, \Omega}+\left\|\gamma_{1}\left(\boldsymbol{v}_{h}-\boldsymbol{v}\right)\right\|_{1, \Omega}\right) \\
& \leq c_{3}\left(c_{4}\|h\|_{\infty, \Omega}\left\|\gamma_{1}\right\|_{s, \Omega}+c_{5}\left\|\gamma_{1}\right\|_{0, r, \Omega}\left\|\left(\boldsymbol{v}_{h}-\boldsymbol{v}\right)\right\|_{1, q^{*}, \Omega}\right) \\
& \leq c_{3}\left(c_{4}+c_{6}\right)\|h\|_{\infty, \Omega}\left\|\gamma_{1}\right\|_{s, \Omega},
\end{aligned}
$$

concluding that

$$
\left|\left(\boldsymbol{v}_{h, 1}^{\prime}, \boldsymbol{v}_{h, 2}^{\prime}-\boldsymbol{v}_{2}^{\prime}\right)_{0, \Omega}\right| \leq C\|h\|_{\infty, \Omega}\left\|\gamma_{1}\right\|_{s, \Omega}\left\|\gamma_{2}\right\|_{s, \Omega}
$$

The proof of the second inequality is similar.

Lemma 18. Let $k, j \in\{1,2\}$, with $j \neq k$. There exists $C>0$ such that

$$
\left|\left(\gamma_{j} \boldsymbol{v}_{h, k}^{\prime}, \boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right)_{0, \Omega}\right| \leq C\|h\|_{\infty, \Omega}\left\|\gamma_{1}\right\|_{s, \Omega}\left\|\gamma_{2}\right\|_{s, \Omega}
$$

Proof. First, applying Hölder and Friedrichs-Poincaré inequalties, there exist $c_{1}, c_{2}>0$ such that

$$
\begin{aligned}
\left|\left(\gamma_{j} \boldsymbol{v}_{h, k}^{\prime}, \boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right)_{0, \Omega}\right| & \leq c_{1}\left\|\gamma_{j}\right\|_{0, r, \Omega}\left\|\boldsymbol{v}_{h, k}^{\prime}\right\|_{1, q^{*}, \Omega}\left|\boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right|_{1, \Omega} \\
& \leq c_{2}\left\|\gamma_{j}\right\|_{s, \Omega}\left\|\gamma_{k}\right\|_{s, \Omega}\left|\boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right|_{1, \Omega}
\end{aligned}
$$

where $\left[\boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}, \xi_{h}-\xi\right]$ verifies

$$
\begin{aligned}
-\nu \triangle\left(\boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right)-\left[\nabla\left(\boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right)\right] \boldsymbol{a}+\nabla\left(\xi_{h}-\xi\right)+\gamma\left(\boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right) & =\boldsymbol{v}_{h}-\boldsymbol{v}-h \boldsymbol{\varphi}_{h} \\
\operatorname{div}\left(\boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right) & =0 .
\end{aligned}
$$

By repeating the previous procedures, since $\left|\boldsymbol{\varphi}_{h}\right|_{1, \Omega}$ is uniformly bounded, there exist $c_{2}, c_{3}>0$ such that

$$
\begin{aligned}
\nu\left|\boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right|_{1, \Omega} & \leq c_{2}\left(\left\|\boldsymbol{v}_{h}-\boldsymbol{v}\right\|_{1, q, \Omega}+\|h\|_{\infty, \Omega}\left|\boldsymbol{\varphi}_{h}\right|_{1, \Omega}\right) \\
& \leq c_{3}\|h\|_{\infty, \Omega}
\end{aligned}
$$

Thus, there exists $C>0$ such that $\left|\left(\gamma_{j} \boldsymbol{v}_{h, k}^{\prime}, \boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right)_{0, \Omega}\right| \leq C\|h\|_{\infty, \Omega}\left|\gamma_{j}\right|_{s, \Omega}\left|\gamma_{k}\right|_{s, \Omega}$.
Theorem 19. Let $\gamma, \gamma_{1}, \gamma_{2}, h \in L^{\infty}(\Omega)$. There exists $L>0$ such that

$$
\left|\left(J^{\prime \prime}(\gamma+h)-J^{\prime \prime}(\gamma)\right)\left[\gamma_{1}, \gamma_{2}\right]\right| \leq L\|h\|_{\infty, \Omega}\left|\gamma_{1}\right|_{s, \Omega}\left|\gamma_{2}\right|_{s, \Omega}
$$

Proof. Applying triangular inequality, we obtain

$$
\begin{aligned}
\left|\left(J^{\prime \prime}(\gamma+h)-J^{\prime \prime}(\gamma)\right)\left[\gamma_{1}, \gamma_{2}\right]\right| & \leq\left|\left(\boldsymbol{v}_{h, 1}^{\prime}, \boldsymbol{v}_{h, 2}^{\prime}-\boldsymbol{v}_{2}^{\prime}\right)_{0, \Omega}\right|+\left|\left(\boldsymbol{v}_{h, 1}^{\prime}-\boldsymbol{v}_{1}^{\prime}, \boldsymbol{v}_{2}^{\prime}\right)_{0, \Omega}\right| \\
& +\left|\left(\gamma_{1} \boldsymbol{v}_{h, 2}^{\prime}, \boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right)_{0, \Omega}\right|+\left|\left(\gamma_{1}\left(\boldsymbol{v}_{h, 2}^{\prime}-\boldsymbol{v}_{2}^{\prime}\right), \boldsymbol{\varphi}\right)_{0, \Omega}\right| \\
& +\left|\left(\gamma_{2} \boldsymbol{v}_{h, 1}^{\prime}, \boldsymbol{\varphi}_{h}-\boldsymbol{\varphi}\right)_{0, \Omega}\right|+\left|\left(\gamma_{2}\left(\boldsymbol{v}_{h, 1}^{\prime}-\boldsymbol{v}_{1}^{\prime}\right), \boldsymbol{\varphi}\right)_{0, \Omega}\right|,
\end{aligned}
$$

where every term were bounded with estimations of the form $C\|h\|_{\infty, \Omega}\left|\gamma_{1}\right|_{s, \Omega}\left|\gamma_{2}\right|_{s, \Omega}$ (see Lemmas 17 and 18.) In conclusion, there exists $L>0$ such that

$$
\left|\left(J^{\prime \prime}(\gamma+h)-J^{\prime \prime}(\gamma)\right)\left[\gamma_{1}, \gamma_{2}\right]\right| \leq L\|h\|_{\infty, \Omega}\left|\gamma_{1}\right|_{s, \Omega}\left|\gamma_{2}\right|_{s, \Omega} .
$$

Corollary 20. There exists $L>0$ such that, for every $\theta \in[0,1]$

$$
\left|\left(J^{\prime \prime}\left(\theta \gamma^{*}+(1-\theta) \gamma\right)-J^{\prime \prime}\left(\gamma^{*}\right)\right)\left[\gamma-\gamma^{*}, \gamma-\gamma^{*}\right]\right| \leq L\left\|\gamma-\gamma^{*}\right\|_{\infty, \Omega}\left|\gamma-\gamma^{*}\right|_{s, \Omega}^{2}
$$

Finally, a second order sufficent optimality condition is presented and proved.
Theorem 21. Let $s \geq \frac{N}{q^{*}}$ and $\gamma^{*} \in \mathcal{A}$ such that $\gamma^{*}$ verifies the first orden optimality condition. If there exists $\delta>0$ such that

$$
(\forall \gamma \in \mathcal{A})
$$

$$
J^{\prime \prime}\left(\gamma^{*}\right)\left[\gamma-\gamma^{*}, \gamma-\gamma^{*}\right] \geq \delta\left\|\gamma-\gamma^{*}\right\|_{s, \Omega}^{2}
$$

Then, there exist $\sigma, \varepsilon>0$, independent of $\gamma$ and $\gamma^{*}$, such that

$$
(\forall \gamma \in \mathcal{A}) \quad\left\|\gamma-\gamma^{*}\right\|_{\infty, \Omega} \leq \varepsilon \Rightarrow J(\gamma) \geq J\left(\gamma^{*}\right)+\sigma\left\|\gamma-\gamma^{*}\right\|_{s, \Omega}^{2}
$$

In consequence, $J$ has a local minimum at $\gamma^{*}$.
Proof. Applying a Taylor expansion, there exists $\theta \in(0,1)$ such that

$$
J(\gamma)=J\left(\gamma^{*}\right)+J^{\prime}\left(\gamma^{*}\right)\left[\gamma-\gamma^{*}\right]+\frac{1}{2} J^{\prime \prime}\left(\theta \gamma^{*}+(1-\theta) \gamma\right)\left[\gamma-\gamma^{*}, \gamma-\gamma^{*}\right]
$$

Using Corollary 20 , if $\varepsilon \leq \frac{\delta}{2 L}, \sigma=\frac{\delta}{4}$ and $\left\|\gamma-\gamma^{*}\right\|_{\infty, \Omega} \leq \varepsilon$, then

$$
\begin{aligned}
J(\gamma) & =J\left(\gamma^{*}\right)+J^{\prime}\left(\gamma^{*}\right)\left[\gamma-\gamma^{*}\right]+\frac{1}{2} J^{\prime \prime}\left(\theta \gamma^{*}+(1-\theta) \gamma\right)\left[\gamma-\gamma^{*}, \gamma-\gamma^{*}\right] \\
& \geq J\left(\gamma^{*}\right)+\frac{1}{2} J^{\prime \prime}\left(\theta \gamma^{*}+(1-\theta) \gamma\right)\left[\gamma-\gamma^{*}, \gamma-\gamma^{*}\right] \\
& \geq J\left(\gamma^{*}\right)+\frac{1}{2}\left(J^{\prime \prime}\left(\theta \gamma^{*}+(1-\theta) \gamma\right)-J^{\prime \prime}\left(\gamma^{*}\right)\right)\left[\gamma-\gamma^{*}, \gamma-\gamma^{*}\right]+\frac{1}{2} J^{\prime \prime}\left(\gamma^{*}\right)\left[\gamma-\gamma^{*}, \gamma-\gamma^{*}\right] \\
& \geq J\left(\gamma^{*}\right)+\frac{\delta}{2}\left\|\gamma-\gamma^{*}\right\|_{s, \Omega}^{2}-\frac{L}{2}\left\|\gamma-\gamma^{*}\right\|_{\infty, \Omega}\left\|\gamma-\gamma^{*}\right\|_{s, \Omega}^{2} \\
& \geq J\left(\gamma^{*}\right)+\frac{\delta}{2}\left\|\gamma-\gamma^{*}\right\|_{s, \Omega}^{2}-\frac{\delta}{4}\left\|\gamma-\gamma^{*}\right\|_{s, \Omega}^{2}=J\left(\gamma^{*}\right)+\sigma\left\|\gamma-\gamma^{*}\right\|_{s, \Omega}^{2},
\end{aligned}
$$

proving the theorem.
Remark 22. The result obtained in Theorem 21 is conditioned to $s \geq \frac{N}{q^{*}}$. Considering the most pessimistic case, when $q^{*}=2$, it is possible to establish that the simplest spaces for $\gamma^{*}$, where Theorem 21 is valid, are $H^{1}(\Omega)$ and $H^{2}(\Omega)$ when $N=2$ and $N=3$, respectively.

## 6. Numerical results

In this section, the previous analysis is complemented by numerical experiments for the new parameter identification problem. Realistic synthetic cases are analyzed in 2D, which represent a longitudinal section of the structure of a cardiac valve in an arbitrary position. However, the numerical results presented below were obtained using the Navier-Stokes equation. So, the numerical problem to be solved is given by

$$
\begin{array}{rlrl}
\operatorname{minimize} & J(\gamma, \boldsymbol{u})=\frac{1}{2}\left\|\boldsymbol{u}-\boldsymbol{u}_{R}\right\|_{0, \omega}^{2}+\frac{\alpha}{2}\|\gamma\|_{1, \Omega}^{2} & & \\
\text { subject to } & -\nu \triangle \boldsymbol{u}+(\nabla \boldsymbol{u}) \boldsymbol{u}+\nabla p+\gamma \boldsymbol{u} & =\mathbf{0} &  \tag{9}\\
\text { in } \Omega \\
\operatorname{div} \boldsymbol{u} & =0 & & \text { in } \Omega \\
\boldsymbol{u} & =\boldsymbol{u}_{D} & & \text { on } \Gamma_{D} \\
-\nu \frac{\partial \boldsymbol{u}}{\partial \boldsymbol{n}}+p \boldsymbol{n} & =\mathbf{0} & & \text { on } \Gamma_{N} \\
\boldsymbol{u} \in \boldsymbol{H}^{1}(\Omega), \gamma \in H^{1}(\Omega) & & \\
0 \leq \gamma \leq M \text { a.e. in } \Omega & &
\end{array}
$$

This special formulation of $(\sqrt{9})$ is similar to $(2)$ in terms of the resistance term. In the following subsections, the configurations of the reference case is explained, as well as the numerical solutions of the inverse problems associated with MRI images or Doppler ultrasound. Due to our lack of real images, our experiments are fully synthetic.

### 6.1. Reference test

First, a reference geometry is defined. This geometry $\Omega$ represents the area around the aortic valve. The inflow $\Gamma_{I}$ has a parabolic profile following Poiseuille's Law given by

$$
\boldsymbol{u}_{D}(x, y)=-U x(d-x) \boldsymbol{n}
$$

where $\boldsymbol{x}=(x, y)$ are the cartesian coordinates of the domain, $\boldsymbol{n}$ is the outer normal vector and $d$ is the diameter of the inflow, while do-nothing conditions are imposed on the outflow $\Gamma_{O}$ given by

$$
-\nu \frac{\partial \boldsymbol{u}}{\partial \boldsymbol{n}}+p \boldsymbol{n}=\mathbf{0}
$$

In the walls of the structure, represented by $\Gamma_{W}$, the fluid has no-slip conditions given by $\boldsymbol{u}=\mathbf{0}$.


Figure 1: Domain for the numerical tests.

The valves are modeled on the resistance term $\gamma \boldsymbol{u}$ using the function $\gamma$. This function assumes a constant value $M \gg 1$ in the regions where the valve is and assumes the value 0 where the valve is not. In order to define $\gamma$, a parabolic figure is drawn on each side $\Omega$ with an approximate thickness of 1 mm . When the valves are fully closed, they are symmetrical with respect to the vertical axis of symmetry. When the valves are open, these parabolas are rotated with respect to a reference system whose origin is at the point where the valve coincides with the walls of the structure given by $\Gamma_{W}$. Since blood flow is modeled, the kinematic viscosity is considered equal to $\nu=0.035 \mathrm{~cm}^{2} / \mathrm{s}$, the density is assumed as $\rho=1 \mathrm{~g} / \mathrm{cm}^{3}$, and the dimensions $d=2 \mathrm{~cm}$ and $U=30 \mathrm{~cm} / \mathrm{s}$, resulting in a peak Reynolds number on the inlet of

$$
\operatorname{Re}=\frac{U d}{\nu}=1714
$$

The Navier-Stokes equations are discretized using the finite element method (FEM) with Taylor-Hood elements $\left(\mathbb{P}_{2}\right.$ for the velocity $\boldsymbol{u}$ and $\mathbb{P}_{1}$ for the pressure) on an unstructured triangular mesh. The mesh used was generated by domain triangulation with $h=0.05 \mathrm{~cm}$, which corresponds to 9950 elements and 4976 nodes. The solver is implemented using the finite element library FEniCS [2] with the default configuration. To solve the nonlinear problems, a Newton's method was used. The resistance $\gamma$ is defined as a discontinuous function with discrete values in each element, assuming the value of $M=10^{8}$ if the element intersects the valve or, otherwise, assumes the value 0 . We define the set $\mathcal{O}$, that represents the valve inside of $\Omega$, by

$$
\mathcal{O}=\left\{\boldsymbol{x} \in \Omega \mid \gamma(\boldsymbol{x})=10^{8}\right\} .
$$



Figure 2: Plots of unstructured mesh (left) and $\gamma$ (right).


Figure 3: Reference solutions $\boldsymbol{u}$ and $p$.

### 6.2. Numerical solution of the inverse problem

Using the solution of the reference test as reference solution $\boldsymbol{u}_{R}$, the following version of the model problem is solved numerically using FEniCS and dolfin-adjoint

$$
\begin{array}{rlrl}
\operatorname{minimize} & J(\gamma, \boldsymbol{u})=\frac{1}{2}\left\|\boldsymbol{u}-\boldsymbol{u}_{R}\right\|_{0, \omega}^{2}+\frac{\alpha}{2}\|\gamma\|_{0, \Omega}^{2}+\frac{\beta}{2}|\gamma|_{1, \Omega}^{2} & & \\
-\nu \triangle \boldsymbol{u}+(\nabla \boldsymbol{u}) \boldsymbol{u}+\nabla p+\gamma \boldsymbol{u} & =\mathbf{0} & & \text { in } \Omega  \tag{11}\\
\text { subject to } & \operatorname{div} \boldsymbol{u} & =0 & \\
\text { in } \Omega \\
\boldsymbol{u} & =\boldsymbol{u}_{D} & & \text { on } \Gamma_{I} \\
\boldsymbol{u} & =\mathbf{0} & & \text { on } \Gamma_{W} \\
-\nu \frac{\partial \boldsymbol{u}}{\partial \boldsymbol{n}}+p \boldsymbol{n} & =\mathbf{0} & & \text { on } \Gamma_{N} \\
& \boldsymbol{u} \in \boldsymbol{H}^{1}(\Omega), \gamma \in H^{1}(\Omega) & & \\
0 \leq \gamma \leq C \text { a.e. in } \Omega . & &
\end{array}
$$

For this example, we considered a measurement area $\omega=\Omega$ and the values $C=10^{3}, \alpha=10^{-4}$ and $\beta=10^{-8}$. The use of two different weights for the norm and seminorm is consistent with the theoretical analysis of the previous sections, so this problem has a solution. The dolfin-adjoint library 26 allows to implement automatic derivation of the discrete adjoint equations for pde models and implement minimization algorithms from the Python 3 libraries. In particular, the L-BFGS-B algorithm (see Section 4.3 in [28]) was used with the following stopping criteria on the step $k$

$$
\frac{\left|J\left(\gamma_{k}\right)-J\left(\gamma_{k+1}\right)\right|}{\max \left\{\left|J\left(\gamma_{k}\right)\right|,\left|J\left(\gamma_{k+1}\right)\right|, 1\right\}} \leq 10^{-6} \quad \text { or } \quad \max _{j \in\{1, \ldots, n\}}\left\{\left|\left(\nabla J_{k}\right)_{j}\right|\right\} \leq 10^{-6}
$$

where $\left(\nabla J_{k}\right)_{j}$ is the $j$-th component of the projected gradient on the step $k$. To start the algorithm, $\gamma_{0}=0$ was used as the inital solution.

As a way to define a valve reconstruction algorithm sketch, we follow these steps.

1. We defined an axis that crosses the domain from the inflow to the outflow.
2. For a uniform discretization of the axis, we defined perpendicular lines.
3. The solution $\gamma^{*}$ obtained by the algorithm is interpolated on each of the lines. Three points are selected on each side of the axis. The first and second point are the limits of an interval where $\nabla \gamma^{*} \cdot \boldsymbol{n}$ has the maximim positive values with $1 \%$ of tolerance. The third point is the local maximum closest to the interval.
4. An average is obtained between the three points.
5. A polyline is drawn on each side of the axis. Each polyline passes through all the average points.

Numerical results are presented in Figure 4. The polyline is drawn in white, which presents a great approximation to the interface between the different values of the reference given by $\gamma$.

The optimal $\gamma^{*}$ has values close to 0 in the areas before and after the valves. On the other hand, in the interior area where there are no valves, the optimal solution takes values close to 0 . Likewise, the magnitude and direction of $\boldsymbol{u}^{*}$ is similar to $\boldsymbol{u}_{\boldsymbol{R}}$, where $\boldsymbol{u}^{*}$ corresponds to the optimal state.


Figure 4: Real $\gamma$ and reconstructed valve, optimal $\gamma$, optimal $\boldsymbol{u}$, and reference solution $\boldsymbol{u}_{R}$ (from left to right). Reference test, 148 iterations.

It is necessary to corroborate that this algorithm is able to solve the inverse problem measuring only a part of $\boldsymbol{u}_{R}$ given by the reference velocity. In this case, choosing $\omega$ as the area where the valves should be (see Figure 5 right), given by

$$
\omega=\Omega \cap\left\{(x, y) \in \mathbb{R}^{2} \mid 1 \leq y \leq 3\right\}
$$

the expected result should be similar to that found above.
Figure 5 presents the numerical results, the polyline and the references. In the case of the reference solution $\boldsymbol{u}_{R}$, a pink rectangle was drawn that allows delimiting the measurement area $\omega$. In particular, this problem required more iterations than the case with measurements on $\Omega$, obtaining very similar results.


Figure 5: Real $\gamma$ and reconstructed valve, optimal $\gamma$, optimal $\boldsymbol{u}$, and reference solution $\boldsymbol{u}_{R}$ (from left to right). Reference test with subdomain, 240 iterations.

### 6.3. Measurements of MRI type

Using the reference solutions, it is possible to generate synthetic measurements that represent the behavior of MRI velocity type measurements. A 2D plane is chosen on which the velocity is measured in one specified direction $\boldsymbol{d} \in \mathbb{R}^{2}$, with $|\boldsymbol{d}|=1$. Then, the measurement is given by

$$
u_{R}=\boldsymbol{u}_{R} \cdot \boldsymbol{d}
$$

The MRI type velocity measurement data is represented in a mesh of uniform quadrilaterals of $2 \mathrm{~mm} \times 2 \mathrm{~mm}$. They are obtained by projecting $u_{R}$ to the $\mathbb{Q}_{0}$ finite element space, given by piecewise constant discontinuous
functions on the quadrilateral mesh. In practical applications, the input speed is unknown and must be estimated. Then, a direction $\boldsymbol{d}$ for this MRI can be chosen as that which is an inner normal to $\Gamma_{I}$. To determine $U$, the projection of $u_{R}$ in $L_{2}\left(\Gamma_{I}\right)$ with respect to $U x(2-x)$ is calculated by

$$
U=\frac{\int_{\Gamma_{I}} u_{R}[x(2-x)] d S}{\int_{\Gamma_{I}}[x(2-x)]^{2} d S} .
$$

Thus, $U$ approaches in a least squares sense. Figure 6 shows the synthetic MRI generated from the reference solution.


Figure 6: Reference $\boldsymbol{u}_{R}$ (left) and synthetic MRI type velocity measurement ( $u_{R}$, right).

The new problem to solve is given by

$$
\begin{array}{rlrl}
\operatorname{minimize} & J(\gamma, \boldsymbol{u})=\frac{1}{2}\left\|\left(\boldsymbol{u}-\boldsymbol{u}_{R}\right) \cdot \boldsymbol{d}\right\|_{0, \Omega}^{2}+\frac{\alpha}{2}\|\gamma\|_{0, \Omega}^{2}+\frac{\beta}{2}|\gamma|_{1, \Omega}^{2} & \\
-\nu \triangle \boldsymbol{u}+(\nabla \boldsymbol{u}) \boldsymbol{u}+\nabla p+\gamma \boldsymbol{u} & =\mathbf{0} & & \text { in } \Omega  \tag{13}\\
\text { subject to } & \operatorname{div} \boldsymbol{u}=0 & \text { in } \Omega \\
\boldsymbol{u} & =\boldsymbol{u}_{D} & & \text { on } \Gamma_{I} \\
\boldsymbol{u} & =\mathbf{0} & & \text { on } \Gamma_{W} \\
& -\nu \frac{\partial \boldsymbol{u}}{\partial \boldsymbol{n}}+p \boldsymbol{n}=\mathbf{0} & & \text { on } \Gamma_{N} \\
\boldsymbol{u} \in \boldsymbol{H}^{1}(\Omega), \gamma \in H^{1}(\Omega) & \\
0 \leq \gamma \leq C \text { a.e. in } \Omega . &
\end{array}
$$

where the values $C=10^{3}, \alpha=10^{-4}$ and $\beta=10^{-8}$ were used. It is possible to prove the existence of solution of this problem in the same way as in the proof of the theorem. Figure 7shows the numerical results and the references. The parameter $\gamma$ has some jumps that coincide with the vertical lines between the MRI voxels, while the polyline tends to have segments parallel to those vertical lines, but it acceptably approximates the space between the valves.


Figure 7: Real $\gamma$ and reconstructed valve, optimal $\gamma$, optimal $\boldsymbol{u}$, and synthetic MRI (from left to right). MRI noiseless.

The white noise intensity in the velocity measurements from MRI is proportional to the velocity encoding parameter (also called VENC [17]) of the scan. This parameter is configured with a value greater than the maximum expected velocity, in order to eliminate signal aliasing. Then, the noise in all voxels is proportional to the maximum velocity in the measurement area. In the clinical practice it can be expected that highquality MRI contains a velocity noise of $10 \%$ of the maximum velocity in each voxel [17]. Gaussian noises were added to this MRI with a standard deviation of $5 \%, 10 \%$ and $20 \%$ of the maximum value of $u_{R}$.

Figure 8 shows the results of this experiment with a $5 \%$ of Guassian noise, but changing the weights to $\alpha=10^{-4}$ and $\beta=10^{-6}$ in terms to decrease the effects of noise. The results are similar to the experiment without noise in terms to the tendency of the polyline to approximate the valve shape and draw lines parallel to the voxels.


Figure 8: Real $\gamma$ and reconstructed valve, optimal $\gamma$, optimal $\boldsymbol{u}$, and synthetic MRI (from left to right). MRI with $5 \%$ of noise.

This approximation seems weaker as noise increases, in the sense that the polyline has a lower quality in its approximation and that the gamma function tends to overfit the data. The following two figures show the result of the experiment with $10 \%$ and $20 \%$ of Gaussian noise, respectively. The noise is exactly the same than the $5 \%$ Gaussian noise case, but increasing the level noise.


Figure 9: Real $\gamma$ and reconstructed valve, optimal $\gamma$, optimal $\boldsymbol{u}$, and synthetic MRI (from left to right). MRI with $10 \%$ of noise.


Figure 10: Real $\gamma$ and reconstructed valve, optimal $\gamma$, optimal $\boldsymbol{u}$, and synthetic MRI (from left to right). MRI with $20 \%$ of noise.

Table 1 shows the mean square error (MSE) between the reconstructed valve given by the polylines obtained using MRI in Figures $7,8,9$ and 10 , and the polyline obtained in the reference test (see Figure 4). To quantify this error, we consider only the points of the polyline in Figure 4 that are at a distance less than or equal to 0.5 mm from $\mathcal{O}$. There are minor differences between the valve reconstructions for the cases with a noise level of $5 \%$ and $10 \%$. However, the quality of the reconstruction decreases when the level noise increases up to $20 \%$.

| Noise level | MSE |
| :---: | :---: |
| $0 \%$ | $4.4680 \cdot 10^{-3}$ |
| $5 \%$ | $8.6308 \cdot 10^{-3}$ |
| $10 \%$ | $8.1992 \cdot 10^{-3}$ |
| $20 \%$ | $1.5016 \cdot 10^{-2}$ |

Table 1: MSE of reconstructed valves using MRI with different noise levels.

Figure 11 shows the reconstructed polyline for three experiments with independent noises at same level of noise. Each polyline color represents the final result of this experiment with a Gaussian noise independent
of the others, whose amplitude was adjusted for the noise level of $5 \%, 10 \%$ and $20 \%$. We can see that the rebuilt valves are more irregular as the noise level increases, especially for $20 \%$ of level noise, which is consistent with Table 1.


Figure 11: Comparison of reconstructed valves with three different noises, with $5 \%, 10 \%$ and $20 \%$ of noise (from left to right).

### 6.4. Measurements of ultrasound imaging type

It is possible to generate synthetic measurements similar to ultrasound images. From a given focal point, directional velocities are measured at each point in a domain with the form of a circular sector with center in the focal point. The chosen directions $\boldsymbol{d}(\boldsymbol{x})$ are given by unit vectors parallel to the vector that joins the measurement point with the focal point. Then, the measurement is given by

$$
u_{R}=\boldsymbol{u}_{R} \cdot \boldsymbol{d}(x)
$$

The measurement data is represented in a structured triangular mesh of 2116 nodes and 4050 triangles with $h=0.067 \mathrm{~cm}$. The measurements are obtained by interpolation of $u_{R}$ to the $\mathbb{P}_{1}$ finite element space defined on the structured mesh. Again, the input speed is unknown and must be estimated from a preliminary measurement. To determine $U$, the projection of $u_{R}$ in $L_{2}\left(\Gamma_{I}\right)$ with respect to $U x(2-x)[\boldsymbol{n} \cdot \boldsymbol{d}(x)]$ is calculated by

$$
U=\frac{\int_{\Gamma_{I}} u_{R}[x(2-x)][\boldsymbol{n} \cdot \boldsymbol{d}(\boldsymbol{x})] d S}{\int_{\Gamma_{I}}[x(2-x)]^{2}[\boldsymbol{n} \cdot \boldsymbol{d}(\boldsymbol{x})]^{2} d S}
$$



Figure 12: Ultrasound imaging domain, mesh and synthetic ultrasound imaging.

The new problem to solve is given by

$$
\begin{array}{rlrl}
\operatorname{minimize} & J(\gamma, \boldsymbol{u})=\frac{1}{2}\left\|\left(\boldsymbol{u}-\boldsymbol{u}_{R}\right) \cdot \boldsymbol{d}(\boldsymbol{x})\right\|_{0, \omega}^{2}+\frac{\alpha}{2}\|\gamma\|_{0, \Omega}^{2}+\frac{\beta}{2}|\gamma|_{1, \Omega}^{2} & & \\
-\nu \triangle \boldsymbol{u}+(\nabla \boldsymbol{u}) \boldsymbol{u}+\nabla p+\gamma \boldsymbol{u} & =\mathbf{0} & & \text { in } \Omega \\
\text { subject to } & \operatorname{div} \boldsymbol{u}=0 & & \text { in } \Omega  \tag{15}\\
\boldsymbol{u} & =\boldsymbol{u}_{D} & & \text { on } \Gamma_{I} \\
\boldsymbol{u} & =\mathbf{0} & & \text { on } \Gamma_{W} \\
& & & \\
& -\nu \frac{\partial \boldsymbol{u}}{\partial \boldsymbol{n}}+p \boldsymbol{n}=\mathbf{0} & & \text { on } \Gamma_{N} \\
\boldsymbol{u} \in \boldsymbol{H}^{1}(\Omega), \gamma \in H^{1}(\Omega) & & \\
0 \leq \gamma \leq C \text { a.e. in } \Omega & &
\end{array}
$$

where $\omega$ is the subdomain on $\Omega$ covered by the ultrasound imaging and the values $C=10^{3}, \alpha=10^{-4}$ and $\beta=10^{-8}$ were used. This problem has solution, and the proof of the existence of solution is similar to the proof of the theorem. Figure 13 presents the numerical results, the polyline and the references. The results obtained are very similar to those obtained in the reference tests with measurements in the entire domain and in a subdomain.


Figure 13: Real $\gamma$ and reconstructed valve, optimal $\gamma$, optimal $\boldsymbol{u}$, and synthetic ultrasound imaging(from left to right).

There are multiple sources of noise in ultrasound images, such as reverberation, ghosting, or fake echo, among others. Therefore, it is not possible to define a simplified structure of noise in this type of images [10, 24. In order to study the effect of the noise, we assume an additive Gaussian noise that is proportional to the maximum of $\left|u_{R}\right|$ in the measurement area. From [23], an experiment with a straight-vessel phantom had a $9.6 \%$ of standard deviation. Then, is acceptable to simulate an ultrasound imaging with a velocity noise of $10 \%$ of the maximum velocity in each node. Gaussian noises were added to this ultrasound imaging with a standard deviation of $10 \%$ and $20 \%$ of the maximum of $\left|u_{R}\right|$.

Figures 14 and 15 show the results of this experiment with a $10 \%$ and $20 \%$ of Guassian noise, respectively, but changing the weights to $\alpha=10^{-4}$ and $\beta=10^{-6}$ in terms to decrease the effects of noise. The results are similar to the experiment without noise in terms to the tendency of the polyline to approximate the valve shape.


Figure 14: Real $\gamma$ and reconstructed valve, optimal $\gamma$, optimal $\boldsymbol{u}$, and synthetic ultrasound imaging(from left to right) with $10 \%$ of noise.


Figure 15: Real $\gamma$ and reconstructed valve, optimal $\gamma$, optimal $\boldsymbol{u}$, and synthetic ultrasound imaging(from left to right) with $20 \%$ of noise.

Table 2 shows the mean square error (MSE) between the reconstructed valve given by the polylines obtained using MRI in Figures 13,14 and 15 , and the polyline obtained in the reference test with measurement in a subdomain (see Figure 5). The MSE was quantified by the same way as in Table 1. There are minor differences between the valve reconstructions for the cases with a noise level of $10 \%$ and $20 \%$.

| Noise level | MSE |
| :---: | :---: |
| $0 \%$ | $3.0985 \cdot 10^{-3}$ |
| $10 \%$ | $8.1993 \cdot 10^{-3}$ |
| $20 \%$ | $8.6772 \cdot 10^{-3}$ |

Table 2: MSE of reconstructed valves using ultrasound imaging with different noise levels.

## 7. Conclusions

We have presented a new parameter identification problem for the Oseen equation, contributing to the detection of obstacles in fluid dynamic studies. For this problem, the existence of solution and optimality
conditions were determined. The study of optimality conditions is necessary to validate the use of optimization algorithms. In particular, the second order optimality condition allowed to specify the $H^{s}$ space where the parameter $\gamma$ belongs.

The numerical experiments developed in this work allow us to conjecture that it would be possible to extend this analysis for the Navier-Stokes equation. Indeed, the numerical tests without noise had satisfactory results in terms of rebuilding the simulated valve. Furthermore, experiments with Gaussian noise show that disturbances in the reference image are reflected in bounded changes in the parameter and in the reconstructed valve. The quality of the solutions is worse as noise increases, as expected.

There are several ways to deepen this work. From an analytical perspective, it is necessary to repeat the theoretical analysis for Problem (9), in the same way as we worked with Problem (1), looking for obtaining similar results. Furthermore, it is required to study the uniqueness of the solution and stability results for the control. Based on the favorable numerical results, the next step is to perform experiments with 3D domains, including structures and real images (both MRI and ultrasound imaging), whether obtained from phantoms or real patients. The 3D case is of medical interest, since it will contribute to simplify the detection of defects in the function of aortic valves, among other structures. Therefore, designing an algorithm to recreate the aortic valve in 3 D is also one of the future improvements.

## Acknowledgements

Jorge Aguayo was partially funded by the National Agency for Research and Development (ANID) / Scholarship Program / BECA DOCTORADO NACIONAL /2018-21180642. Cristóbal Bertoglio was funded by the European Research Council (ERC) under the European Union's Horizon 2020 research and innovation programme (grant agreement No 852544 - CardioZoom). Axel Osses was partially funded by Fondecyt Grant 1191903, Basal Program PFB-03 (AFB170001), MathAmsud 18-MATH-04, FONDAP/15110009 and Millennium Nucleus for Cardiovascular Magnetic Resonance.

## References

[1] Abraham, F., Behr, M., Heinkenschloss, M., 2004. The effect of stabilization in finite element methods for the optimal boundary control of the Oseen equations. Finite Elements in Analysis and Design 41, 229-251. doi 10.1016/j.finel. 2004.06.001
[2] Alnæs, M., Blechta, J., Hake, J., Johansson, A., Kehlet, B., Logg, A., Richardson, C., Ring, J., Rognes, M.E., Wells, G.N., 2015. The FEniCS Project Version 1.5. Archive of Numerical Software Vol 3. doi 10.11588/ANS.2015.100.20553.
[3] Alvarez, C., Conca, C., Friz, L., Kavian, O., Ortega, J., 2005. Identification of immersed obstacles via boundary measurements. Inverse Problems 21, 1531.
[4] Aßmann, U., Rösch, A., 2013. Identification of an unknown parameter function in the main part of an elliptic partial differential equation. Zeitschrift für Analysis und ihre Anwendungen 32. doi 10.4171/ZAA/1479
[5] Astorino, M., Hamers, J., Shadden, S.C., Gerbeau, J.F., 2012. A robust and efficient valve model based on resistive immersed surfaces. International Journal for Numerical Methods in Biomedical Engineering 28, 937-959. doi 10.1002/ cnm. 2474
[6] Auriault, J.L., 2008. On the domain of validity of Brinkman's Equation. Transport in Porous Media 79, $215-223$. doi $10.1007 / \mathrm{s} 11242-008-9308-7$.
[7] Ballerini, A., 2010. Stable determination of an immersed body in a stationary Stokes fluid. Inverse Problems $26,125015$.
[8] Bertoglio, C., Caiazzo, A., Bazilevs, Y., Braack, M., Esmaily, M., Gravemeier, V., Marsden, A.L., Pironneau, O., VignonClementel, I.E., Wall, W.A., 2017a. Benchmark problems for numerical treatment of backflow at open boundaries. International Journal for Numerical Methods in Biomedical Engineering 34, e2918. doi 10.1002/cnm. 2918
[9] Bertoglio, C., Nuñez, R., Galarce, F., Nordsletten, D., Osses, A., 2017b. Relative pressure estimation from velocity measurements in blood flows: State-of-the-art and new approaches. International Journal for Numerical Methods in Biomedical Engineering 34, e2925. doi $10.1002 / \mathrm{cnm} .2925$
[10] Bjærum, S., Torp, H., Kristoffersen, K., 2002. Clutter filter design for ultrasound color flow imaging. IEEE Transactions on Ultrasonics, Ferroelectrics and Frequency Control 49, 204-216. doi 10.1109/58.985705.
[11] Braack, M., Mucha, P.B., 2014. Directional Do-Nothing Condition for the Navier-Stokes Equations. Journal of Computational Mathematics 32, 507-521. doi $10.4208 / \mathrm{jcm} .1405-\mathrm{m} 4347$.
[12] Brinkman, H.C., 1949. A calculation of the viscous force exerted by a flowing fluid on a dense swarm of particles. Flow, Turbulence and Combustion 1. doi 10.1007/bf02120313.
[13] Caiazzo, A., Fernández, M.A., Gerbeau, J.F., Martin, V., 2011. Projection schemes for fluid flows through a porous interface. SIAM Journal on Scientific Computing 33, 541-564. doi 10.1137/100788124.
[14] Caubet, F., 2013. Instability of an Inverse Problem for the Stationary Navier-Stokes Equations. SIAM Journal on Control and Optimization 51, 2949-2975.
[15] Caubet, F., Dambrine, M., Kateb, D., Timimoun, C., 2013. A Kohn-Vogelius formulation to detect an obstacle immersed in a fluid. Inverse Problems \& Imaging 7, 123. doi 10.3934/ipi.2013.7.123
[16] Ciarlet, P.G., 2013. Linear and Nonlinear Functional Analysis with Applications. Society for Industrial and Applied Mathematics, USA.
[17] Dyverfeldt, P., Bissell, M., Barker, A.J., Bolger, A.F., Carlhääll, C.J., Ebbers, T., Francios, C.J., Frydrychowicz, A., Geiger, J., Giese, D., Hope, M.D., Kilner, P.J., Kozerke, S., Myerson, S., Neubauer, S., Wieben, O., Markl, M., 2015. 4D flow cardiovascular magnetic resonance consensus statement. Journal of Cardiovascular Magnetic Resonance 17. doi 10.1186/s12968-015-0174-5
[18] Fattouch, K., Lancellotti, P., Vannan, M., Speziale, G., 2018. Advances in Treatments for Aortic Valve and Root Diseases. Springer. Doi.org/10.1007/978-3-319-66483-5.
[19] Fedele, M., Faggiano, E., Dedè, L., Quarteroni, A., 2017. A patient-specific aortic valve model based on moving resistive immersed implicit surfaces. Biomechanics and Modeling in Mechanobiology 16, 1779-1803. doi 10.1007/s10237-017-0919-1
[20] Gröger, K., 1989. A $W^{1, p}$-Estimate for solutions to mixed boundary value problems for second order elliptic differential equtions. Mathematische Annalen 283, 679-688. URL: http://eudml.org/doc/164533
[21] Iung, B., Vahanian, A., 2011. Epidemiology of valvular heart disease in the adult. Nature Reviews Cardiology 8, $162-172$.
[22] John, V., 2016. Finite Element Methods for Incompressible Flow Problems. Springer International Publishing. doi 10. 1007/978-3-319-45750-5
[23] Leow, C.H., Bazigou, E., Eckersley, R.J., Yu, A.C.H., Weinberg, P.D., Tang, M.X., 2015. Flow Velocity Mapping Using Contrast Enhanced High-Frame-Rate Plane Wave Ultrasound and Image Tracking: Methods and Initial in-Vitro and in-Vivo Evaluation. Ultrasound in Medicine \& Biology 41, 2913-2925. doi 10.1016/j.ultrasmedbio.2015.06.012
[24] Li, Y.L., Hyun, D., Abou-Elkacem, L., Willmann, J.K., Dahl, J.J., 2016. Visualization of Small-Diameter Vessels by Reduction of Incoherent Reverberation With Coherent Flow Power Doppler. IEEE Transactions on Ultrasonics, Ferroelectrics, and Frequency Control 63, 1878-1889. doi 10.1109/tuffc.2016.2616112
[25] Mayo, J.R., Leipsic, J.A., 2009. Radiation Dose in Cardiac CT. American Journal of Roentgenology 192, 646-653. doi 10.2214/ajr.08. 2066
[26] Mitusch, S., Funke, S., Dokken, J., 2019. dolfin-adjoint 2018.1: automated adjoints for FEniCS and Firedrake. Journal of Open Source Software 4, 1292. doi $10.21105 /$ joss. 01292
[27] Nolte, D., Bertoglio, C., 2019. Reducing the impact of geometric errors in flow computations using velocity measurements. International Journal for Numerical Methods in Biomedical Engineering, e3203doi 10.1002/cnm. 3203
[28] de los Reyes, J.C., 2015. Numerical PDE-Constrained Optimization. Springer International Publishing. doi 10.1007/ 978-3-319-13395-9

## Appendix A. Proof of Theorem 7

Theorem. $A$ is Fréchet-differentiable. For each $\gamma, \gamma_{1} \in \mathcal{A}$, if $A(\gamma)=[\boldsymbol{v}, p]$, then the Fréchet derivative $A^{\prime}(\gamma)\left[\gamma_{1}\right]=\left[\boldsymbol{v}_{1}^{\prime}, p_{1}^{\prime}\right]$ can be described by the weak solution of the following problem

$$
\begin{aligned}
-\nu \triangle \boldsymbol{v}_{1}^{\prime}+\left(\nabla \boldsymbol{v}_{1}^{\prime}\right) \boldsymbol{a}+\nabla p_{1}^{\prime}+\gamma \boldsymbol{v}_{1}^{\prime} & =-\gamma_{1}(\boldsymbol{v}+\boldsymbol{G}) & & \text { in } \Omega \\
\operatorname{div} \boldsymbol{v}_{1}^{\prime} & =0 & & \text { in } \Omega \\
\boldsymbol{v}_{1}^{\prime} & =\mathbf{0} & & \text { on } \Gamma_{D} \\
-\nu \frac{\partial \boldsymbol{v}_{1}^{\prime}}{\partial \boldsymbol{n}}+p_{1}^{\prime} \boldsymbol{n} & =\mathbf{0} & & \text { on } \Gamma_{N} .
\end{aligned}
$$

Proof. Let $\gamma, \gamma_{1} \in \mathcal{A}$, it will be proved that there is a linear application $D: L^{\infty}(\Omega) \rightarrow \mathcal{H}$ such that

$$
A\left(\gamma+\gamma_{1}\right)-A(\gamma)=D\left(\gamma_{1}\right)+r\left(\gamma, \gamma_{1}\right),
$$

where

$$
\left\|\gamma_{1}\right\|_{\infty, \Omega} \rightarrow 0 \Rightarrow \frac{\left\|r\left(\gamma, \gamma_{1}\right)\right\|_{\mathcal{H}}}{\left\|\gamma_{1}\right\|_{\infty, \Omega}} \rightarrow 0
$$

Let $\left[\boldsymbol{v}_{1}, p_{1}\right]=A\left(\gamma+\gamma_{1}\right)$. Taking $D\left(\gamma_{1}\right)=\left[\boldsymbol{v}_{1}^{\prime}, p_{1}^{\prime}\right]$ and $r\left(\gamma, \gamma_{1}\right)=[\delta \boldsymbol{v}, \delta p]=\left[\boldsymbol{v}_{1}-\boldsymbol{v}-\boldsymbol{v}_{1}^{\prime}, p_{1}-p-p_{1}^{\prime}\right]$, it is possible to see that $D\left(\gamma_{1}\right)$ is linear. Also, $r\left(\gamma, \gamma_{1}\right)$ is solution of the problem

$$
\begin{aligned}
-\nu \triangle(\delta \boldsymbol{v})+\left(\nabla \delta \boldsymbol{v}_{1}\right) \boldsymbol{a}+\nabla(\delta p)+\gamma(\delta \boldsymbol{v}) & =\gamma_{1}\left(v-v_{1}\right) & & \text { in } \Omega \\
\operatorname{div}(\delta \boldsymbol{v}) & =0 & & \text { in } \Omega \\
(\delta \boldsymbol{v}) & =\mathbf{0} & & \text { on } \Gamma_{D} \\
-\nu \frac{\partial(\delta \boldsymbol{v})}{\partial \boldsymbol{n}}+(\delta p) \boldsymbol{n} & =\mathbf{0} & & \text { on } \Gamma_{N} .
\end{aligned}
$$

Thus, there exists a constant $c_{1}>0$, independent of $\gamma$ and $\gamma_{1}$ (see Lemma 5.8 on [22]), such that

$$
\left\|r\left(\gamma, \gamma_{1}\right)\right\|_{\mathcal{H}}=\nu|\delta \boldsymbol{v}|_{1, \Omega}+\|\delta p\|_{0, \Omega} \leq c_{1}\left\|\gamma_{1}\right\|_{0, \Omega}\left\|\boldsymbol{v}-\boldsymbol{v}_{1}\right\|_{0, \Omega} .
$$

subtracting the equations of $\boldsymbol{v}$ y $\boldsymbol{v}_{1}$ and using integration by parts,

$$
\nu\left|\boldsymbol{v}-\boldsymbol{v}_{1}\right|_{1, \Omega}^{2}+\left(\gamma\left(\boldsymbol{v}-\boldsymbol{v}_{1}\right), \boldsymbol{v}-\boldsymbol{v}_{1}\right)_{0, \Omega}=\left(\gamma_{1}\left(\boldsymbol{G}+\boldsymbol{v}_{1}\right), \boldsymbol{v}-\boldsymbol{v}_{1}\right)_{0, \Omega} .
$$

Applying Cauchy-Schwarz and Friedrichs-Poincaré inequalities, there exists $c_{2}>0$ such that

$$
\begin{aligned}
\nu\left|\boldsymbol{v}-\boldsymbol{v}_{1}\right|_{1, \Omega}^{2} & \leq\left(\gamma_{1}\left(\boldsymbol{G}+\boldsymbol{v}_{1}\right), \boldsymbol{v}-\boldsymbol{v}_{1}\right)_{0, \Omega} \\
& \leq\left\|\gamma_{1}\right\|_{\infty, \Omega}\left\|\boldsymbol{G}+\boldsymbol{v}_{1}\right\|_{0, \Omega}\left\|\boldsymbol{v}-\boldsymbol{v}_{1}\right\|_{0, \Omega} \\
& \leq c_{2}\left\|\gamma_{1}\right\|_{\infty, \Omega}\left\|\boldsymbol{G}+\boldsymbol{v}_{1}\right\|_{0, \Omega}\left|\boldsymbol{v}-\boldsymbol{v}_{1}\right|_{1, \Omega} .
\end{aligned}
$$

But, using Friedrichs-Poincaré inequality again, there exists $c_{3}>0$ such that $\left\|\boldsymbol{G}+\boldsymbol{v}_{1}\right\|_{0, \Omega} \leq c_{3}\|\boldsymbol{G}\|_{1, \Omega}$ (see Remark 5). Then, there exists $C_{1}>0$ such that

$$
\nu\left|\boldsymbol{v}-\boldsymbol{v}_{1}\right|_{1, \Omega} \leq C_{1}\left\|\gamma_{1}\right\|_{\infty, \Omega}\|\boldsymbol{G}\|_{0, \Omega}
$$

Finally, applying Friedrichs-Poincaré inequality, there exists a constant $C>0$ such that

$$
\left\|r\left(\gamma, \gamma_{1}\right)\right\|_{\mathcal{H}}=\nu|\delta \boldsymbol{v}|_{1, \Omega}+\|\delta p\|_{0, \Omega} \leq C\left\|\gamma_{1}\right\|_{\infty, \Omega}^{2}\|\boldsymbol{G}\|_{0, \Omega}
$$

Thus,

$$
\frac{\left\|r\left(\gamma, \gamma_{1}\right)\right\|_{\mathcal{H}}}{\left\|\gamma_{1}\right\|_{\infty, \Omega}} \leq C\left\|\gamma_{1}\right\|_{0, \Omega}\|\boldsymbol{G}\|_{0, \Omega} \xrightarrow{\left\|\gamma_{1}\right\|_{\infty, \Omega} \rightarrow 0} 0,
$$

proving the theorem.

## Appendix B. Proof of Theorem 14

Theorem. The map $\gamma \longmapsto A^{\prime}(\gamma)\left[\gamma_{1}\right]$ from $L^{\infty}(\Omega)$ to $\mathcal{H}$ for each $\gamma_{1} \in \mathcal{A}$ is Fréchet-differentiable. Let $\gamma, \gamma_{1}, \gamma_{2} \in \in \mathcal{A}$, the Fréchet derivative of $A^{\prime}(\gamma)\left[\gamma_{1}\right]$ on $\gamma_{2}$ direction is given by $A^{\prime \prime}(\gamma)\left[\gamma_{1}, \gamma_{2}\right]=\left[\boldsymbol{v}^{\prime \prime}, p^{\prime \prime}\right]$, where $\left[\boldsymbol{v}^{\prime \prime}, p^{\prime \prime}\right]$ is the unique weak solution of the problem

$$
\begin{align*}
-\nu \triangle \boldsymbol{v}^{\prime \prime}+\left(\nabla \boldsymbol{v}^{\prime \prime}\right) \boldsymbol{a}+\nabla p^{\prime \prime}+\gamma \boldsymbol{v}^{\prime \prime} & =-\left(\gamma_{2} \boldsymbol{v}_{1}^{\prime}+\gamma_{1} \boldsymbol{v}_{2}^{\prime}\right) & & \text { in } \Omega  \tag{B.1}\\
\operatorname{div} \boldsymbol{v}^{\prime \prime} & =0 & & \text { in } \Omega \\
\boldsymbol{v}^{\prime \prime} & =\mathbf{0} & & \text { on } \Gamma_{D} \\
-\nu \frac{\partial \boldsymbol{v}^{\prime \prime}}{\partial \boldsymbol{n}}+p^{\prime \prime} \boldsymbol{n} & =\mathbf{0} & & \text { on } \Gamma_{N} .
\end{align*}
$$

and $A^{\prime}(\gamma) \gamma_{j}=\left[\boldsymbol{v}_{j}^{\prime}, p_{j}^{\prime}\right]$, for $j \in\{1,2\}$.
Proof. Let $\gamma, \gamma_{1} \in \mathcal{A}$, it will be proved that there is a linear application $D_{2}: L^{\infty}(\Omega) \rightarrow \mathcal{H}$ such that

$$
A^{\prime}\left(\gamma+\gamma_{2}\right)\left[\gamma_{1}\right]-A^{\prime}(\gamma)\left[\gamma_{1}\right]=D_{2}\left(\gamma, \gamma_{1}\right)\left[\gamma_{2}\right]+r\left(\gamma, \gamma_{1}, \gamma_{2}\right)
$$

where

$$
\left\|\gamma_{2}\right\|_{\infty, \Omega} \rightarrow 0 \Rightarrow \frac{\left\|r\left(\gamma, \gamma_{1}, \gamma_{2}\right)\right\|_{\mathcal{H}}}{\left\|\gamma_{2}\right\|_{\infty, \Omega}} \rightarrow 0
$$

Let $A^{\prime}\left(\gamma+\gamma_{2}\right)\left[\gamma_{1}\right]=\left[\boldsymbol{w}_{1}^{\prime}, q_{1}^{\prime}\right], A^{\prime}(\gamma)\left[\gamma_{1}\right]=\left[\boldsymbol{v}_{1}^{\prime}, p_{1}^{\prime}\right]$ (see Theorem 7), $A(\gamma)=[\boldsymbol{v}, p]$ and $A\left(\gamma+\gamma_{2}\right)=[\boldsymbol{w}, q]$ (see Definition 4). Defining $D\left(\gamma, \gamma_{1}\right)\left[\gamma_{2}\right]=\left[\boldsymbol{v}^{\prime \prime}, p^{\prime \prime}\right]$ as function of $\gamma_{2}$ and

$$
r\left(\gamma, \gamma_{1}, \gamma_{2}\right)=[\delta \boldsymbol{v}, \delta p]=\left[\boldsymbol{w}_{1}^{\prime}-\boldsymbol{v}_{1}^{\prime}-\boldsymbol{v}^{\prime \prime}, q_{1}^{\prime}-p_{1}^{\prime}-p^{\prime \prime}\right]
$$

it is possble to see that $D\left(\gamma, \gamma_{1}\right)$ is linear and $r\left(\gamma, \gamma_{1}, \gamma_{2}\right)$ is solution of the problem

$$
\begin{aligned}
-\nu \triangle(\delta \boldsymbol{v})+(\nabla \delta \boldsymbol{v}) \boldsymbol{a}+\nabla(\delta p)+\gamma(\delta \boldsymbol{v}) & =\gamma_{2}\left(\boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right)+\gamma_{1}\left(\boldsymbol{v}_{2}^{\prime}+\boldsymbol{v}-\boldsymbol{w}\right) & & \text { in } \Omega \\
\operatorname{div}(\delta \boldsymbol{v}) & =0 & & \text { in } \Omega \\
(\delta \boldsymbol{v}) & =\mathbf{0} & & \text { on } \Gamma_{D} \\
-\nu \frac{\partial(\delta \boldsymbol{v})}{\partial \boldsymbol{n}}+(\delta p) \boldsymbol{n} & =\mathbf{0} & & \text { on } \Gamma_{N} .
\end{aligned}
$$

Then, applying Friedrichs-Poincaré inequality, there exists $c_{1}>0$, independent of $\gamma, \gamma_{1}, \gamma_{2}$ such that

$$
\nu|\delta \boldsymbol{v}|_{1, \Omega}+\|\delta p\|_{0, \Omega} \leq c_{1}\left(\left\|\gamma_{2}\right\|_{\infty, \Omega}\left|\boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right|_{1, \Omega}+\left\|\gamma_{1}\right\|_{\infty, \Omega}\left|\boldsymbol{v}_{2}^{\prime}+\boldsymbol{v}-\boldsymbol{w}\right|_{1, \Omega}\right) .
$$

Later, subtracting the equations of $\boldsymbol{v}_{1}^{\prime}$ and $\boldsymbol{w}_{1}^{\prime}$, testing with $\boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}$ and $p_{1}^{\prime}-q_{1}$, respectively, and applying Friedrichs-Poincaré inequality, there exists $c_{2}>0$ such that

$$
\begin{aligned}
\nu\left|\boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right|_{1, \Omega}^{2} & \leq\left(\gamma_{2} \boldsymbol{w}_{1}^{\prime}, \boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right)_{0, \Omega}+\left(\gamma_{1}(\boldsymbol{w}-\boldsymbol{v}), \boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right)_{0, \Omega} \\
& \leq\left\|\gamma_{2}\right\|_{\infty, \Omega}\left\|\boldsymbol{w}_{1}^{\prime}\right\|_{0, \Omega}\left\|\boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right\|_{0, \Omega}+\left\|\gamma_{1}\right\|_{\infty, \Omega}\|\boldsymbol{w}-\boldsymbol{v}\|_{0, \Omega}\left\|\boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right\|_{0, \Omega} \\
& \leq c_{2}\left(\left\|\gamma_{2}\right\|_{\infty, \Omega}\left|\boldsymbol{w}_{1}^{\prime}\right|_{1, \Omega}+\left\|\gamma_{1}\right\|_{\infty, \Omega}|\boldsymbol{w}-\boldsymbol{v}|_{1, \Omega}\right)\left|\boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right|_{1, \Omega} .
\end{aligned}
$$

Hence,

$$
\nu\left|\boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right|_{1, \Omega} \leq c_{2}\left(\left\|\gamma_{2}\right\|_{\infty, \Omega}\left|\boldsymbol{w}_{1}^{\prime}\right|_{1, \Omega}+\left\|\gamma_{1}\right\|_{\infty, \Omega}|\boldsymbol{w}-\boldsymbol{v}|_{1, \Omega}\right)
$$

Analogously, from the equations of $\boldsymbol{w}$ and $\boldsymbol{v}$, there exists $c_{2}>0$ such that

$$
\begin{aligned}
\nu|\boldsymbol{w}-\boldsymbol{v}|_{1, \Omega}^{2} & \leq-\left(\gamma_{2}(\boldsymbol{v}+\boldsymbol{G}), \boldsymbol{w}-\boldsymbol{v}\right)_{0, \Omega} \\
& \leq\left\|\gamma_{2}\right\|_{\infty, \Omega}\|\boldsymbol{v}+\boldsymbol{G}\|_{0, \Omega}\|\boldsymbol{w}-\boldsymbol{v}\|_{0, \Omega} \\
& \leq c_{2}\left\|\gamma_{2}\right\|_{\infty, \Omega}|\boldsymbol{v}+\boldsymbol{G}|_{1, \Omega}|\boldsymbol{w}-\boldsymbol{v}|_{1, \Omega}
\end{aligned}
$$

but there exists $c_{3}>0$ such that $|\boldsymbol{v}|_{1, \Omega} \leq c_{3}\|\boldsymbol{G}\|_{1, \Omega}$ (see Remark 5). Then,

$$
\nu|\boldsymbol{w}-\boldsymbol{v}|_{1, \Omega} \leq c_{2} c_{3}\left\|\gamma_{2}\right\|_{\infty, \Omega}\|\boldsymbol{G}\|_{1, \Omega}
$$

Futhermore, there exists $c_{4}>0$ such that

$$
\nu\left|\boldsymbol{w}_{1}^{\prime}\right|_{1, \Omega} \leq c_{4}\|\boldsymbol{G}\|_{1, \Omega}
$$

In conclusion, there exists $C_{1}>0$ such that

$$
\nu\left|\boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right|_{1, \Omega} \leq C_{1}\left\|\gamma_{1}\right\|_{\infty, \Omega}\left\|\gamma_{2}\right\|_{\infty, \Omega}\|\boldsymbol{G}\|_{1, \Omega}
$$

Repeating this analysis for the term $\boldsymbol{v}_{2}^{\prime}+\boldsymbol{v}-\boldsymbol{w}$, there exists $c_{5}>0$ such that

$$
\begin{aligned}
\nu\left|\boldsymbol{v}_{2}^{\prime}+\boldsymbol{v}-\boldsymbol{w}\right|_{0, \Omega} & \leq c_{5}\left\|\gamma_{2}\right\|_{\infty, \Omega}|\boldsymbol{w}-\boldsymbol{v}|_{1, \Omega} \\
& \leq c_{5}\left\|\gamma_{2}\right\|_{\infty, \Omega}^{2}\|\boldsymbol{G}\|_{1, \Omega} .
\end{aligned}
$$

Then, there exists $C_{2}>0$ such that

$$
\begin{aligned}
\left\|r\left(\gamma, \gamma_{1}, \gamma_{2}\right)\right\|_{\mathcal{H}} & =\nu|\delta \boldsymbol{v}|_{1, \Omega}+\|\delta p\|_{0, \Omega} \\
& \leq c_{1}\left(\left\|\gamma_{2}\right\|_{\infty, \Omega}\left|\boldsymbol{v}_{1}^{\prime}-\boldsymbol{w}_{1}^{\prime}\right|_{0, \Omega}+\left\|\gamma_{1}\right\|_{\infty, \Omega}\left|\boldsymbol{v}_{2}^{\prime}+\boldsymbol{v}-\boldsymbol{w}\right|_{1, \Omega}\right) \\
& \leq c_{1}\left(c_{2}\left\|\gamma_{1}\right\|_{\infty, \Omega}\left\|\gamma_{2}\right\|_{\infty, \Omega}^{2}\|\boldsymbol{G}\|_{1, \Omega}+c_{5} c_{2}\left\|\gamma_{1}\right\|_{\infty, \Omega}\left\|\gamma_{2}\right\|_{\infty, \Omega}^{2}\|\boldsymbol{G}\|_{1, \Omega}\right) \\
& \leq C_{2}\left\|\gamma_{1}\right\|_{\infty, \Omega}\left\|\gamma_{2}\right\|_{\infty, \Omega}^{2}\|\boldsymbol{G}\|_{1, \Omega} .
\end{aligned}
$$

In conclusion,

$$
\frac{\left\|r\left(\gamma, \gamma_{1}, \gamma_{2}\right)\right\|_{\mathcal{H}}}{\left\|\gamma_{2}\right\|_{\infty, \Omega}} \leq C_{2}\left\|\gamma_{1}\right\|_{\infty, \Omega}\left\|\gamma_{2}\right\|_{\infty, \Omega}\|\boldsymbol{G}\|_{0, \Omega} \xrightarrow{\left\|\gamma_{2}\right\|_{\infty, \Omega} \rightarrow 0} 0,
$$

proving the theorem


[^0]:    Email addresses: jaguayo@dim.uchile.cl (Jorge Aguayo), c.a.bertoglio@rug.nl (Cristóbal Bertoglio), axosses@dim.uchile.cl (Axel Osses)

