# A smooth extension method for transmission problems 

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# A smooth extension method for transmission problems 

Astrid Decoene, ${ }^{*}$ Sébastien Martin $\dagger$, Fabien Vergnet ${ }^{\ddagger}$

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#### Abstract

In this work, we present a numerical method for the resolution of transmission problems with non-conformal meshes which preserves the optimal rates of convergence in space. The Smooth Extension method is a fictitious domain approach based on a control formulation stated as a minimization problem, that we prove to be equivalent to the initial transmission problem. Formulated as a minimization problem, the transmission problem can be solved with standard finite element function spaces and usual optimization algorithms. The method is applied to different transmission problems (Laplace, Stokes and a fluid-structure interaction problem) and compared to standard finite element methods.


Keywords. transmission problem, interface problem, partial differential equations with discontinuous coefficients, fictitious domain, optimal rates of convergence, finite element method, smooth extension method

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## 1 Introduction

The numerical simulation of transmission problems (also known as interface problems or problems with discontinuous coefficients) is of major importance for the mathematical study of many physical and living systems. As examples, we can mention the study of composite materials (see [29]), the flow of multiphasic fluids or the swimming of bacteria in a viscous fluid (see [16]). There exists a wide variety of numerical methods for this type of problems, which can be categorized in two classes: fitted and unfitted mesh techniques.

Fitted mesh methods are successful when dealing with stationary problems or problems with domains ongoing moderate displacements. In the later, the mesh is moved using a regular extension of the displacement of the interface, which is only possible in practice if this displacement is small enough. For example, in the case of fluid-structure interaction problems it can be done using an arbitrary Lagrangian-Eulerian (ALE) description of the fluid (see [17]). The main advantage of fitted mesh techniques is their optimal convergence rate in space. However, when transmission problems involve large deformations, a much more appropriate method is the use of unfitted meshes, in which the interface deforms independently of a background fixed mesh.

Unfitted mesh techniques use a non-conformal mesh, i.e. a mesh that does not fit the interface. Among these methods we can reference the immersed boundary method (see [31]), fictitious domain methods with Lagrange multipliers (see [5], [33]) or penalization terms (see [24]), the extended-finite element method (XFEM, see [27], [21] and [20]), the Nitsche-XFEM method (see [2]) and the fat boundary method (FBM, see [7]). The first two techniques do not converge with optimal orders in space (see [22], [23], [32], [26]), because of the discrete treatment of transmission conditions. The XFEM and the Nitsche-XFEM overcome this issue, but they come with a cost: the XFEM is known to lack robustness; the Nitsche-XFEM method circumvents theses difficulties but, still, all cut-FEM based methods require a specific evaluation of the interface intersections, which can be difficult, especially in three space dimensions (see [9]). An interesting fictitious domain method is also presented in [25], where the optimal rates of convergence are obtained for the Poisson problem with Dirichlet boundary conditions on a mesh that does not fit the boundary of the domain. The originality of this work consists in avoiding integration over the elements cut by the boundary of the problem domain, while preserving the optimal rates of convergence. The FBM is a fictitious domain method designed to recover the optimal convergence at any order. It is well suited for elliptic stationary problems or problems involving rigid domains, but it is not straightforward to adapt it to deformable moving materials. An other class of unfitted mesh techniques is the control based approach presented in [3], initially developed to solve boundary value problems in complex geometries and which is based on an optimal control formulation. In [3], the authors consider the Poisson problem for the Laplace operator with an included obstacle on which Dirichlet and Neumann boundary conditions are applied. The idea of the method is to extend the problem inside the obstacle and use the right-hand side of the equations as a control to impose the boundary conditions on the frontier of the obstacle. This method has also been used to treat boundary conditions in a fictitious domain approach for the Helmoltz equations in [3], [4] and [30]. In [18] and [19], the authors present the smooth extension method, an extension of the control based method to the resolution of fluid-structure interaction problems involving rigid particles and a viscous fluid. This method has the advantages to recover optimal convergence in space
and to be rather simple to implement (it does not necessitate any mesh adaptation or local enrichment of the function spaces).

Here, we aim to generalize this smooth extension method to transmission problems, i.e. to a class of problems where the behavior of the obstacle is also described by partial differential equations with transmission conditions on the interface. Then, the smooth extension method presented in this chapter is a finite element method based on a control approach to solve elliptic transmission problems with unfitted meshes and which recovers the optimal convergence rates. This method is designed to be suitable for transmission problems involving several small materials included in a bigger material, such bubbles in a fluid, rigid and deformable particles in a fluid or problems of conductivity and elasticity of non-homogeneous materials. Because it is an unfitted mesh technique, it enables to precisely mesh the small materials without necessarily considering a thin mesh on the whole domain and thus, speed up the numerical resolutions of such problems. It also enables to use a Cartesian mesh on the bigger domain, which permits the utilization of fast solvers (e.g. FFT solvers).

The remaining of the chapter is organized as follows. In section 2, we state the smooth extension method for the Laplace transmission problem. section 3 is devoted to the presentation of numerical experiments using the methodology described in section 2 . In section 4 , we extend the method to other transmission problems: the Stokes transmission problem and a fluid-structure interaction problem involving a linear elastic structure and a viscous fluid.

## 2 Presentation of the method

The smooth extension method presented here is devoted to the numerical simulation of transmission problems. It is a fictitious domain method which enables to recover the optimal order of convergence in space by smoothly extending a part of the exact solution to the whole domain. In this section we focus on a toy model, the Laplace transmission problem, in order to properly explain the main steps of the method. This section is divided in four parts. In subsection 2.1, we state the smooth extension formulation of the problem as a control problem and explain how to recover the exact solution from it. Then, in subsection 2.2 we rewrite this control problem as a minimization problem and prove the equivalence of the two formulations. In subsection 2.3 we discuss the advantage of this method compared to the classical finite element method. Finally, in subsection 2.4, we extend the method to the particular case of strictly included domains.

### 2.1 The smooth extension formulation

Let $n>0$ and $\Omega$ be a domain of $\mathbb{R}^{n}$ that satisfies the following set of hypotheses:
i) Domain $\Omega$ is a bounded connected Lipschitz domain of $\mathbb{R}^{n}$.
ii) Domain $\Omega$ is divided in two subdomains, $\Omega_{1}$ and $\Omega_{2}$, which have Lipschitz boundaries.
iii) The interface $\Gamma=\partial \Omega_{1} \cap \partial \Omega_{2}$ is not empty.
iv) The remaining boundaries $\Gamma_{1}=\partial \Omega_{1} \backslash \Gamma$ and $\Gamma_{2}=\partial \Omega_{2} \backslash \Gamma$ are not empty.

The problem we are interested in, for now, is the coupled Laplace problem with homogeneous Dirichlet boundary conditions, also called Laplace transmission problem or Laplace problem with discontinuous coefficients. It is the simplest coupled system of partial differential equations that we can think of and, given two positive real constants $\mu_{1}$ and $\mu_{2}$
such that $\mu_{1} \neq \mu_{2}$ and two source terms $f_{1} \in L^{2}\left(\Omega_{1}\right)$ and $f_{2} \in L^{2}\left(\Omega_{2}\right)$, it writes:

$$
\left\{\begin{array}{rlrl}
\text { find } u_{1}: \Omega_{1} \rightarrow \mathbb{R} & \text { and } u_{2}: \Omega_{2} \rightarrow \mathbb{R} \text { such that }  \tag{1a}\\
-\mu_{1} \Delta u_{1} & =f_{1}, & & \text { in } \Omega_{1}, \\
u_{1} & =0, & & \text { on } \Gamma_{1}, \\
-\mu_{2} \Delta u_{2} & =f_{2}, & & \text { in } \Omega_{2}, \\
u_{2} & =0, & & \text { on } \Gamma_{2}, \\
u_{1} & =u_{2}, & & \text { on } \Gamma \\
\mu_{1} \nabla u_{1} \cdot n_{1} & =-\mu_{2} \nabla u_{2} \cdot n_{2}, & & \text { on } \Gamma .
\end{array}\right.
$$

This problem is completely equivalent to the more classical formulation of the Laplace problem, written in the whole domain $\Omega$,

$$
\left\{\begin{aligned}
\text { find } u: \Omega \rightarrow \mathbb{R} & \text { such that } \\
-\operatorname{div}(\mu \nabla u) & =f \text { in } \Omega \\
u & =0 \quad \text { on } \quad \partial \Omega,
\end{aligned}\right.
$$

where $\mu$ and $f$ are defined by

$$
\begin{array}{r}
\mu=\left\{\begin{array}{lll}
\mu_{1} & \text { in } & \Omega_{1} \\
\mu_{2} & \text { in } & \Omega_{2}
\end{array}\right. \\
f=\left\{\begin{array}{lll}
f_{1} & \text { in } & \Omega_{1} \\
f_{2} & \text { in } & \Omega_{2}
\end{array}\right.
\end{array}
$$

However, the formulation (1) has the advantage of making the coupling between the two subproblems (1a) and (1b) clear. The so-called coupling conditions are detailed in equations (1c) and physically represent the continuity of the field and the continuity of the normal constraint through the interface $\Gamma$. Moreover, $n_{1}$ (resp. $n_{2}$ ) is the unit exterior normal vector of $\Omega_{1}$ (resp. $\Omega_{2}$ ). The solution of problem (1) will be searched in $V_{1} \times V_{2}$, where these two functional spaces are defined by

$$
\begin{aligned}
V_{1} & =\left\{v_{1} \in H^{1}\left(\Omega_{1}\right) ; v_{\left.1\right|_{\Gamma_{1}}}=0\right\} \\
V_{2} & =\left\{v_{2} \in H^{1}\left(\Omega_{2}\right) ; v_{\left.2\right|_{\Gamma_{2}}}=0\right\} .
\end{aligned}
$$

In addition, the dual spaces of $V_{1}$ and $V_{2}$ will be denoted by $V_{1}^{\prime}$ and $V_{2}^{\prime}$.
Let $v$ be a distribution in $\mathcal{D}(\Omega)$ and suppose that $u_{1}$ and $u_{2}$ are sufficiently regular. We can formally multiply the first equation in (1a) by $v_{\mid \Omega_{1}}$ and the first equation in (1b) by $v_{\mid \Omega_{2}}$ and integrate respectively over $\Omega_{1}$ and $\Omega_{2}$. After an integration by part, we obtain

$$
\begin{array}{r}
\mu_{1} \int_{\Omega_{1}} \nabla u_{1} \cdot \nabla v_{\mid \Omega_{1}}-\int_{\Gamma}\left(\mu_{1} \nabla u_{1} \cdot n_{1}\right) v+\mu_{2} \int_{\Omega_{2}} \nabla u_{2} \cdot \nabla v_{\mid \Omega_{2}}-\int_{\Gamma}\left(\mu_{2} \nabla u_{2} \cdot n_{2}\right) v \\
=\int_{\Omega_{1}} f_{1} v_{\mid \Omega_{1}}+\int_{\Omega_{2}} f_{2} v_{\mid \Omega_{2}}, \quad \forall v \in \mathcal{D}(\Omega)
\end{array}
$$

Using the second transmission condition in (1c) i.e., that $\mu_{1} \nabla u_{1} \cdot n_{1}=-\mu_{2} \nabla u_{2} \cdot n_{2}$ on $\Gamma$, it follows that

$$
\mu_{1} \int_{\Omega_{1}} \nabla u_{1} \cdot \nabla v_{\mid \Omega_{1}}+\mu_{2} \int_{\Omega_{2}} \nabla u_{2} \cdot \nabla v_{\mid \Omega_{2}}=\int_{\Omega_{1}} f_{1} v_{\mid \Omega_{1}}+\int_{\Omega_{2}} f_{2} v_{\mid \Omega_{2}}, \quad \forall v \in \mathcal{D}(\Omega)
$$

Then, introducing the space

$$
\mathcal{V}=\left\{\left(v_{1}, v_{2}\right) \in V_{1} \times V_{2} ; v_{\left.1\right|_{\Gamma}}=v_{\left.2\right|_{\Gamma}}\right\}
$$

we can define the weak formulation of problem (1):

$$
\left\{\begin{array}{l}
\text { find }\left(u_{1}, u_{2}\right) \text { in } \mathcal{V} \text { such that, }  \tag{2}\\
\mu_{1} \int_{\Omega_{1}} \nabla u_{1} \cdot \nabla v_{1}+\mu_{2} \int_{\Omega_{2}} \nabla u_{2} \cdot \nabla v_{2}=\int_{\Omega_{1}} f_{1} v_{1}+\int_{\Omega_{2}} f_{2} v_{2}, \forall\left(v_{1}, v_{2}\right) \in \mathcal{V} .
\end{array}\right.
$$

Problem (2) is well-posed, since $f_{1} \in L^{2}\left(\Omega_{1}\right)$ and $f_{2} \in L^{2}\left(\Omega_{2}\right)$, according to the LaxMilgram theorem. Furthermore, the unique solution of problem (2) will be denoted by ( $\bar{u}_{1}, \bar{u}_{2}$ ).

At this point, we can precise in what sense the solution of (2) is also solution of the initial problem (1). Taking test functions $v_{1}$ in $\mathcal{D}\left(\Omega_{1}\right)$ and $v_{2}$ in $\mathcal{D}\left(\Omega_{2}\right)$ in (2) gives that

$$
-\mu_{1} \Delta \bar{u}_{1}=f_{1} \text { and }-\mu_{2} \Delta \bar{u}_{2}=f_{2} \text { in } L^{2} .
$$

It also implies that $\nabla \bar{u}_{1}$ belongs to $H_{\text {div }}\left(\Omega_{1}\right)$ and that $\nabla \bar{u}_{2}$ belongs to $H_{\text {div }}\left(\Omega_{2}\right)$, where $H_{\text {div }}\left(\Omega_{1}\right)$ and $H_{\text {div }}\left(\Omega_{2}\right)$ are the spaces defined by

$$
H_{\operatorname{div}}(X)=\left\{\sigma \in\left(L^{2}(X)\right)^{n} ; \operatorname{div}(v) \in L^{2}(X)\right\},
$$

and where $X$ stands for either $\Omega_{1}$ or $\Omega_{2}$. In particular, we are able to give a weak sense to the normal derivatives of $\bar{u}_{1}$ and $\bar{u}_{2}$. Let $X$ stand for either $\Omega_{1}$ or $\Omega_{2}$ and let $\eta$ be the unit exterior normal vector of $X$. Let $\Lambda=H_{00}^{1 / 2}(\Gamma)$ be the image of $H_{\partial X \backslash \Gamma}^{1}(X)$ by the trace operator on the interface $\Gamma$, i.e. the space of functions in $H^{1 / 2}(\Gamma)$ whose extension by zero on $\partial X \backslash \Gamma$ belongs to $H^{1 / 2}(\partial X)$. Then, For all $\sigma$ in $H_{\text {div }}(X)$, we have the following Stokes formula:

$$
\begin{equation*}
\int_{X} \sigma \cdot \nabla v+\int_{X} \operatorname{div}(\sigma) v=\left\langle\gamma_{\eta}(\sigma), v\right\rangle_{\Lambda^{\prime}, \Lambda}, \quad \forall v \in H_{\partial X \backslash \Gamma}^{1}(X) . \tag{3}
\end{equation*}
$$

where $\Lambda^{\prime}=\left(H_{00}^{1 / 2}(\Gamma)\right)^{\prime}$ is the dual space of $\Lambda$ and $\gamma_{\eta}$ denotes the normal trace operator on $\Gamma$. Then, $\bar{u}_{1}$ verifies $-\mu_{1} \Delta \bar{u}_{1}=f_{1}$ a.e. in $\Omega_{1}, \bar{u}_{2}$ verifies $-\mu_{2} \Delta \bar{u}_{2}=f_{2}$ a.e. in $\Omega_{2}$, the first transmission condition, $\bar{u}_{1}=\bar{u}_{2}$ on $\Gamma$, is included in the functional space $\mathcal{V}$ and the second one is verified in a weak sense:

$$
\begin{equation*}
\left\langle\mu_{1} \gamma_{n_{1}}\left(\nabla \bar{u}_{1}\right), v\right\rangle_{\Lambda^{\prime}, \Lambda}=-\left\langle\mu_{2} \gamma_{n_{2}}\left(\nabla \bar{u}_{2}\right), v\right\rangle_{\Lambda^{\prime}, \Lambda}, \quad \forall v \in \Lambda . \tag{4}
\end{equation*}
$$

Now, we present the smooth extension method applied to problem (1). This method consists in $i$ ) extending the problem in $\Omega_{1}$ into a problem defined in the whole domain $\Omega$, ii) relaxing the condition of equality in the functional space $\mathcal{V}$ and $i i i$ ) finding a control term in $\Omega_{2}$ to enforce the condition of equality on $\Gamma$. We denote by $g$ this control term, which should belong to $V_{2}^{\prime}$, and we define $\bar{g}^{\Omega}$, the extension of $g$ in the whole domain $\Omega$ such that

$$
\left\langle\bar{g}^{\Omega}, v\right\rangle_{H^{-1}(\Omega), H_{0}^{1}(\Omega)}:=\left\langle g, v_{\left.\right|_{\Omega_{2}}}\right\rangle_{V_{2}^{\prime}, V_{2}}, \forall v \in H_{0}^{1}(\Omega) .
$$

In a similar way, as $f_{1}$ belongs to $L^{2}\left(\Omega_{1}\right), \overline{f_{1}}$ denotes the extension of $f_{1}$ by 0 over the whole space $\Omega$.

Formally, we define the smooth extension problem associated to problem (1) as the problem of finding a suitable control $g$ in $V_{2}^{\prime}$, such that the solution of the following problem,

$$
\left\{\begin{array}{rlrlr}
\text { find } u_{1}: \Omega \rightarrow \mathbb{R} & \text { and } u_{2}: \Omega_{2} \rightarrow \mathbb{R} \text { such that }  \tag{5a}\\
-\mu_{1} \Delta u_{1} & =\bar{f}_{1}+\bar{g}^{\Omega}, & & \text { in } \Omega, \\
u_{1} & =0, & & \text { on } \partial \Omega, \\
-\mu_{2} \Delta u_{2} & =f_{2}, & & \text { in } \Omega_{2}, \\
u_{2} & =0, & & \text { on } \Gamma_{2}, \\
\mu_{2} \nabla u_{2} \cdot n_{2} & =-\mu_{1} \nabla u_{1} \cdot n_{1}, & & \text { on } \Gamma,
\end{array}\right.
$$

satisfies the equality

$$
\left(u_{\left.1\right|_{\Omega_{1}}}, u_{2}\right)=\left(\bar{u}_{1}, \bar{u}_{2}\right)
$$

where $\left(\bar{u}_{1}, \bar{u}_{2}\right)$ is the solution of the initial problem (1). Again, let $v$ be in $\mathcal{D}(\Omega)$ and suppose that $u_{1}$ and $u_{2}$ are sufficiently regular. Moreover, we assume for the moment that $g$ belongs to $L^{2}\left(\Omega_{2}\right)$ in order to do formal computations. We multiply equation (5a) by $v$ and equation (5b) by $v_{\mid \Omega_{2}}$ and integrate respectively over $\Omega$ and $\Omega_{2}$. After an integration by part and using the Neumann condition on $\Gamma$, we find

$$
\begin{aligned}
\mu_{1} \int_{\Omega} \nabla u_{1} \cdot \nabla v & =\int_{\Omega_{1}} f_{1} v_{\mid \Omega_{1}}+\int_{\Omega_{2}} g v_{\mid \Omega_{2}}, \\
\mu_{2} \int_{\Omega_{2}} \nabla u_{2} \cdot \nabla v_{\mid \Omega_{2}} & =\int_{\Omega_{2}} f_{2} v_{\mid \Omega_{2}}-\int_{\Gamma}\left(\mu_{1} \nabla u_{1} \cdot n_{1}\right) v
\end{aligned}
$$

Furthermore, we remark that

$$
\begin{aligned}
-\int_{\Gamma}\left(\mu_{1} \nabla u_{1} \cdot n_{1}\right) v & =\int_{\Gamma}\left(\mu_{1} \nabla u_{1} \cdot n_{2}\right) v, \\
& =\mu_{1} \int_{\Omega_{2}} \Delta u_{1 \mid \Omega_{2}} v_{\mid \Omega_{2}}+\mu_{1} \int_{\Omega_{2}} \nabla u_{1 \mid \Omega_{2}} \cdot \nabla v_{\mid \Omega_{2}}, \\
& =-\int_{\Omega_{2}} g v_{\mid \Omega_{2}}+\mu_{1} \int_{\Omega_{2}} \nabla u_{1 \mid \Omega_{2}} \cdot \nabla v_{\mid \Omega_{2}} .
\end{aligned}
$$

Thus, we define the weak formulation of the smooth extension problem associated to problem (2), that makes sense for $g$ in $V_{2}^{\prime}$, as the problem of finding a suitable control $g$ in $V_{2}^{\prime}$, such that the solution of the following problem,

$$
\left\{\begin{array}{rlr}
\text { find }\left(u_{1}, u_{2}\right) \text { in } H_{0}^{1}(\Omega) \times V_{2} \text { such that }  \tag{6}\\
\mu_{1} \int_{\Omega} \nabla u_{1} \cdot \nabla v_{1}= & \int_{\Omega_{1}} f_{1} v_{1 \mid \Omega_{1}}+\left\langle g, v_{1 \mid \Omega_{2}}\right\rangle_{V_{2}^{\prime}, V_{2}^{\prime}}, & \forall v_{1} \in H_{0}^{1}(\Omega), \\
\mu_{2} \int_{\Omega_{2}} \nabla u_{2} \cdot \nabla v_{2}= & \int_{\Omega_{2}} f_{2}, v_{2}-\left\langle g, v_{2}\right\rangle_{V_{2}^{\prime}, V_{2}} & \\
& +\mu_{1} \int_{\Omega_{2}} \nabla u_{1 \mid \Omega_{2}} \cdot \nabla v_{2}, & \forall v_{2} \in V_{2} .
\end{array}\right.
$$

satisfies the equality

$$
\begin{equation*}
\left(u_{\left.1\right|_{\Omega_{1}}}, u_{2}\right)=\left(\bar{u}_{1}, \bar{u}_{2}\right) \tag{7}
\end{equation*}
$$

where ( $\bar{u}_{1}, \bar{u}_{2}$ ) is the solution of the initial weak problem (2). By standard arguments and considering the regularity of $f_{1}$ and $f_{2}$, it is straightforward to prove that problem (6) is well-posed for every $g$ in $V_{2}^{\prime}$, using the Lax-Milgram theorem. Consequently, we denote by $\left(u_{1}^{g}, u_{2}^{g}\right)$ its unique solution. Problem (6) is weakly coupled, in the sense that, given $g$ in $V_{2}^{\prime}, u_{1}^{g}$ is obtained independently of $u_{2}^{g}$ and $u_{2}^{g}$ is obtained afterwards from $g$ and $u_{1}^{g}$. We remark that finding such a control $g$ enables to directly obtain the solution of the initial coupled problem (2). Then, in what follows, we will prove the existence of at least one control $g$ and detail the process of finding it.

With these notations we state the following theorem, showing the existence of a control $g$ for which the couple $\left(u_{1}^{g}, u_{2}^{g}\right)$, the unique solution of (6), verifies the condition (7).
Theorem 2.1. Let $\Omega$ be a domain that satisfies Assumption $\left(H_{1}\right)$. Consider $f_{1}$ in $L^{2}\left(\Omega_{1}\right), f_{2}$ $i L^{2}\left(\Omega_{2}\right)$ and let $\left(\bar{u}_{1}, \bar{u}_{2}\right)$ be the unique solution of problem (2). Then, there exists a function $g$ in $V_{2}^{\prime}$ such that the solution $\left(u_{1}^{g}, u_{2}^{g}\right)$ of problem (6) verifies (7).

Proof. We can construct an extension operator $E$ that extend $\bar{u}_{1}$ into the whole space $H_{0}^{1}(\Omega)$. Indeed, consider the operator defined by

$$
E \bar{u}_{1}=\left\{\begin{array}{lll}
\bar{u}_{1} & \text { in } & \Omega_{1} \\
\bar{u}_{2} & \text { in } & \Omega_{2}
\end{array} .\right.
$$

Since $\bar{u}_{1}$ belongs to $V_{1}, \bar{u}_{2}$ belongs to $V_{2}$ and $\bar{u}_{1}=\bar{u}_{2}$ on the interface $\Gamma$, then $E \bar{u}_{1}$ is an extension of $\bar{u}_{1}$ which belongs to $H_{0}^{1}(\Omega)$. Furthermore, $\Delta\left(E \bar{u}_{1}\right)$ belongs to $L^{2}(\Omega)$.

Then, we construct $g$ in $V_{2}^{\prime}$ such that

$$
\begin{array}{rlr}
\langle g, v\rangle_{V_{2}^{\prime}, V_{2}}= & \mu_{1} \int_{\Omega_{2}} \nabla\left(E \bar{u}_{1}\right)_{\mid \Omega_{2}} \cdot \nabla v+\left\langle\mu_{1} \gamma_{n_{1}}\left(\nabla E \bar{u}_{1}\right), v\right\rangle_{\Lambda^{\prime}, \Lambda}, \\
= & -\mu_{1} \int_{\Omega_{2}} \Delta\left(E \bar{u}_{1}\right)_{\Omega_{2}} v+\left\langle\mu_{1} \gamma_{n_{2}}\left(\nabla E \bar{u}_{1}\right), v\right\rangle_{\Lambda^{\prime}, \Lambda}  \tag{8}\\
& +\left\langle\mu_{1} \gamma_{n_{1}}\left(\nabla E \bar{u}_{1}\right), v\right\rangle_{\Lambda^{\prime}, \Lambda}, & \forall v \in V_{2} .
\end{array}
$$

Using the Stokes formula (3), it follows that the extension $E \bar{u}_{1}$ in $H_{0}^{1}(\Omega)$ verifies

$$
\begin{aligned}
\mu_{1} \int_{\Omega} \nabla\left(E \bar{u}_{1}\right) \cdot \nabla v_{1}= & \mu_{1} \int_{\Omega_{1}} \nabla \bar{u}_{1} \cdot \nabla v_{1 \mid \Omega_{1}}+\mu_{1} \int_{\Omega_{2}} \nabla\left(E \bar{u}_{1}\right)_{\mid \Omega_{2}} \cdot \nabla v_{1 \mid \Omega_{2}}, \\
= & -\mu_{1} \int_{\Omega_{1}} \Delta \bar{u}_{1} v_{1 \mid \Omega_{1}}+\left\langle\mu_{1} \gamma_{n_{1}}\left(\nabla E \bar{u}_{1}\right), v_{1 \mid \Omega_{2}}\right\rangle_{\Lambda^{\prime}, \Lambda} \\
& +\mu_{1} \int_{\Omega_{2}} \nabla\left(E \bar{u}_{1}\right)_{\mid \Omega_{2}} \cdot \nabla v_{1 \mid \Omega_{2}}, \\
= & \int_{\Omega_{1}} f_{1} v_{1 \mid \Omega_{1}}+\left\langle g, v_{1 \mid \Omega_{2}}\right\rangle_{V_{2}^{\prime}, V_{2}}, \quad \forall v_{1} \in H_{0}^{1}(\Omega) .
\end{aligned}
$$

Similarly, the function $\bar{u}_{2}$ in $V_{2}$ verifies, using (3) and (4),

$$
\begin{aligned}
\mu_{2} \int_{\Omega_{2}} \nabla \bar{u}_{2} \cdot \nabla v_{2} & =-\mu_{2} \int_{\Omega_{2}} \Delta \bar{u}_{2} v_{2}+\left\langle\mu_{2} \gamma_{n_{2}}\left(\nabla \bar{u}_{2}\right), v_{2}\right\rangle_{\Lambda^{\prime}, \Lambda} \\
& =\left\langle f_{2}, v_{2}\right\rangle_{V_{2}^{\prime}, V_{2}}-\left\langle\mu_{1} \gamma_{n_{1}}\left(\nabla \bar{u}_{1}\right), v_{2}\right\rangle_{\Lambda^{\prime}, \Lambda} \\
& =\int_{\Omega_{2}} f_{2} v_{2}+\mu_{1} \int_{\Omega_{2}} \nabla\left(E \bar{u}_{1}\right)_{\Omega_{2}} \cdot \nabla v_{2}-\left\langle g, v_{2}\right\rangle_{V_{2}^{\prime}, V_{2}}, \quad \forall v_{2} \in V_{2} .
\end{aligned}
$$

Finally, we can conclude that the couple $\left(E \bar{u}_{1}, \bar{u}_{2}\right)$ is the solution of problem (6). Thus

$$
\left(u_{1}^{g}, u_{2}^{g}\right)=\left(E \bar{u}_{1}, \bar{u}_{2}\right)
$$

and by construction of the operator $E$, the condition (7) is verified since $\left(u_{1 \mid \Omega_{1}}^{g}, u_{2}^{g}\right)=$ $\left(\bar{u}_{1}, \bar{u}_{2}\right)$.

Remark 1. The control $g$ is not unique. In fact, the extension operator $E$ constructed in the preceding proof can be defined in different ways, leading to the construction of a different control $g$ from an other extension of $\bar{u}_{1}$.
Remark 2. If $\bar{u}_{1}$ is of regularity $H^{2}\left(\Omega_{1}\right)$ and if $E \bar{u}_{1}$ is an extension which preserves this regularity on the whole domain $\Omega$, then the following weak transmission condition holds:

$$
\left\langle\mu_{1} \gamma_{n_{2}}\left(\nabla E \bar{u}_{1}\right), v\right\rangle_{\Lambda^{\prime}, \Lambda}=-\left\langle\mu_{1} \gamma_{n_{1}}\left(\nabla E \bar{u}_{1}\right), v\right\rangle_{\Lambda^{\prime}, \Lambda} .
$$

Thus, because $\Delta\left(E \bar{u}_{1}\right)_{\mid \Omega_{2}}$ belongs to $L^{2}\left(\Omega_{2}\right)$, we see from the definition of $g$ in (8) that $g$ belongs to $L^{2}\left(\Omega_{2}\right)$ and can be identified to $-\mu_{1} \Delta\left(E \bar{u}_{1}\right)_{\Omega_{2}}$. This justifies the use of the smooth extension method because the numerical approximation of the solution of the smooth extension problem (6) will converge in space with optimal rates of convergence, using $P 1$ finite elements, whereas the numerical approximation of the solution of the transmission problem (2) will not in the general case. We refer to subsection 2.3 for details. Moreover, in what follows, we will explicit some cases where such a regular extension can be constructed (in subsection 2.3 and subsection 2.4).

Remark 3. For the presentation of the smooth extension method, we considered the Laplace transmission problem with homogeneous Dirichlet boundary conditions. However, the extension of our results to the non-homogeneous case is straight forward and amounts to study a problem with homogeneous Dirichlet boundary conditions, where only the external forces are modified.
Remark 4. The hypothesis that $\Gamma_{1}$ and $\Gamma_{2}$ should not be empty could also be weakened. For example, the case where $\Omega_{2}$ is strictly included in $\Omega_{1}$ could be particularly interesting to study. In this situation $u_{2}^{g}$ is not unique in $V_{2}$ but can be searched such that it has a zero mean value. This is the subject of subsection 2.4.

### 2.2 Formulation as an optimisation problem

We saw in the previous subsection that a suitable control $g$ can be obtained by extending $u_{1}$ in the whole domain. In practice, the control $g$ can not be constructed in such a direct manner, since the solution ( $\bar{u}_{1}, \bar{u}_{2}$ ) of the initial problem (2) is unknown. To that matter, we introduce the following cost function, defined from $V_{2}^{\prime}$ to $\mathbb{R}^{+}$,

$$
\begin{equation*}
J(g)=\frac{1}{2} \int_{\Gamma}\left|u_{1}^{g}-u_{2}^{g}\right|^{2}, \tag{9}
\end{equation*}
$$

where $\left(u_{1}^{g}, u_{2}^{g}\right)$ is the unique solution of problem (6). The main idea to overcome this issue is to write the problem of finding a function $g$ such that the solution $\left(u_{1}^{g}, u_{2}^{g}\right)$ of the problem (6) verifies the equation (7), as a minimization problem on the function $J$. A first link between the minimization of $J$ and the research of a good control $g$ is given by the following theorem.
Theorem 2.2. Let $g$ be in $V_{2}^{\prime}$ be such that the solution $\left(u_{1}^{g}, u_{2}^{g}\right)$ of (6) satisfies condition (7). Then, $g$ is a minimizer of $J$.

Proof. Let $g$ be such a control, which exists according to Theorem 2.1. We know that the couple ( $u_{\left.1\right|_{\Omega_{1}}}^{g}, u_{2}^{g}$ ) is the unique solution of the weak initial problem (2). In particular, this couple satisfies the constraint of continuity through the interface $\Gamma$, which writes $u_{1 \mid \Gamma}^{g}=u_{2 \mid \Gamma}^{g}$ and implies that $J(g)=0$, so that $g$ is in fact a minimizer of $J$.

Remark 5. Theorem 2.2, along with Theorem 2.1, proves the existence of several minimizers of $J$. Actually, every control $g$ such that the solution $\left(u_{1}^{g}, u_{2}^{g}\right)$ of (6) verifies (7) is a suitable minimizer. At the end of this section, we will prove the reciprocal statement of Theorem 2.2, i.e. that for every minimizer $g$ of $J,\left(u_{1}^{g}, u_{2}^{g}\right)$ verifies the condition (7). In other words, finding a minimizer of $J$ would enable us to obtain the solution of the transmission problem (1).

The problem we face now is the minimization of the cost function $J$. To do that, classical methods such as gradient methods or quasi-Newton methods require the computation of the gradient of the cost function with respect to the control $g$. For that purpose, we will use the adjoint approach (see [12]), a suitable method to compute the gradient of a cost function which depends on the solution of a system of differential equations. The idea is the following: knowing that, for any $g$ in $V_{2}^{\prime}, J(g)$ is obtained by solving problem (6) and computing the explicit formula $\frac{1}{2} \int_{\Gamma}\left|u_{1}^{g}-u_{2}^{g}\right|^{2}$, we will prove that we can compute the gradient $\nabla J(g)$ by solving a system of linear partial differential equations, called the adjoint equations, and evaluating $\nabla J(g)$ with an explicit formula. The key here is to remark that the minimization of $J(g)$ can be seen as the minimization of the real-valued function

$$
\begin{aligned}
H_{0}^{1}(\Omega) \times V_{2} & \rightarrow \mathbb{R}^{+} \\
\left(v_{1}, v_{2}\right) & \mapsto \frac{1}{2} \int_{\Gamma}\left|v_{1 \mid \Gamma}-v_{2 \mid \Gamma}\right|^{2},
\end{aligned}
$$

under the constraint that $\left(v_{1}, v_{2}\right)$ is solution of (6). Thus, it is indicated to introduce the Lagrangian function associated to this constrained optimization problem, defined from

$$
V_{2}^{\prime} \times\left(H_{0}^{1}(\Omega) \times V_{2}\right) \times\left(H_{0}^{1}(\Omega) \times V_{2}\right)
$$

to $\mathbb{R}$ by

$$
\begin{align*}
& \mathcal{L}\left(g,\left(v_{s}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right) \\
&= \frac{1}{2}\left\|v_{1}-v_{2}\right\|_{L^{2}(\Gamma)}^{2}+\mu_{1} \int_{\Omega} \nabla v_{1} \cdot \nabla \lambda_{1}-\mu_{1} \int_{\Omega_{2}} \nabla v_{1 \mid \Omega_{2}} \cdot \nabla \lambda_{2} \\
& \quad+ \mu_{2} \int_{\Omega_{2}} \nabla v_{2} \cdot \nabla \lambda_{2}-\int_{\Omega_{1}} f_{1} \lambda_{1 \mid \Omega_{1}}-\int_{\Omega_{2}} f_{2} \lambda_{2}-\left\langle g, \lambda_{\left.1\right|_{\Omega_{2}}}\right\rangle_{V_{2}^{\prime}, V_{2}}  \tag{10}\\
& \quad+\left\langle g, \lambda_{2}\right\rangle_{V_{2}^{\prime}, V_{2}}
\end{align*}
$$

where the Lagrangian multipliers $\lambda_{1}$ and $\lambda_{2}$ are called the adjoint variables of $v_{1}$ and $v_{2}$, associated to the state equations (5a)-(5b). We also introduce the so-called adjoint equations, defined for all $g$ in $V_{2}^{\prime}$ by

$$
\begin{align*}
& \left\langle\frac{\partial \mathcal{L}}{\partial v_{1}}\left(g,\left(u_{1}^{g}, u_{2}^{g}\right),\left(\lambda_{1}, \lambda_{2}\right)\right), \delta v_{1}\right\rangle_{H^{-1}(\Omega), H_{0}^{1}(\Omega)}
\end{align*}=0, \quad \forall \delta v_{1} \in H_{0}^{1}(\Omega), ~ \begin{cases}\left\langle\frac{\partial \mathcal{L}}{\partial v_{2}}\left(g,\left(u_{1}^{g}, u_{2}^{g}\right),\left(\lambda_{1}, \lambda_{2}\right)\right), \delta v_{2}\right\rangle_{V_{2}^{\prime}, V_{2}} & =0, \quad \forall \delta v_{2} \in V_{2},\end{cases}
$$

where $\left(u_{1}^{g}, u_{2}^{g}\right)$ is the solution of the smooth extension problem (6). The computation of the differential forms $\frac{\partial \mathcal{L}}{\partial v_{1}}$ and $\frac{\partial \mathcal{L}}{\partial v_{2}}$ will be detailed later on.

Remark 6. The Lagrangian $\mathcal{L}$ is differentiable with respect to $v_{1}$ and $v_{2}$ because all its terms are quadratic or linear with respect to each of them. Thus, equations (11) are well defined.

With all these notations, we state the following theorem adapted from [12], where the use of the Lagrangian function $\mathcal{L}$ is made clear in providing a convenient way to compute the gradient of $J$.

Theorem 2.3. The mapping $g \in V_{2}^{\prime} \mapsto J(g) \in \mathbb{R}$ is differentiable and its gradient $\nabla J(g) \in$ $V_{2}^{\prime \prime}$ is given, by

$$
\begin{equation*}
\langle\nabla J(g), \delta g\rangle_{V_{2}^{\prime \prime}, V_{2}^{\prime}}=\left\langle\delta g, \lambda_{2}^{g}-\lambda_{\left.1\right|_{\Omega_{2}}}^{g}\right\rangle_{V_{2}^{\prime}, V_{2}}, \quad \forall \delta g \in V_{2}^{\prime} \tag{12}
\end{equation*}
$$

where $\left(\lambda_{1}^{g}, \lambda_{2}^{g}\right)$ verifies the adjoint equations (11).

To prove this theorem, we first need to prove the following result.

Lemma 2.1. There exists a unique $\left(\lambda_{1}^{g}, \lambda_{2}^{g}\right)$ in $H_{0}^{1}(\Omega) \times V_{2}$, solution of the adjoint equations (11).

Proof. We start by computing the partial derivatives of $\mathcal{L}$ with respect to $v_{1}$ and $v_{2}$.

Let $\varepsilon>0$. For all $g \in V_{2}^{\prime}$, and for all $\left(v_{1}, v_{2}\right),\left(\delta v_{1}, \delta v_{2}\right)$ and $\left(\lambda_{1}, \lambda_{2}\right)$ in $H_{0}^{1}(\Omega) \times V_{2}$,

$$
\begin{aligned}
\mathcal{L}\left(g,\left(v_{1}+\right.\right. & \left.\left.\varepsilon \delta v_{1}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right) \\
= & \frac{1}{2}\left\|v_{1}-v_{2}+\varepsilon \delta v_{1}\right\|_{L^{2}(\Gamma)}^{2}+\mu_{1} \int_{\Omega} \nabla\left(v_{1}+\varepsilon \delta v_{1}\right) \cdot \nabla \lambda_{1}+\mu_{2} \int_{\Omega_{2}} \nabla v_{2} \cdot \nabla \lambda_{2} \\
& -\mu_{1} \int_{\Omega_{2}} \nabla\left(v_{1 \mid \Omega_{2}}+\varepsilon \delta v_{1}\right) \cdot \nabla \lambda_{2}-\int_{\Omega_{1}} f_{1} \lambda_{\left.1\right|_{\Omega_{1}}}-\int_{\Omega_{2}} f_{2} \lambda_{2}-\left\langle g, \lambda_{1 \Omega_{\Omega_{2}}}\right\rangle_{V_{2}^{\prime}, V_{2}} \\
+ & \left\langle g, \lambda_{2}\right\rangle_{V_{2}^{\prime}, V_{2}} \\
= & \mathcal{L}\left(g,\left(v_{1}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right)+\varepsilon \mu_{1} \int_{\Omega} \nabla\left(\delta v_{1}\right) \cdot \nabla \lambda_{1}-\varepsilon \mu_{1} \int_{\Omega_{2}} \nabla\left(\delta v_{1 \mid \Omega_{2}}\right) \cdot \nabla \lambda_{2} \\
& +\varepsilon \int_{\Gamma}\left(v_{1}-v_{2}\right) \delta v_{1}+\frac{\varepsilon^{2}}{2}\left\|\delta v_{1}\right\|_{L^{2}(\Gamma)}^{2}
\end{aligned}
$$

It leads to define $\frac{\partial \mathcal{L}}{\partial v_{1}}\left(g,\left(v_{1}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right) \in H^{-1}(\Omega)$ as:

$$
\begin{aligned}
\left\langle\frac{\partial \mathcal{L}}{\partial v_{1}}\left(g,\left(v_{1}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right), \delta v_{1}\right\rangle_{H^{-1}(\Omega), H_{0}^{1}(\Omega)}= & \mu_{1} \int_{\Omega} \nabla\left(\delta v_{1}\right) \cdot \nabla \lambda_{1}-\mu_{1} \int_{\Omega_{2}} \nabla\left(\delta v_{1 \mid \Omega_{2}}\right) \cdot \nabla \lambda_{2} \\
& +\int_{\Gamma}\left(v_{1}-v_{2}\right) \delta v_{1}, \quad \forall \delta v_{1} \in H_{0}^{1}(\Omega)
\end{aligned}
$$

Similarly, we define the derivative of the Lagrangian function $\mathcal{L}$ with respect to $v_{2}$ evaluated at point $\delta v_{2}, \frac{\partial \mathcal{L}}{\partial v_{2}}\left(g,\left(v_{1}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right) \in V_{2}^{\prime}$ as:

$$
\left\langle\frac{\partial \mathcal{L}}{\partial v_{2}}\left(g,\left(v_{1}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right), \delta v_{2}\right\rangle_{V_{2}^{\prime}, V_{2}}=\mu_{2} \int_{\Omega_{2}} \nabla\left(\delta v_{2}\right) \cdot \nabla \lambda_{2}-\int_{\Gamma}\left(v_{1}-v_{2}\right) \delta v_{2}, \quad \forall \delta v_{2} \in V_{2}
$$

Then, we can deduce the adjoint problem:

$$
\begin{cases}\text { find }\left(\lambda_{1}, \lambda_{2}\right) \text { in } H_{0}^{1}(\Omega) \times V_{2} \text { such that, }  \tag{13}\\ \mu_{1} \int_{\Omega} \nabla \lambda_{1} \cdot \nabla v_{1}=\mu_{1} \int_{\Omega_{2}} \nabla \lambda_{2} \cdot \nabla v_{1 \mid \Omega_{2}}-\int_{\Gamma}\left(u_{1}^{g}-u_{2}^{g}\right) v_{1}, & \forall v_{1} \in H_{0}^{1}(\Omega) \\ \mu_{2} \int_{\Omega_{2}} \nabla \lambda_{2} \cdot \nabla v_{2}=\int_{\Gamma}\left(u_{1}^{g}-u_{2}^{g}\right) v_{2}, & \forall v_{2} \in V_{2}\end{cases}
$$

Again, we can use the Lax-Milgram Theorem to show that problem (13) admits a unique solution $\left(\lambda_{1}^{g}, \lambda_{2}^{g}\right)$.

We now have all the tools we need to prove Theorem 2.3.
Proof of Theorem 2.3. The solution $\left(u_{1}^{g}, u_{2}^{g}\right)$ of problem (6) is unique for every control $g$ in $V_{2}^{\prime}$. Thus, we can define the so called direct mapping

$$
\begin{align*}
\phi: \quad V_{2}^{\prime} & \rightarrow H_{0}^{1}(\Omega) \times V_{2} \\
g & \mapsto\left(\phi_{1}(g), \phi_{2}(g)\right)=\left(u_{1}^{g}, u_{2}^{g}\right) \tag{14}
\end{align*}
$$

Because of the linearity of the equations in (6), the mapping $\phi$ is linear, thus differentiable on $V_{2}^{\prime}$. Similarly, the mapping which for every $\left(u_{1}, u_{2}\right)$ in $H_{0}^{1}(\Omega) \times V_{2}$ associates the quadratic functional $\int_{\Gamma}\left|u_{1}-u_{2}\right|^{2}$ is differentiable. By composition, it follows that the mapping $g \mapsto J(g)$ is differentiable on $V_{2}^{\prime}$.

Taking $\left(v_{1}, v_{2}\right)=\left(u_{1}^{g}, u_{2}^{g}\right)$ the Lagrangian (10) reduces to

$$
\mathcal{L}\left(g,\left(u_{1}^{g}, u_{2}^{g}\right),\left(\lambda_{1}, \lambda_{2}\right)\right)=J(g), \forall g \in V_{2}^{\prime}, \forall\left(\lambda_{1}, \lambda_{2}\right) \in H_{0}^{1}(\Omega) \times V_{2}
$$

Then, we differentiate this previous equality with respect to $g$ for a fixed couple $\left(\lambda_{1}, \lambda_{2}\right)$ and we obtain, using the chain rule, the following equation:

$$
\begin{aligned}
<\nabla J(g), \delta g>_{V_{2}^{\prime \prime}, V_{2}^{\prime}}= & \left\langle\frac{\partial \mathcal{L}}{\partial g}\left(g,\left(u_{1}^{g}, u_{2}^{g}\right),\left(\lambda_{1}, \lambda_{2}\right)\right), \delta g\right\rangle_{V_{2}^{\prime \prime}, V_{2}} \\
& +\left\langle\frac{\partial \mathcal{L}}{\partial v_{1}}\left(g,\left(u_{1}^{g}, u_{2}^{g}\right),\left(\lambda_{1}, \lambda_{2}\right)\right), \frac{\partial \phi_{1}}{\partial g}(g) \cdot \delta g\right\rangle_{H^{-1}(\Omega) H_{0}^{1}(\Omega)} \\
& +\left\langle\frac{\partial \mathcal{L}}{\partial v_{2}}\left(g,\left(u_{1}^{g}, u_{2}^{g}\right),\left(\lambda_{1}, \lambda_{2}\right)\right), \frac{\partial \phi_{2}}{\partial g}(g) \cdot \delta g\right\rangle_{V_{2}^{\prime}, V_{2}}
\end{aligned}
$$

for all $\delta g \in V_{2}^{\prime}$. Considering $\left(\lambda_{1}, \lambda_{2}\right)=\left(\lambda_{1}^{g}, \lambda_{2}^{g}\right)$, the solution of the adjoint equations (13) (see Lemma 2.1), this reduces to

$$
<\nabla J(g), \delta g>_{V_{2}^{\prime \prime}, V_{2}^{\prime}}=\left\langle\frac{\partial \mathcal{L}}{\partial g}\left(g,\left(u_{1}^{g}, u_{2}^{g}\right),\left(\lambda_{1}^{g}, \lambda_{2}^{g}\right)\right), \delta g\right\rangle_{V_{2}^{\prime \prime}, V_{2}^{\prime}}, \quad \forall \delta g \in V_{2}^{\prime}
$$

It remains to compute the derivative of $\mathcal{L}$ with respect to $g$ at point $\delta g$ : let $\varepsilon>0$, for all $g$ and $\delta g$ in $V_{2}^{\prime}$, and for all $\left(v_{1}, v_{2}\right)$ and $\left(\lambda_{1}, \lambda_{2}\right)$ in $H_{0}^{1}(\Omega) \times V_{2}$,

$$
\begin{aligned}
\mathcal{L}\left(g+\varepsilon \delta g,\left(v_{1}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right)= & \mathcal{L}\left(g,\left(v_{1}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right)-\varepsilon\left\langle\delta g, \lambda_{\left.1\right|_{\Omega_{2}}}\right\rangle_{V_{2}^{\prime}, V_{2}} \\
& +\varepsilon\left\langle\delta g, \lambda_{2}\right\rangle_{V_{2}^{\prime}, V_{2}}
\end{aligned}
$$

which leads to the definition of $\frac{\partial \mathcal{L}}{\partial g}\left(g,\left(v_{1}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right) \in V_{2}^{\prime \prime}$ as:

$$
\left\langle\frac{\partial \mathcal{L}}{\partial g}\left(g,\left(v_{1}, v_{2}\right),\left(\lambda_{1}, \lambda_{2}\right)\right), \delta g\right\rangle_{V_{2}^{\prime \prime}, V_{2}^{\prime}}=\left\langle\delta g, \lambda_{2}-\lambda_{\left.1\right|_{\Omega_{2}}}\right\rangle_{V_{2}^{\prime}, V_{2}}, \quad \forall \delta g \in V_{2}^{\prime}
$$

Finally, we obtain an explicit expression for the gradient of $J$, which writes

$$
<\nabla J(g), \delta g>_{V_{2}^{\prime \prime}, V_{2}^{\prime}}=\left\langle\delta g, \lambda_{2}^{g}-\lambda_{\left.1\right|_{\Omega_{2}}}^{g},\right\rangle_{V_{2}^{\prime}, V_{2}}, \quad \forall \delta g \in V_{2}^{\prime}
$$

Now that we have an explicit expression for $\nabla J$, we can prove the reciprocal statement of Theorem 2.2.

Theorem 2.4. Let $g$ be a minimizer of $J$ in $V_{2}^{\prime}$. Then, the solution $\left(u_{1}^{g}, u_{2}^{g}\right)$ of (6) satisfies condition (7).

Proof. As a minimizer of $J, g$ verifies the equality $<\nabla J(g), \delta g>_{V_{2}^{\prime \prime}, V_{2}^{\prime}}=0$, for all $\delta g$ in $V_{2}^{\prime}$. According to Theorem 2.3, it follows that the couple $\left(\lambda_{1}^{g}, \lambda_{2}^{g}\right)$ verifies

$$
\left\langle\delta g, \lambda_{2}^{g}-\lambda_{\left.1\right|_{\Omega_{2}}}^{g}\right\rangle_{V_{2}^{\prime}, V_{2}}=0, \quad \forall \delta g \in V_{2}^{\prime}
$$

It comes that $\lambda_{\left.1\right|_{\Omega_{2}}}^{g}=\lambda_{2}^{g}$ in $V_{2}$ and the couple $\left(\lambda_{1 \mid \Omega_{1}}^{g}, \lambda_{2}^{g}\right)$ belongs to the space $\mathcal{V}$. Yet, the couple $\left(\lambda_{1}^{g}, \lambda_{2}^{g}\right)$ is solution of problem (13). In particular, for all $\left(v_{1}, v_{2}\right)$ in the space $\tilde{\mathcal{V}}$ defined by

$$
\tilde{\mathcal{V}}=\left\{\left(v_{1}, v_{2}\right) \in H_{0}^{1}(\Omega) \times V_{2} ; v_{1 \mid \Omega_{1}}=v_{2}\right\}
$$

the two equations in (13) write

$$
\begin{align*}
\mu_{1} \int_{\Omega} \nabla \lambda_{1}^{g} \cdot \nabla v_{1} & =\mu_{1} \int_{\Omega_{2}} \nabla \lambda_{2}^{g} \cdot \nabla v_{1 \mid \Omega_{2}}-\int_{\Gamma}\left(u_{1}^{g}-u_{2}^{g}\right) v_{1}  \tag{15}\\
\mu_{2} \int_{\Omega_{2}} \nabla \lambda_{2}^{g} \cdot \nabla v_{2} & =\int_{\Gamma}\left(u_{1}^{g}-u_{2}^{g}\right) v_{1}
\end{align*}
$$

Summing the two equations in (15) and using the fact that $\lambda_{\left.1\right|_{\Omega_{2}}}^{g}=\lambda_{2}^{g}$, we find that

$$
\mu_{1} \int_{\Omega_{1}} \nabla \lambda_{\left.1\right|_{\Omega_{1}}}^{g} \cdot \nabla v_{1 \mid \Omega_{1}}+\mu_{2} \int_{\Omega_{2}} \nabla \lambda_{2}^{g} \cdot \nabla v_{2}=0, \quad \forall\left(v_{1}, v_{2}\right) \in \tilde{\mathcal{V}}
$$

Moreover, let $\left(v_{1}, v_{2}\right)$ be in $\mathcal{V}$. We can construct an extension of $v_{1}$ in the whole space $H_{0}^{1}(\Omega)$ using $v_{2}$, as in the proof of Theorem 2.1 and we still denote $E v_{1}$ this extension. Then, the couple $\left(E v_{1}, v_{2}\right)$ belongs to $\tilde{\mathcal{V}}$ and it follows that the couple $\left(\lambda_{1 \mid \Omega_{1}}^{g}, \lambda_{2}^{g}\right)$ verifies the equation

$$
\mu_{1} \int_{\Omega_{1}} \nabla \lambda_{1 \mid \Omega_{1}}^{g} \cdot \nabla v_{1}+\mu_{2} \int_{\Omega_{2}} \nabla \lambda_{2}^{g} \cdot \nabla v_{2}=0, \quad \forall\left(v_{1}, v_{2}\right) \in \mathcal{V}
$$

We conclude that the couple $\left(\lambda_{1 \mid \Omega_{1}}^{g}, \lambda_{2}^{g}\right)$ is solution of a weak problem similar to (2) but with no external force and, thus, is the zero of $\mathcal{V}$. Especially, this implies that $\lambda_{2}^{g}=0$ and, according to the second equation of (15), that

$$
\int_{\Gamma}\left(u_{1}^{g}-u_{2}^{g}\right) v_{2}=0, \quad \forall v_{2} \in V_{2}
$$

Because $u_{\left.1\right|_{\Gamma}}^{g}$ and $u_{\left.2\right|_{\Gamma}}^{g}$ belong to $H^{1 / 2}(\Gamma)$ the previous equality implies that $u_{\left.1\right|_{\Gamma}}^{g}=u_{\left.2\right|_{\Gamma}}^{g}$ in $H^{1 / 2}(\Gamma)$, so that the couple $\left(u_{\left.1\right|_{\Omega_{1}}}^{g}, u_{2}^{g}\right)$ belongs to the space $\mathcal{V}$. As the unique solution of problem (6), the pair $\left(u_{1}^{g}, u_{2}^{g}\right)$ verifies, in particular, for all $\left(v_{1}, v_{2}\right)$ in $\tilde{\mathcal{V}}$,

$$
\begin{align*}
\mu_{1} \int_{\Omega} \nabla u_{1}^{g} \cdot \nabla v_{1} & =\int_{\Omega_{1}} f_{1} v_{\left.1\right|_{\Omega_{1}}}+\left\langle g, v_{2}\right\rangle_{V_{2}^{\prime}, V_{2}} \\
\mu_{2} \int_{\Omega_{2}} \nabla u_{2}^{g} \cdot \nabla v_{2} & =\int_{\Omega_{2}} f_{2} v_{2}-\left\langle g, v_{2}\right\rangle_{V_{2}^{\prime}, V_{2}}+\mu_{1} \int_{\Omega_{2}} \nabla u_{1}^{g} \cdot \nabla v_{\left.1\right|_{\Omega_{2}}} \tag{16}
\end{align*}
$$

Summing the two equations in (16) it follows that the pair $\left(u_{\left.1\right|_{\Omega_{1}}}^{g}, u_{2}^{g}\right)$, which belongs to $\mathcal{V}$, verifies the equation

$$
\mu_{1} \int_{\Omega_{1}} \nabla u_{1 \mid \Omega_{1}}^{g} \cdot \nabla v_{1 \mid \Omega_{1}}+\mu_{2} \int_{\Omega_{2}} \nabla u_{2}^{g} \cdot \nabla v_{2}=\int_{\Omega_{1}} f_{1} v_{1 \mid \Omega_{1}}+\int_{\Omega_{2}} f_{2} v_{2}, \quad \forall\left(v_{1}, v_{2}\right) \in \tilde{\mathcal{V}}
$$

As before, for all $\left(v_{1}, v_{2}\right)$ in $\mathcal{V}$, the couple $\left(E v_{1}, v_{2}\right)$ belongs to $\tilde{\mathcal{V}}$ and, finally, the couple $\left(u_{1 \mid \Omega_{1}}^{g}, u_{2}^{g}\right)$ verifies the equation

$$
\mu_{1} \int_{\Omega_{1}} \nabla u_{\left.1\right|_{\Omega_{1}}}^{g} \cdot \nabla v_{1}+\mu_{2} \int_{\Omega_{2}} \nabla u_{2}^{g} \cdot \nabla v_{2}=\int_{\Omega_{1}} f_{1} v_{1}+\int_{\Omega_{2}} f_{2} v_{2}, \quad \forall\left(v_{1}, v_{2}\right) \in \mathcal{V}
$$

Thus, $\left(u_{\left.1\right|_{\Omega_{1}}}^{g}, u_{2}^{g}\right)$ is the unique solution to the initial coupled problem (2) and verifies, $a$ fortiori, condition (7).

So, the minimization of the function $J$ in $V_{2}^{\prime}$ is equivalent to the resolution of the problem of finding a suitable control $g$ such that the solution $\left(u_{1}^{g}, u_{2}^{g}\right)$ verifies (7) and thus, to the resolution of the initial coupled problem (2). As we will explain it in the next section, this minimization formulation offers, for the purpose of the numerical resolution of problem (2), the advantage of conserving the optimal rates of convergence, even with non-conformal meshes.

Remark 7. Theorem 2.4 also shows that all extrema of $J$ correspond to global minimizers. Actually, in Theorem 2.4, it is enough to assume that $g$ verifies $<\nabla J(g), \delta g>_{V_{2}^{\prime \prime}, V_{2}^{\prime}}=0$ i.e., that $J(g)$ is a local extremum of $J$, to conclude that the solution $\left(u_{1}^{g}, u_{2}^{g}\right)$ of (6) verifies the condition (7). According to Theorem 2.2, this implies that $g$ is a global minimizer of $J$.

### 2.3 On convergence rates for the numerical method

The main difficulty with the numerical simulation of the Laplace transmission problem, that we will also find in other transmission problems, is to recover the optimal rate of convergence when the solution is more regular. Actually, even if $f_{1}$ belongs to $L^{2}\left(\Omega_{1}\right), f_{2}$ to $L^{2}\left(\Omega_{2}\right)$ and the boundaries of $\Omega_{1}$ and $\Omega_{2}$ are smooth, the solution of problem (2) is not of regularity $H^{2}$ in the whole domain $\Omega$, because of the jump in its gradient across the interface $\Gamma$. As a consequence, the standard finite element method with $P 1$ elements does not enable to recover the optimal rates of convergence in $L^{2}$ and $H^{1}$-norms, unless we use a conformal mesh that fully represents the interface $\Gamma$ (see Figure 5).

Nevertheless, with such a regularity on the data and the boundaries, one can show that the solution of problem (2) is partially of regularity $H^{2}$, in the following sense: $\bar{u}_{1}$ belongs to $H^{2}\left(\Omega_{1}\right) \cap V_{1}$ and $\bar{u}_{2}$ to $H^{2}\left(\Omega_{2}\right) \cap V_{2}$ (see [14, Theorem 5.2.1]). Thus, if we can extend $\bar{u}_{1}$ in the whole domain $\Omega$ with regularity $H^{2}$, we know how to construct a suitable control $g$ in $L^{2}\left(\Omega_{2}\right)$ such that the solution of the smooth extension problem (6) is also of regularity $H^{2}: u_{1}^{g}$ belongs to $H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$ and $u_{2}^{g}$ belongs to $H^{2}\left(\Omega_{2}\right) \cap V_{2}$.

The advantage to consider the smooth extension problem (6), instead of problem (2), is now clear: the approximations of the two functions $u_{1}^{g}$ and $u_{2}^{g}$ are defined on two independent meshes and the two subproblems in (6) can be solved consecutively using standard finite element function spaces. These make the smooth extension method a fictitious domain method with optimal rates of convergence, even with non-conformal meshes.

To precise more the possibility to actually extend a function $u_{1}$ which belongs to $H^{2}\left(\Omega_{1}\right)$ into the whole space $\Omega$ with the same regularity, we state the following theorem, which requires some (rather weak) assumptions on the domains.

Theorem 2.5. Let $\Omega \in \mathbb{R}^{n}$, with $n \in\{2,3\}$, be a bounded open set with Lipschitz boundary $\partial \Omega$, and let $\Gamma$ be an interface that divides $\Omega$ into two bounded open connected subdomains $\Omega_{1}$ and $\Omega_{2}$ with Lipschitz boundaries. Consider a function $u_{1}$ in $H^{2}\left(\Omega_{1}\right) \cap V_{1}$. Moreover, if $n=3$, we assume that $\Gamma \cap \partial \Omega$ is a curve of regularity $C^{2}$. If $\Gamma$ and $\partial \Omega$ have at least $C^{2}$-regularity in a neighborhood of each element of $\Gamma \cap \partial \Omega$ (a curve in three space dimensions and a point in two space dimensions), there exists a regular extension $u \in H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$ such that $u=u_{1}$ in $\Omega_{1}$.

Remark 8. The case where $\partial \Omega_{2}=\Gamma$ is particularly easy to study and is the subject of subsection 2.4. Thus we will only consider the case where $\partial \Omega_{1} \cap \partial \Omega$ and $\partial \Omega_{2} \cap \partial \Omega$ are both non empty in the following proof.

Proof. Suppose that $\partial \Omega_{1} \cap \partial \Omega$ and $\partial \Omega_{2} \cap \partial \Omega$ are both non empty. For the sake of simplicity, we consider the two-dimensional case, i.e. $\Omega \subset \mathbb{R}^{2}$, but the extension to the three-dimensional case can be performed in a similar manner. In two space dimension, the interface $\Gamma$ intersects the boundary $\partial \Omega$ in two points that we denote by $x_{1}$ and $x_{2}$, i.e. $\Gamma \cap \partial \Omega=\left\{x_{1}, x_{2}\right\}$.

For $i \in\{1,2\}$, let us denote by $B_{i}=B\left(x_{i}, \varepsilon_{i}\right)$ the open ball of radius $\varepsilon_{i}$ and center $x_{i}$. Moreover, we consider two open balls $B_{i}^{+}=B\left(x_{i}, \varepsilon_{i}+\varepsilon_{i} / 4\right)$ and $B_{i}^{-}=B\left(x_{i}, \varepsilon_{i}-\varepsilon_{i} / 4\right)$, such that $B_{i}^{-} \subset B_{i} \subset B_{i}^{+}$. Furthermore, we should ensure that, for $i \in\{1,2\}, \varepsilon_{i}$ is small enough such that $\Gamma \cap B_{i}^{+}$and $\partial \Omega \cap B_{i}^{+}$are included in the neighborhood of $x_{i}$ where both


Figure 1: Two-dimensional example of a geometry considered for the proof of the smooth extension.
frontiers $\Gamma$ and $\partial \Omega$ have $C^{2}$-regularity. A quite general example of such a configuration is illustrated in Figure 1.

Let $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{R}^{+}\right)$be a real-valued function satisfying

$$
\varphi(x)= \begin{cases}1 & \text { if } x<3 / 4 \\ 0 & \text { if } x>1\end{cases}
$$

Then we can define, for $i \in\{1,2\}$, the functions

$$
\begin{aligned}
\varphi_{i}: \Omega_{1} & \rightarrow[0,1] \\
x & \mapsto \varphi\left(\frac{\left|x-x_{i}\right|}{\varepsilon}\right)
\end{aligned}
$$

whose regularity is $\mathcal{C}^{\infty}\left(\Omega_{1}\right)$, and the functions

$$
\begin{aligned}
\xi_{i}: \Omega_{1} & \rightarrow \mathbb{R} \\
x & \mapsto
\end{aligned} u_{1}(x) \varphi_{i}(x) .
$$

Thus, we can write the following decomposition for the function $u_{1}$ on $\Omega_{1}$ :

$$
u_{1}=u_{1} \varphi_{1}+u_{1} \varphi_{2}+u_{1}\left(1-\varphi_{1}-\varphi_{2}\right)=\xi_{1}+\xi_{2}+\xi_{0}
$$

where $\xi_{0}=u_{1}\left(1-\varphi_{1}-\varphi_{2}\right)$. The problem to construct a regular extension for $u_{1}$ is equivalent to extend the three functions $\xi_{0}, \xi_{1}$ and $\xi_{2}$ to the whole domain $\Omega$. Moreover, let us remark that we have

$$
\begin{array}{llll}
\xi_{i}(x)=u_{1}(x), & \forall x \in B_{i}^{-} \cap \Omega_{1}, & \forall i \in\{0,1\}, \\
\xi_{i}(x)=0, & \forall x \in \Omega_{1} \backslash\left(B_{i} \cap \Omega_{1}\right), & \forall i \in\{0,1\},
\end{array}
$$

and

$$
\begin{array}{lll}
\xi_{0}(x)=0, & \forall x \in B_{i}^{-} \cap \Omega_{1}, & \forall i \in\{0,1\}, \\
\xi_{0}(x)=u_{1}(x), & \forall x \in \Omega_{1} \backslash\left(B_{i} \cap \Omega_{1}\right), & \forall i \in\{0,1\} .
\end{array}
$$



Figure 2: Construction of the function $\xi_{0}$ and its regular extension, $\tilde{\xi}_{0}$.

We begin with the function $\xi_{0}$. For that matter, we make use of Stein's theorem ( $\left[1\right.$, Theorem 5.24]), which states that, for any bounded Lipschitz domain $\mathcal{O} \subset \mathbb{R}^{n}$, there exists a total extension operator, i.e. an extension operator from $H^{m}(\mathcal{O})$ into the whole space $H^{m}\left(\mathbb{R}^{n}\right)$, for all $m \geq 0$. Then, $\Omega_{1}$ being a bounded Lipschitz domain, there exists an extension operator from $H^{2}\left(\Omega_{1}\right)$ into the whole space $H^{2}\left(\mathbb{R}^{2}\right)$. We denote by $E \xi_{0}$ the restriction to $\Omega$ of such a regular extension of $\xi_{0}$ to $\mathbb{R}^{2}$. This extension is regular enough but it is not zero on the boundary $\partial \Omega \cap \partial \Omega_{2}$. Therefore, we introduce a compact set $\mathfrak{V}$ such that $\left(\partial \Omega \cap \partial \Omega_{2}\right) \subset \mathfrak{V}$ and $\mathfrak{V} \cap \Omega_{1}=\emptyset$, and an open neighborhood of $\mathfrak{V}$, denoted by $\mathfrak{U}$, such that $\left(\mathfrak{U} \cap \Omega_{1}\right) \subset\left(B_{1}^{-} \cup B_{2}^{-}\right)$. We also define a smooth cut-off function $\psi: x \in \Omega \rightarrow \psi(x)$ such that $\psi(x)=0$ in $\mathfrak{V}$, and $\psi(x)=1$ in $\Omega \backslash \mathfrak{U}$. Then, the function $\tilde{\xi}_{0}=\psi E \xi_{0}$ is an extension of $\xi_{0}$ to the whole domain $\Omega$ of regularity $H^{2}$, which satisfies $\tilde{\xi_{0}}=0$ on $\partial \Omega$. The construction of $\tilde{\xi}_{0}$ is represented in Figure 2.

Let us now proceed with the construction of a smooth extension of $\xi_{i}$, for $i \in\{1,2\}$. First, since $\xi_{i}=0$ in $\Omega_{1} \backslash\left(B_{i} \cap \Omega_{1}\right)$, we can extend it by 0 in $\Omega_{2} \backslash\left(B_{i}^{+} \cap \Omega_{2}\right)$. Then, it remains to construct the extension in $B_{i}^{+} \cap \Omega_{2}$. Due to the $C^{2}$-regularity of the interface $\Gamma$ and the boundary $\partial \Omega$ inside the balls $B_{1}^{+}$and $B_{2}^{+}$(if $\varepsilon_{1}$ and $\varepsilon_{2}$ are chosen small enough), there exist two $C^{2}$-diffeomorphisms $\chi_{1}$ and $\chi_{2}$ that map $B_{1}^{+}$and $B_{2}^{+}$respectively into the open unit square $Q$, and such that

$$
\begin{array}{rlrl}
\chi_{i}\left(x_{i}\right) & =0, & & \forall i \in\{1,2\}, \\
\chi_{i}\left(\Gamma \cap B_{i}^{+}\right) & =\{0\} \times] 0, \frac{1}{2}[, & & \forall i \in\{1,2\}, \\
\chi_{i}\left(\partial \Omega \cap B_{i}^{+}\right) & =]-\frac{1}{2}, \frac{1}{2}[\times\{0\}, & \forall i \in\{1,2\} .
\end{array}
$$

For $i \in\{1,2\}$, we define $\hat{\xi}_{i}=\xi_{i} \circ \chi_{i}^{-1}$ in

$$
\left.\chi_{i}\left(B_{i}^{+} \cap \Omega_{1}\right)=\right]-\frac{1}{2}, 0[\times] 0, \frac{1}{2}[,
$$

which belongs to $H^{2}\left(\chi_{i}\left(B_{i}^{+} \cap \Omega_{1}\right)\right)$ (see [1, Theorem 3.41]) and, recalling that $\xi_{i}=0$ in $\left(B_{i}^{+} \backslash B_{i}\right) \cap \Omega_{1}$, it follows that $\hat{\xi}_{i}=0$ in $\chi_{i}\left(\left(B_{i}^{+} \backslash B_{i}\right) \cap \Omega_{1}\right)$. Then, a regular extension
of $\hat{\xi}_{i}$ to the whole square $Q$ can be obtained using Babič's extension (see [6] or [1, Theorem 5.19]), defined for $i \in\{1,2\}$ by

$$
E \hat{\xi}_{i}(x, y)=\left\{\begin{array}{ll}
\hat{\xi}(x, y), & \text { if } x \leq 0, \\
-3 \hat{\xi}(-x, y)+4 \hat{\xi}\left(-\frac{x}{2}, y\right), & \text { if } x>0,
\end{array} \quad \forall(x, y) \in\right]-\frac{1}{2}, \frac{1}{2}[\times] 0, \frac{1}{2}[.
$$

Clearly, for $i \in\{1,2\}$, the extension $E \hat{\xi}_{i}$ and its first derivative are continuous through the interface $\{0\} \times] 0, \frac{1}{2}\left[\right.$, such that $E \hat{\xi}_{i}$ well defines an extension of $\hat{\xi}_{i}$ in $H^{2}\left(\chi_{i}\left(B_{i}^{+} \cap \Omega\right)\right)$. However, the extension $E \hat{\xi}_{i}$ is not zero on the exterior frontier $\left.\frac{1}{2} \times\right] 0, \frac{1}{2}[$, thus we multiply it by a smooth cut-off function that is zero for $x>1 / 2-\varepsilon$ and equal to 1 for $x<1 / 2-2 \varepsilon$, for some $\varepsilon<1 / 4$. Let us denote by $\bar{\xi}_{i}$ the product of $E \hat{\xi}_{i}$ by this cut-off function, which belongs to $H^{2}\left(\chi_{i}\left(B_{i}^{+} \cap \Omega\right)\right)$ and is zero on $\chi_{i}\left(\partial \Omega \cap B_{i}^{+}\right)$. Mapping this extension into $B_{i}^{+}$ we obtain $\bar{\xi}_{i} \circ \chi_{i}$, which then is an extension of $\xi_{i}$ in $B_{i}^{+} \cap \Omega$ of regularity $H^{2}$. Now, we can construct the smooth extension of $\xi_{i}$ to the whole domain $\Omega$, denoted $\tilde{\xi}_{i}$ and defined by

$$
\begin{aligned}
\tilde{\xi}_{i} & =\xi_{i}, & & \text { in } \Omega_{1}, \\
\tilde{\xi}_{i} & =0, & & \text { in } \Omega_{2} \backslash\left(B_{i}^{+} \cap \Omega_{2}\right), \\
\tilde{\xi}_{i} & =\bar{\xi}_{i} \circ \chi_{i}, & & \text { in } B_{i}^{+} \cap \Omega .
\end{aligned}
$$

Thus, for $i \in\{1,2\}, \tilde{\xi}_{i}$ is an extension of $\xi_{i}$ which belongs to $H^{2}(\Omega)$ and is zero on $\partial \Omega$.
Finally, we define

$$
\tilde{u}=\tilde{\xi}_{0}+\tilde{\xi}_{1}+\tilde{\xi}_{2}, \text { in } \Omega
$$

and it follows that $\tilde{u}$ belongs to $H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$ and $\tilde{u}=u_{1}$ in $\Omega_{1}$.
Remark 9. The proof of Theorem 2.5 in three dimensions of space follows the same path. We start by considering three tubular neighborhoods of the intersection $\Gamma \cap \partial \Omega$ with circular cross sections and decompose the function $u_{1}$ in two regular functions $\xi_{0}$ and $\xi_{1}$. The function $\xi_{0}$ is zero inside the smallest neighborhood and is equal to $u_{1}$ far from the intersection $\Gamma \cap \partial \Omega$, while the function $\xi_{1}$ is equal to $u_{1}$ inside the smallest neighborhood and is zero far from $\Gamma \cap \partial \Omega$. On one hand, the extension of the function $\xi_{0}$ is done using Stein's extension theorem. On the other hand, the extension of the function $\xi_{1}$ is conducted by transforming the largest tubular neighborhood into the unit torus with square cross section with a $C^{2}$-diffeomorphism. Then, Babič's extension theorem is used in the unit torus and the resulting extension is brought back to the initial domain. In both cases, the extensions are multiplied by cut-off functions to ensure that they are zero on the exterior frontier $\partial \Omega$.

To conclude, the smooth extension formulation of transmission problems is possible as soon as the domains have Lipschitz regularities. Moreover, if their frontiers are more regular, such that the hypotheses of Theorem 2.5 are satisfied, then the smooth extension method applied to the numerical resolution of transmission problems has optimal rates of convergence.

A well-suited situation where the domains only have Lipschitz regularities is presented in the next subsection.

### 2.4 Particular case for strictly included domains

We explained in the previous subsection that the smooth extension method preserves the optimal order of convergence if there exists a $H^{2}$-extension of the solution of the Laplace transmission problem from $\Omega_{1}$ into the whole space $H_{0}^{1}(\Omega)$. However, this may not always be possible, depending on the geometry of the problem. A most suitable case is, for
example, the one where $\Omega_{2}$ is strictly included in $\Omega_{1}$. Indeed, we will show that the construction of such a regular extension is always possible in this case, so that the smooth extension method could always converge with optimal rates. Then, in this subsection we relax the hypothesis $i v$ ) in Assumption $\left(H_{1}\right)$, i.e. we suppose that $\Omega_{2}$ is strictly included in $\Omega_{1}$ and that $\Gamma_{1}=\partial \Omega$. Then, the set of hypotheses in this subsection is the following:
i) Domain $\Omega$ is a bounded connected Lipschitz domain of $\mathbb{R}^{n}$.
ii) Domain $\Omega$ is divided in two subdomains, $\Omega_{1}$ and $\Omega_{2}$, which have Lipschitz boundaries.
$\left(H_{2}\right)$
iii) The interface is $\Gamma=\partial \Omega_{2}$.
$i v)$ The remaining boundary of $\Omega_{1}$ is $\Gamma_{1}=\partial \Omega$.
Again, we are able to define a smooth extension formulation for the Laplace transmission problem.

We still consider the Laplace problem presented in subsection 2.1:

$$
\left\{\begin{array}{r}
\text { find } u: \Omega \rightarrow \mathbb{R} \text { such that }  \tag{17}\\
-\operatorname{div}(\mu \nabla u)=f \quad \text { in } \quad \Omega \\
u=0 \quad \text { on } \quad \partial \Omega
\end{array}\right.
$$

where $\mu$ and $f$ are defined by

$$
\begin{aligned}
\mu & =\left\{\begin{array}{lll}
\mu_{1} & \text { in } & \Omega_{1} \\
\mu_{2} & \text { in } & \Omega_{2}
\end{array}\right. \\
f & =\left\{\begin{array}{lll}
f_{1} & \text { in } & \Omega_{1} \\
f_{2} & \text { in } & \Omega_{2}
\end{array}\right.
\end{aligned}
$$

For the same reasons, problem (17) admits a unique solution $\bar{u}$ in $H_{0}^{1}(\Omega)$.
Let us define the following subspace of $H^{1}\left(\Omega_{2}\right)^{\prime}$ :

$$
V=\left\{v \in H^{1}\left(\Omega_{2}\right)^{\prime},\langle v, 1\rangle_{H^{1}\left(\Omega_{2}\right)^{\prime}, H^{1}\left(\Omega_{2}\right)}=0\right\}
$$

The smooth extension problem associated to problem (17) is the problem of finding a suitable control $g$ in $V$, such that the solution of the following weak problem,

$$
\left\{\begin{array}{rlrl}
\text { find } u_{1} \in H_{0}^{1}(\Omega) \text { and } u_{2} \in H^{1}\left(\Omega_{2}\right) \cap L_{0}^{2}\left(\Omega_{2}\right) \text { such that } &  \tag{18}\\
\mu_{1} \int_{\Omega} \nabla u_{1} \cdot \nabla v_{1}= & \int_{\Omega_{1}} f_{1} v_{1 \mid \Omega_{1}}+\left\langle g, v_{1 \mid \Omega_{2}}\right\rangle_{H^{1}\left(\Omega_{2}\right)^{\prime}, H^{1}\left(\Omega_{2}\right)} & \\
& +\frac{1}{\left|\Omega_{2}\right|} \int_{\Omega_{2}} f_{2} \int_{\Omega_{2}} v_{1 \mid \Omega_{2}}, & \forall v_{1} \in H_{0}^{1}(\Omega) \\
\mu_{2} \int_{\Omega_{2}} \nabla u_{2} \cdot \nabla v_{2}= & \int_{\Omega_{2}} f_{2} v_{2}+\left\langle\gamma_{n_{2}}\left(\mu_{1} \nabla u_{1}\right), v\right\rangle_{H^{-1 / 2}(\Gamma), H^{1 / 2}(\Gamma)}, & \forall v_{2} \in H^{1}\left(\Omega_{2}\right)
\end{array}\right.
$$

enables to recover the solution $\bar{u}$ of the initial Laplace problem (17).
Theorem 2.6. Let $\Omega$ be a domain that satisfies Assumption $\left(H_{2}\right)$. For every $f_{1} \in L^{2}\left(\Omega_{1}\right)$, every $f_{2} \in L^{2}\left(\Omega_{2}\right)$ and every $g \in V$, problem (18) admits a unique solution $\left(u_{1}^{g}, u_{2}^{g}\right)$ in $H_{0}^{1}(\Omega) \times\left(H^{1}\left(\Omega_{2}\right) \cap L_{0}^{2}\left(\Omega_{2}\right)\right)$. Moreover, there exists a control $g$ such that $\left(u_{1}^{g}, u_{2}^{g}\right)$ satisfies the following equality:

$$
\begin{equation*}
u_{1}^{g}-\frac{1}{|\Gamma|} \int_{\Gamma} u_{1}^{g}=u_{2}^{g}-\frac{1}{|\Gamma|} \int_{\Gamma} u_{2}^{g}, \quad \text { on } \Gamma \tag{19}
\end{equation*}
$$

and we can recover the solution of the initial Laplace problem (17):

$$
\begin{align*}
\bar{u}_{\mid \Omega_{1}} & =u_{1 \mid \Omega_{1}}^{g} \\
\bar{u}_{\mid \Omega_{2}} & =u_{2}^{g}+\frac{1}{|\Gamma|}\left(\int_{\Gamma}\left(u_{1}^{g}-u_{2}^{g}\right)\right) \tag{20}
\end{align*}
$$

Proof. Problem (18) is well-posed under the compatibility condition

$$
\int_{\Omega_{2}} f_{2}=-\left\langle\gamma_{n_{2}}\left(\mu_{1} \nabla u_{1}\right), 1\right\rangle_{H^{-1 / 2}(\Gamma), H^{1 / 2}(\Gamma)}
$$

Let us prove that this condition is satisfied. Using the Stokes formula (3), we know that, for all $v \in H_{0}^{1}(\Omega)$,

$$
\begin{aligned}
&\left\langle\gamma_{n_{2}}\right.\left.\left(\mu_{1} \nabla u_{1}\right), v\right\rangle_{H^{-1 / 2}(\Gamma), H^{1 / 2}(\Gamma)} \\
& \quad=\int_{\Omega_{1}} f_{1} v-\mu_{1} \int_{\Omega_{1}} \nabla u_{1} \cdot \nabla v \\
& \quad= \mu_{1} \int_{\Omega_{2}} \nabla u_{1} \cdot \nabla v-\langle g, v\rangle_{H^{1}\left(\Omega_{2}\right)^{\prime}, H^{1}\left(\Omega_{2}\right)}-\frac{1}{\left|\Omega_{2}\right|} \int_{\Omega_{2}} f_{2} \int_{\Omega_{2}} v
\end{aligned}
$$

In particular, taking $v=1$ and remembering that $g$ belongs to $V$ it follows that

$$
\left\langle\gamma_{n_{2}}\left(\mu_{1} \nabla u_{1}\right), 1\right\rangle_{H^{-1 / 2}(\Gamma), H^{1 / 2}(\Gamma)}=-\int_{\Omega_{2}} f_{2}
$$

which is the desired equality and prove the first part of Theorem 2.6.
Now, we prove the existence of a suitable control $g$ in $V$ such that the equality (19) holds. We define

$$
\begin{aligned}
& \bar{u}_{1}=\bar{u}_{\mid \Omega_{1}}, \\
& \bar{u}_{2}=\bar{u}_{\mid \Omega_{2}},
\end{aligned}
$$

and construct an extension of $\bar{u}_{1}$ in the whole space $H_{0}^{1}(\Omega)$. This is possible according to the Stein extension theorem (see [1, Theorem 5.24]) which states that, for any bounded Lipschitz domain $\mathcal{O}$ of $\mathbb{R}^{n}$, there exists a total extension operator, i.e. an extension operator from $H^{m}(\mathcal{O})$ into the whole space $H^{m}\left(\mathbb{R}^{n}\right)$, for all $m \geq 0$. We apply this theorem to $\bar{u}_{1}$ and we denote by $E \bar{u}_{1}$ the restriction to $\Omega$ of its total extension. Thus, if $\bar{u}_{1}$ belongs to $H^{m}\left(\Omega_{1}\right) \cap V_{1}$ for $m \geq 1$, the extension $E \bar{u}_{1}$ belongs to $H^{m}(\Omega) \cap H_{0}^{1}(\Omega)$. In particular, if $\bar{u}_{1}$ belongs to $H^{2}\left(\Omega_{1}\right) \cap V_{1}, E \bar{u}_{1}$ belongs to $H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$. Then, we construct $g$ in $V$ such that

$$
\begin{align*}
\langle g, v\rangle_{V_{2}^{\prime}, V_{2}}= & \mu_{1} \int_{\Omega_{2}} \nabla\left(E \bar{u}_{1}\right) \cdot \nabla v-\frac{1}{\left|\Omega_{2}\right|} \int_{\Omega_{2}} f_{2} \int_{\Omega_{2}} v  \tag{21}\\
& +\mu_{1}\left\langle\gamma_{n_{1}}\left(\nabla E \bar{u}_{1}\right), v\right\rangle_{H^{-1 / 2}(\Gamma), H^{1 / 2}(\Gamma)}, \quad \forall v \in H^{1}\left(\Omega_{2}\right)
\end{align*}
$$

Using the Stokes formula (3), it follows that the extension $E \bar{u}_{1}$ satisfies, for all $v_{1} \in H_{0}^{1}(\Omega)$,

$$
\begin{aligned}
\mu_{1} \int_{\Omega} \nabla\left(E \bar{u}_{1}\right) \cdot \nabla v_{1} & =\mu_{1} \int_{\Omega_{1}} \nabla\left(E \bar{u}_{1}\right) \cdot \nabla v_{1}+\mu_{1} \int_{\Omega_{2}} \nabla\left(E \bar{u}_{1}\right) \cdot \nabla v_{1} \\
& =\int_{\Omega_{1}} f_{1} v_{1}+\mu_{1}\left\langle\gamma_{n_{1}}\left(\nabla E \bar{u}_{1}\right), v\right\rangle_{H^{-1 / 2}(\Gamma), H^{1 / 2}(\Gamma)}+\mu_{1} \int_{\Omega_{2}} \nabla\left(E \bar{u}_{1}\right) \cdot \nabla v_{1} \\
& =\int_{\Omega_{1}} f_{1} v_{1}+\langle g, v\rangle_{V_{2}^{\prime}, V_{2}}+\frac{1}{\left|\Omega_{2}\right|} \int_{\Omega_{2}} f_{2} \int_{\Omega_{2}} v
\end{aligned}
$$

Similarly, the function $\tilde{u}_{2}=\bar{u}_{2}-\frac{1}{\left|\Omega_{2}\right|} \int_{\Omega_{2}} \bar{u}_{2}$ belongs to $H^{1}\left(\Omega_{2}\right) \cap L_{0}^{2}\left(\Omega_{2}\right)$ and satisfies, for all $v_{2} \in H^{1}\left(\Omega_{2}\right)$,

$$
\begin{aligned}
\mu_{2} \int_{\Omega_{2}} \nabla \tilde{u}_{2} \cdot \nabla v_{2} & =\mu_{2} \int_{\Omega_{2}} \nabla \bar{u}_{2} \cdot \nabla v_{2} \\
& =\int_{\Omega_{2}} f_{2} v_{2}+\mu_{1}\left\langle\gamma_{n_{1}}\left(\nabla E \bar{u}_{1}\right), v\right\rangle_{H^{-1 / 2}(\Gamma), H^{1 / 2}(\Gamma)}
\end{aligned}
$$

It follows that $\left(E \bar{u}_{1}, \tilde{u}_{2}\right)$ is the solution of (18), i.e.

$$
\left(u_{1}^{g}, u_{2}^{g}\right)=\left(E \bar{u}_{1}, \tilde{u}_{2}\right)
$$

Moreover, on $\Gamma$ we have

$$
\begin{aligned}
u_{1}^{g}-\frac{1}{|\Gamma|} \int_{\Gamma} u_{1}^{g} & =\bar{u}_{2}-\frac{1}{|\Gamma|} \int_{\Gamma} \bar{u}_{2} \\
& =u_{2}^{g}+\frac{1}{\left|\Omega_{2}\right|} \int_{\Omega_{2}} \bar{u}_{2}-\frac{1}{|\Gamma|} \int_{\Gamma}\left(u_{2}^{g}+\frac{1}{\left|\Omega_{2}\right|} \int_{\Omega_{2}} \bar{u}_{2}\right) \\
& =u_{2}^{g}-\frac{1}{|\Gamma|} \int_{\Gamma} u_{2}^{g}
\end{aligned}
$$

which proves equality (19).
Finally, the solution $\left(\bar{u}_{1}, \overline{u_{2}}\right)$ of the Laplace problem (17) can be recover from $\left(u_{1}^{g}, u_{2}^{g}\right)$ :

$$
\begin{aligned}
\bar{u}_{1} & =u_{1 \mid \Omega_{1}}^{g} \\
\bar{u}_{2} & =u_{2}^{g}+C,
\end{aligned}
$$

where the constant $C$ can be determined by integrating the difference $\bar{u}_{2}-u_{2}^{g}$ on $\Gamma$ :

$$
\int_{\Gamma} C=\int_{\Gamma}\left(\bar{u}_{2}-u_{2}^{g}\right)=\int_{\Gamma}\left(u_{1}^{g}-u_{2}^{g}\right)
$$

which leads to

$$
C=\frac{1}{|\Gamma|} \int_{\Gamma}\left(u_{1}^{g}-u_{2}^{g}\right)
$$

Theorem 2.6 implies that, if we can construct a suitable control $g$, then we are able to recover the solution of the Laplace transmission problem (17) from the resolution of the smooth extension problem (18). To obtain such a control, we can solve the following optimization problem:

$$
\begin{equation*}
\min _{g \in V} J(g)=\frac{1}{2} \int_{\Gamma}\left|u_{1}^{g}-u_{1}^{g}+\frac{1}{|\Gamma|} \int_{\Gamma}\left(u_{2}^{g}-u_{1}^{g}\right)\right|^{2} \tag{22}
\end{equation*}
$$

Problem (22) admits at least one solution, according to Theorem 2.6, which corresponds to a zero of the cost function $J$. As before, one can show that if $g$ is a minimizer of $J$, the solution $\left(u_{1}^{g}, u_{2}^{g}\right)$ of problem (18) can be used to recover the solution of the Laplace transmission problem, with formulas (20). The proof is similar to what we have done in subsection 2.2 and is let to the reader.

## 3 Validation of the method

In subsection 2.2 we have proved the equivalence between the initial coupled problem and the smooth extension formulation stated as a minimization problem. Based on this result, we detail the numerical procedure used to solve the smooth extension formulation of the Laplace transmission problem presented above. In particular, in subsection 3.1, we explain how to minimize the function $J$ defined in (9) using its gradient. Thereafter, we present in subsection 3.2 , subsection 3.3 and subsection 3.4 some numerical experiments obtained through this process.

### 3.1 Numerical procedure for the smooth extension method

To solve the smooth extension formulation of problem (2), we consider an optimization problem whose solution allows to directly recover the solution of the initial transmission problem. In the case of the Laplace transmission problem, it gives the right control term $g$ for which the solution $u_{1}^{g}$ of problem (6) is a smooth extension of $\bar{u}_{1}$ in the whole domain $\Omega$. Then, instead of directly finding $\bar{u}_{1}$ and $\bar{u}_{2}$ by solving a Laplace problem with discontinuous coefficients, one solves a minimization problem on the control $g$. As already explained, this formulation is advantageous for the numerical resolution because it allows the use of nonconformal meshes on $\Omega$. Thus, the smooth extension method is a fictitious domain method in the sense that the various problems appearing in the numerical resolution process are solved on two meshes, one for $\Omega$ and an other for $\Omega_{2}$ which are not conformal.

An explicit formula is provided for the computation of the gradient of the cost function to minimize (see Theorem 2.3), which enables to treat the minimization problem with a classical descent method. The general algorithm that we use is the following: we choose an initial guess $g_{0}$ for the control term. For each iteration $k$ of the gradient algorithm, we first solve problem (6) with $g=g_{k}$ and obtain the couple ( $u_{1}^{g_{k}}, u_{2}^{g_{k}}$ ). Then, we solve the adjoint problem (13) with $g=g_{k}$ to get the adjoint variables $\left(\lambda_{1}^{g_{k}}, \lambda_{2}^{g_{k}}\right)$. The gradient $\nabla J\left(g_{k}\right)$ is computed using the explicit formula given by Theorem 2.3. Finally, the control is updated using the chosen optimization algorithm. The general formula for the update can be written

$$
g_{k+1}=g_{k}-\rho_{k} \nabla J\left(g_{k}\right),
$$

where $\rho_{k}$ is either a real positive parameter or a matrix, depending on the chosen optimization algorithm. This process is summarized in the Algorithm 1.

```
Algorithm 1 Implementation of the smooth extension method
    Choose an initial guess \(g_{0}\) for the control term.
    \(k=0\).
    while exit criteria are not satisfied do
        Compute the solution \(\left(u_{1}^{g_{k}}, u_{2}^{g_{k}}\right)\) of problem (6) with \(g=g_{k}\).
        Compute the solution \(\left(\lambda_{1}^{g_{k}}, \lambda_{2}^{g_{k}}\right)\) of the adjoint problem (13) with \(g=g_{k}\).
        Compute the gradient \(\nabla J\left(g_{k}\right)=\lambda_{2}^{g_{k}}-\lambda_{1 \mid \Omega_{2}}^{g_{k}}\).
        Update the control \(g_{k+1}=g_{k}-\rho_{k} \nabla J\left(g_{k}\right)\), with \(\rho_{k}\) depending on the chosen gradient
        method.
        Update the number of iterations \(k=k+1\).
    end while
```

In practice, the exit criteria for a descent method usually concern the norm of the difference between two successive solutions $\left\|g_{k+1}-g_{k}\right\|$ divided by the update coefficient $\rho_{k}$. The choice of the initial guess $g_{0}$ is of minor importance for the convergence of the algorithm since, as we explained in Remark 7, every extremum of the cost function $J$ corresponds to one of its minimizers.

### 3.2 Test case 1: the transmission Laplace problem in the unit square

Let $\Omega$ be the unit square of $\mathbb{R}^{2}$, divided in two pieces by a vertical segment which represents the interface $\Gamma$ (see Figure 3). We denote by $x_{\Gamma}$ the position of this interface on the $x$-axis. We aim to apply the smooth extension method to the Laplace transmission problem (1) with homogeneous Dirichlet boundary conditions, studied in subsection 2.1. Because the interface $\Gamma$ meets the boundary $\partial \Omega$ with a right-angled corner and because $f 1$ and $f_{2}$ are constants, the solution of (1), denoted by $\bar{u}$, is partially of regularity $H^{2}$, i.e. $\bar{u}_{\Omega_{1}}$ belongs


Figure 3: Two dimensional representation of the geometry for the Laplace transmission problems (1).
to $H^{2}\left(\Omega_{1}\right)$ and $\bar{u}_{\mid \Omega_{2}}$ belongs to $H^{2}\left(\Omega_{2}\right)$. Moreover, because $\Gamma$ is a straight line, we can construct a $H^{2}$-extension of $\bar{u}_{\mid \Omega_{1}}$ in the whole space $H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$. This is done using the Babič extension (see [6] or [1, Theorem 5.19]), defined for all $0<x_{\Gamma}<1$ and for all $0 \leq x, y \leq 1$ by

$$
E \bar{u}_{\mid \Omega_{1}}(x, y)= \begin{cases}\bar{u}(x, y), & \text { if } x \leq x_{\Gamma} \\ \frac{x_{\Gamma}-2}{x_{\Gamma}} \bar{u}\left(\frac{x_{\Gamma}(x-1)}{x_{\Gamma}-1}, y\right)+\frac{2}{x_{\Gamma}} \bar{u}\left(\frac{1}{2} \frac{x_{\Gamma}(x-1)}{x_{\Gamma}-1}, y\right), & \text { if } x>x_{\Gamma}\end{cases}
$$

It is straightforward to see that $E \bar{u}_{\Omega_{1}}$ and its first order partial derivatives are continuous through $\Gamma$, then $E \bar{u}_{\mid \Omega_{1}}$ well defines an extension of $\bar{u}_{\mid \Omega_{1}}$ in $H^{2}(\Omega)$. Furthermore, we can see that

$$
E \bar{u}_{\mid \Omega_{1}}(x, 0)=E \bar{u}_{\mid \Omega_{1}}(x, 1)=0 .
$$

Then, it remains to multiply this extension by a cut function such that it is zero on $\partial \Omega$. To do that we consider two real values $\alpha$ and $\beta$ in $] 0,1[$ such that $\alpha<\beta$ and we define the following cut function:

$$
\chi(x)= \begin{cases}1, & \text { if } x_{\Gamma} \leq x \leq x_{\alpha} \\ \frac{\left(2 x^{3}-3\left(x_{\alpha}+x_{\beta}\right) x^{2}+6 x_{\alpha} x_{\beta} x+x_{\beta}^{2}\left(x_{\beta}-3 x_{\alpha}\right)\right)}{\left(x_{\beta}-x_{\alpha}\right)^{3}}, & \text { if } x_{\alpha}<x \leq x_{\beta} \\ 0, & \text { if } x_{\beta}<x \leq 1\end{cases}
$$

where

$$
x_{\alpha}=x_{\Gamma}+\alpha\left(1-x_{\Gamma}\right), x_{\beta}=x_{\Gamma}+\beta\left(1-x_{\Gamma}\right)
$$

Because $\chi$ is regular enough, it follows that $E \bar{u}_{\mid \Omega_{1} \chi}$ belongs to $H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$. Thus, all conditions are satisfied so that the smooth extension method applied to this problem converges with optimal rates.

The numerical values of all parameters are $\mu_{1}=1, \mu_{2}=10, f_{1}=1, f_{2}=1$ and $x_{\Gamma}=$ 0.57. An approximation of the solution $\bar{u}$ of (1), obtained with the classical finite element method, is represented in Figure 4. Because of the jump of its gradient through $\Gamma$, the entire field $\bar{u}$ does not belong to $H^{2}(\Omega)$. As a consequence, the use of the classical finite element method to approach the solution of (1) leads to different rates of convergence when


Figure 4: Numerical solution for the Laplace problem (1) with $\mu_{1}=1, \mu_{2}=10, f_{1}=1$, $f_{2}=1$ and $x_{\Gamma}=0.57$. We compare the reference solution $\bar{u}$ obtained with the standard finite element method on a fine mesh (a) to the one obtained through the smooth extension method (c). The fields $u_{1}^{g}(\mathrm{~b})$ and $u_{2}^{g}(\mathrm{~d})$ are superposed (c) to show the continuity through the interface $\Gamma$ despite the use of non-conformal meshes. For the SEM, the mesh is $16 \times 16$.


Figure 5: Rates of convergence in space for the finite element method ( $P 1$ elements) applied to the resolution of the Laplace problem (1) with $\mu_{1}=1, \mu_{2}=10, f_{1}=1, f_{2}=1$ and $x_{\Gamma}=0.57$. The solution with a conformal mesh (a) is compare to the one with a non-conformal mesh (b)
refining the space discretization, whether the mesh of the domain fits the interface $\Gamma$ or not. Indeed, in Figure 5 we can observe that if the mesh fits the interface $\Gamma$, the rates of convergence are of order 2 for the $L^{2}$-norm and of order 1 for the $H_{0}^{1}$-norm, using $P 1$ elements. These are the classical rates of convergence with $P 1$ elements for a $H^{2}$ solution of a Laplace problem, which is not the case here. Then, this result is only due to the fact that the interface $\Gamma$ is well represented by the mesh. On the other hand, if the mesh does not fit the interface, we recover degraded rates of convergence: here we find a rate of order 1 in $L^{2}$-norm and 0.5 in $H_{0}^{1}$-norm.

For the numerical resolution with the smooth extension method, we follow Algorithm 1. In particular, we choose the initial guess $g_{0}$ to be zero and use a classical gradient descent with constant parameter $\rho$ to minimize $J$. The value of $\rho$ is chosen such that the gradient method convergences and can be different depending on the mesh size. We use unstructured meshes for both $\Omega$ and $\Omega_{2}$, such that the interface $\Gamma$ is not represented by the mesh on $\Omega$. At each iteration of the method, we solve 4 second-order boundary problems (the direct and adjoint equations) using $P 1$ elements, whose solutions enable us to compute the gradient $\nabla J\left(g_{k}\right)$. Then, we update the control such that,

$$
g_{k+1}=g_{k}-\rho \nabla J\left(g_{k}\right) .
$$

We also compute the residual error,

$$
\begin{equation*}
\frac{\left\|g_{k+1}-g_{k}\right\|}{\rho}, \tag{23}
\end{equation*}
$$

and stop the algorithm if it is smaller that a given tolerance $\varepsilon$. This method shows to easily converge to an optimal control.

The solution of (1) obtained through the minimization of the function $J$ is represented in Figure 4. It is clear than $u_{1 \mid \Omega_{1}}^{g}$ and $u_{2}^{g}$ are approximations of $\bar{u}_{\Omega_{1}}$ and $\bar{u}_{\mid \Omega_{2}}$. Moreover, superposing these two solutions, we can observe that the condition on the equality of $u_{1}^{g}$


Figure 6: Rates of convergence in space for the smooth extension method applied to the resolution of the Laplace transmission problem with homogeneous boundary conditions (1) , with $\mu_{1}=1, \mu_{2}=10, f_{1}=1, f_{2}=2$ and $x_{\Gamma}=0.67$.
and $u_{2}^{g}$ on $\Gamma$ is fulfilled. Refining the mesh, we compute the $L^{2}$ and $H_{0}^{1}$-errors between the couple ( $u_{1}^{g}, u_{2}^{g}$ ) and the reference solution obtained with the classical Finite Element method and using a fine conformal mesh. We obtain the convergence graph presented in Figure 6 where we observe than optimal rates of convergence are conserved by the smooth extension method in this case.

### 3.3 Test case 2: the transmission Laplace problem with a spherical inclusion

Let $\Omega$ be the unit square of $\mathbb{R}^{2}$ and let $\Omega_{2}$ be the disk of center $\left(x_{c}, y_{c}\right)$ and radius $r_{c}$, strictly included in $\Omega$ (see Figure 7). We aim to apply the smooth extension method to the Laplace transmission problem (17) with a strictly included domain $\Omega_{2}$ in $\Omega$, studied in subsection 2.4. Because the interface $\Gamma$ is smooth and $\Omega$ is convex, the solution of (17), denoted $\bar{u}$, is partially of regularity $H^{2}$. Moreover, we have seen in subsection 2.4 that it is possible to construct an extension of $\bar{u}_{\Omega_{1}}$ in the whole space $H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$. Then, this test case is particularly adapted to the use of the smooth extension method, which enables to recover the optimal rates of convergence in space with non-conformal meshes, whereas the classical finite element method does not.

Choosing the set of parameters, $\mu_{1}=1, \mu_{2}=10, f_{1}=1, f_{2}=2,\left(x_{c}, y_{c}\right)=(0.62,0.43)$ and $r_{c}=0.2$, we apply the smooth extension method to this transmission problem and obtain the convergence graph presented in Figure 8, where we observe, one more time, than optimal rates of convergence are preserved by the smooth extension method.

### 3.4 Test case 3: the transmission Laplace problem in a L-shape domain

Let $\Omega$ be a L-shape domain of $\mathbb{R}^{2}$ and let $\Omega_{2}$ be a quadrilateral included in $\Omega$ (see Figure 9). As before, we apply the smooth extension method to the Laplace transmission problem with


Figure 7: Two dimensional representation of the geometry for the Laplace transmission problem (17).


Figure 8: Rates of convergence in space for the smooth extension method applied to the resolution of the Laplace transmission problem with a strictly included subdomain $\Omega_{2}$ (17), for $\mu_{1}=1, \mu_{2}=10, f_{1}=1, f_{2}=2,\left(x_{c}, y_{c}\right)=(0.62,0.43)$ and $r_{c}=0.2$.


Figure 9: Two dimensional representation of the geometry for the Laplace transmission problem (1) in a L-shape domain.
homogeneous Dirichlet boundary condition (1). Because of the geometry of the problem and especially the presence of a reentrant corner, the solution is not regular and only belongs to $H_{0}^{1}(\Omega)$. Thus, there do not exists a $H^{2}$-extension from $\Omega_{1}$ to the whole domain and the method will not converge with optimal order. However, the smooth extension method can also be applied in this case and we will see that it is still better than the Finite Element method with non-conformal mesh (see Table 1). Considering the set of parameters $\mu_{1}=1, \mu_{2}=10, f_{1}=1, f_{2}=1,\left(x_{c}, y_{c}\right)=(0.62,0.7), x_{1}=0.43, x_{2}=$ $0.58,\left(x_{3}, y_{3}\right)=(0.36,0.38)$ and $\left(x_{4}, y_{4}\right)=(0.64,0.49)$, we obtain the convergence graph represented in Figure 9.

To summarise the results obtained for all test cases that we have studied, the convergence rates are presented in Table 1. As a result, the Smooth Extension method applied to the Laplace transmission problem, converges with optimal rates even with non-conformal meshes when the solution of the Laplace transmission problem is partially of regularity $H^{2}$. When the solution is less regular (because of the geometry of the domain or the regularity of the right-hand sides) the Smooth Extension method can also be applied and showed to converge with the same rates than the Finite Element method with conformal mesh. In the next section we will show that this method can be extended to other kind of coupled problems and enables to treat the numerical resolution of more general transmission problems with non-conformal meshes.

## 4 Extension to other coupled problems

In this section we extend the smooth extension method to two other coupled problems: the Stokes transmission problem in subsection 4.1 and a fluid-structure interaction problem in subsection 4.2 .


Figure 10: Rates of convergence in space for the smooth extension method applied to the resolution of the Laplace transmission problem (1) in a L-shape domain.

| Test case | Method | Mesh conformity | conv. rate in $H_{0}^{1}$-norm | conv. rate in $L^{2}$-norm |
| :--- | :---: | :---: | :---: | :---: |
| 1 | FEM | Conformal | 1.11 | 2.08 |
| 1 | FEM | Non-conformal | 0.58 | 1.28 |
| 1 | SEM | Non-conformal | $1.08\left(\right.$ in $\left.\Omega_{1}\right) / 1.14\left(\right.$ in $\left.\Omega_{2}\right)$ | $2.08\left(\right.$ in $\left.\Omega_{1}\right) / 2.1\left(\right.$ in $\left.\Omega_{2}\right)$ |
| 2 | FEM | Conformal | 1.06 | 2.13 |
| 2 | FEM | Non-conformal | 0.58 | 1.04 |
| 2 | SEM | Non-conformal | $1.12\left(\operatorname{in} \Omega_{1}\right) / 1.21\left(\right.$ in $\left.\Omega_{2}\right)$ | $2.1\left(\right.$ in $\left.\Omega_{1}\right) / 2.17\left(\right.$ in $\left.\Omega_{2}\right)$ |
| 3 | FEM | Conformal | 0.96 | 1.7 |
| 3 | FEM | Non-conformal | 0.76 | 1.43 |
| 3 | SEM | Non-conformal | $0.93\left(\right.$ in $\left.\Omega_{1}\right) / 1.2\left(\right.$ in $\left.\Omega_{2}\right)$ | $1.72\left(\right.$ in $\left.\Omega_{1}\right) / 2.08\left(\right.$ in $\left.\Omega_{2}\right)$ |

Table 1: Comparison of the rates of convergence between the finite element method (FEM) with a conformal mesh, the FEM with a non-conformal mesh and the smooth extension method (SEM) with a non-conformal mesh for different test cases.

### 4.1 The Stokes transmission problem

Let $n>0$ and $\Omega$ be a domain of $\mathbb{R}^{n}$ that satisfies the following set of hypotheses:
i) Domain $\Omega$ is a bounded connected Lipschitz domain of $\mathbb{R}^{n}$.
ii) Domain $\Omega$ is divided into two connected Lipschitz subdomains, $\Omega_{1}$ and $\Omega_{2}$.
iii) The interface $\Gamma=\partial \Omega_{1} \cap \partial \Omega_{2}$ is not empty.
iv) The remaining boundaries $\Gamma_{1}=\partial \Omega_{1} \backslash \Gamma$ and $\Gamma_{2}=\partial \Omega_{2} \backslash \Gamma$ are not empty.

Consider two real coefficients, $\mu_{1}$ and $\mu_{2}$, and two external forces, $f_{1}$ in $\left(L^{2}\left(\Omega_{1}\right)\right)^{n}$ and $f_{2}$ in $\left(L^{2}\left(\Omega_{2}\right)\right)^{n}$. Then, the Stokes transmission problem writes:

$$
\left\{\begin{array}{rlrl}
\text { find } u_{1}: \Omega_{1} \rightarrow \mathbb{R}^{n}, p_{1}: \Omega_{1} \rightarrow \mathbb{R}, u_{2}: \Omega_{2} \rightarrow \mathbb{R}^{n} & \text { and } p_{2}: \Omega_{2} & \rightarrow \mathbb{R} \text { such that }  \tag{24a}\\
-\mu_{1} \Delta u_{1}+\nabla p_{1} & =f_{1}, & & \text { in } \Omega_{1}, \\
\operatorname{div}\left(u_{1}\right) & =0, & & \text { in } \Omega_{1}, \\
u_{1} & =0, & & \text { on } \Gamma_{1}, \\
-\mu_{2} \Delta u_{2}+\nabla p_{2} & =f_{2}, & & \text { in } \Omega_{2}, \\
\operatorname{div}\left(u_{2}\right) & =0, & & \text { in } \Omega_{2}, \\
u_{2} & =0, & & \text { on } \Gamma_{2}, \\
u_{1} & =u_{2}, & & \text { on } \Gamma, \\
\left(\mu_{1} \nabla u_{1}-p_{1} I\right) \cdot n_{1} & =-\left(\mu_{2} \nabla u_{2}-p_{2} I\right) \cdot n_{2}, & & \text { on } \Gamma .
\end{array}\right.
$$

Equations (24a) and (24b) are two sets of Stokes equations coupled at the interface $\Gamma$ with the coupling conditions $(24 \mathrm{c})$. These conditions represent the continuity of the fluid velocity and the continuity of the constraints applied by the fluid on $\Gamma$. The vectors $n_{1}$ and $n_{2}$ still denote the unit exterior normal vector of $\Omega_{1}$ and $\Omega_{2}$. Of course, problem (24) is equivalent to the Stokes problem with discontinuous viscosity and external force,

$$
\left\{\begin{align*}
& \text { find } u: \Omega \rightarrow \mathbb{R}^{n} \text { and } p: \Omega \rightarrow \mathbb{R} \text { such that }  \tag{25}\\
&-\operatorname{div}(\mu \nabla u-p I)=f
\end{align*} \quad \text { in } \Omega,\right.
$$

where $\mu$ and $f$ are defined by

$$
\begin{aligned}
\mu & = \begin{cases}\mu_{1} & \text { in } \Omega_{1} \\
\mu_{2} & \text { in } \Omega_{2}\end{cases} \\
f & = \begin{cases}f_{1} & \text { in } \Omega_{1} \\
f_{2} & \text { in } \Omega_{2}\end{cases}
\end{aligned}
$$

We can define a weak formulation of problem (24). Let us introduce the functional spaces

$$
\begin{aligned}
W_{1} & =\left\{v_{1} \in\left(H^{1}\left(\Omega_{1}\right)\right)^{n} ; v_{1 \mid \Gamma_{1}}=0\right\} \\
W_{2} & =\left\{v_{2} \in\left(H^{1}\left(\Omega_{2}\right)\right)^{n} ; v_{2 \mid \Gamma_{2}}=0\right\} \\
\mathcal{W} & =\left\{\left(v_{1}, v_{2}\right) \in W_{1} \times W_{2} ; v_{1 \mid \Gamma}=v_{2 \mid \Gamma}\right\} \\
Q & =\left\{\left(p_{1}, p_{2}\right) \in L^{2}\left(\Omega_{1}\right) \times L^{2}\left(\Omega_{2}\right) ; \frac{1}{|\Omega|}\left(\int_{\Omega_{1}} p_{1}+\int_{\Omega_{2}} p_{2}\right)=0\right\}, \\
\left(H_{\operatorname{div}}(X)\right)^{n} & =\left\{\sigma \in\left(L^{2}(X)\right)^{n \times n} ; \operatorname{div}(\sigma) \in\left(L^{2}(X)\right)^{n}\right\}
\end{aligned}
$$

where $X$ stands for either $\Omega_{1}$ or $\Omega_{2}$. Then, the weak problem associated to problem (24) writes:

$$
\left\{\begin{align*}
& \text { find }\left(u_{1}, u_{2}\right) \in \mathcal{W} \text {, and }\left(p_{1}, p_{2}\right) \in Q \text { such that }  \tag{26}\\
& \mu_{1} \int_{\Omega_{1}} \nabla u_{1}: \nabla v_{1}+\mu_{2} \int_{\Omega_{2}} \nabla u_{2}: \nabla v_{2}-\int_{\Omega_{1}} p_{1} \operatorname{div}\left(v_{1}\right)-\int_{\Omega_{2}} p_{2} \operatorname{div}\left(v_{2}\right)=\int_{\Omega_{1}} f_{1} \cdot v_{1}+\int_{\Omega_{2}} f_{2} \cdot v_{2}, \\
& \forall\left(v_{1}, v_{2}\right) \in \mathcal{W} \\
& \int_{\Omega_{1}} q_{1} \operatorname{div}\left(u_{1}\right)+\int_{\Omega_{2}} q_{2} \operatorname{div}\left(u_{2}\right)=0, \forall\left(q_{1}, q_{2}\right) \in Q
\end{align*}\right.
$$

The well-posedness of problem (25) in $\left(H_{0}^{1}(\Omega)\right)^{n} \times L_{0}^{2}(\Omega)$ is a particular case of [10, Theorem IV.8.1], where $\mu$ (positive) needs to be in $L^{\infty}(\Omega)$ and $f$ belongs to $\left(L^{2}(\Omega)\right)^{n}$. It follows that problem (26) is well-posed and we denote by $\left(\left(\bar{w}_{1}, \bar{w}_{2}\right),\left(\bar{p}_{1}, \bar{p}_{2}\right)\right)$ its unique solution. In particular, $\sigma_{1}=\mu_{1} \nabla u_{1}-p_{1} I$ belongs to $\left(H_{\text {div }}\left(\Omega_{1}\right)\right)^{n}$, $\sigma_{2}=\mu_{2} \nabla u_{2}-p_{2} I$ belongs to $\left(H_{\mathrm{div}}\left(\Omega_{2}\right)\right)^{n}$ and we are able to give a weak sense to the second transmission condition on $\Gamma$ in (24c). Let $X$ stand for either $\Omega_{1}$ or $\Omega_{2}$ and let $\eta$ be the unit exterior normal vector to $X$. Let $\Upsilon=\left(H_{00}^{1 / 2}(\Gamma)\right)^{n}$ be the image of $\left(H_{\partial X \backslash \Gamma}^{1}(X)\right)^{n}$ by the trace operator on $\Gamma$, i.e. the space of functions in $\left(H^{1 / 2}(\Gamma)\right)^{n}$ whose extension by zero on $\partial X \backslash \Gamma$ belongs to $\left(H^{1 / 2}(\partial X)\right)^{n}$. Then, for all $\sigma$ in $\left(H_{\text {div }}(X)\right)^{n}$, we have the following Stokes formula:

$$
\begin{equation*}
\int_{X} \sigma: \nabla v+\int_{X} \operatorname{div}(\sigma) \cdot v=\left\langle\gamma_{\eta}(\sigma), v\right\rangle_{\Upsilon^{\prime}, \Upsilon}, \quad \forall v \in\left(H_{\partial X \backslash \Gamma}^{1}(X)\right)^{n} . \tag{27}
\end{equation*}
$$

where $\Upsilon^{\prime}$ is the dual space of $\Upsilon$. Then, the second transmission condition in (24c) is satisfied in the following sense:

$$
\begin{equation*}
\left\langle\gamma_{n_{1}}\left(\sigma_{1}\right), v\right\rangle_{\Upsilon^{\prime}, \Upsilon}=-\left\langle\gamma_{n_{2}}\left(\sigma_{2}\right), v\right\rangle_{\Upsilon^{\prime}, \Upsilon}, \quad \forall v \in \Upsilon . \tag{28}
\end{equation*}
$$

Now, we present the smooth extension method applied to problem (24). Formally, it writes: find $g$ such that the solution of the following problem,

$$
\left\{\begin{array}{rlrlrl}
\text { find } u_{1}: \Omega_{1} \rightarrow \mathbb{R}^{n}, p_{1}: \Omega_{1} \rightarrow \mathbb{R}, u_{2}: \Omega_{2} \rightarrow \mathbb{R}^{n} & \text { and } p_{2}: \Omega_{2} & \rightarrow \mathbb{R} \text { such that }  \tag{29a}\\
-\mu_{1} \Delta u_{1}+\nabla p_{1} & =\bar{f}_{1}+\bar{g}^{\Omega}, & & \text { in } \Omega, \\
\operatorname{div}\left(u_{1}\right) & =0, & & \text { in } \Omega, \\
u_{1} & =0, & & \text { in } \partial \Omega, \\
-\mu_{2} \Delta u_{2}+\nabla p_{2} & =f_{2}, & & \text { in } \Omega_{2}, \\
\operatorname{div}\left(u_{2}\right) & =0, & & \text { in } \Omega_{2}, \\
u_{2} & =0, & & \text { on } \Gamma_{2}, \\
\left(\mu_{2} \nabla u_{2}-p_{2} I\right) \cdot n_{2} & =\left(\mu_{1} \nabla u_{1}-p_{1} I\right) \cdot n_{2}, & & \text { on } \Gamma,
\end{array}\right.
$$

satisfies the equality

$$
u_{1}=u_{2}, \quad \text { on } \Gamma
$$

Let $v$ be in $(\mathcal{D}(\Omega))^{n}$ and suppose that $u_{1}, p_{1}, u_{2}$ and $p_{2}$ are sufficiently regular. Moreover, we assume for the moment that $g$ belongs to $\left(L^{2}\left(\Omega_{2}\right)\right)^{n}$ in order to do formal computations. Formally, we multiply equation (29a) by $v$ and equation (29b) by $v_{\mid \Omega_{2}}$ and integrate respectively over $\Omega$ and $\Omega_{2}$. After an integration by part and using the Neumann condition on $\Gamma$, we find

$$
\begin{aligned}
\mu_{1} \int_{\Omega} \nabla u_{1}: \nabla v-\int_{\Omega} p_{1} \operatorname{div}(v) & =\int_{\Omega_{1}} f_{1} \cdot v_{\mid \Omega_{1}}+\int_{\Omega_{2}} g \cdot v_{\mid \Omega_{2}} \\
\mu_{2} \int_{\Omega_{2}} \nabla u_{2}: \nabla v_{\mid \Omega_{2}}-\int_{\Omega_{2}} p_{2} \operatorname{div}\left(v_{\mid \Omega_{2}}\right) & =\int_{\Omega_{2}} f_{2} \cdot v_{\mid \Omega_{2}}+\int_{\Gamma}\left(\left(\mu_{1} \nabla u_{1}-p_{1} I\right) n_{2}\right) \cdot v .
\end{aligned}
$$

Moreover, we remark that

$$
\begin{aligned}
\int_{\Gamma}\left(\left(\mu_{1} \nabla u_{1}-p_{1} I\right) n_{2}\right) \cdot v= & \int_{\Omega_{2}}\left(\mu_{1} \Delta u_{1}-\nabla p_{1}\right) \cdot v_{\mid \Omega_{2}}+\mu_{1} \int_{\Omega_{2}} \nabla u_{1}: \nabla v_{\mid \Omega_{2}} \\
& -\int_{\Omega_{2}} p_{1} \operatorname{div}\left(v_{\mid \Omega_{2}}\right), \\
= & -\int_{\Omega_{2}} g \cdot v_{\mid \Omega_{2}}+\mu_{1} \int_{\Omega_{2}} \nabla u_{1}: \nabla v_{\mid \Omega_{2}} \\
& -\int_{\Omega_{2}} p_{1} \operatorname{div}\left(v_{\mid \Omega_{2}}\right) .
\end{aligned}
$$

Thus, we define the weak formulation of the smooth extension problem associated to problem (29), which makes sense for $g \in W_{2}^{\prime}$, as the problem of finding a suitable control $g$ in $W_{2}^{\prime}$, such that the solution of the following problem,

$$
\left\{\begin{array}{cl}
\text { find }\left(u_{1}, p_{1}\right) \in\left(H_{0}^{1}(\Omega)\right)^{n} \times L_{0}^{2}(\Omega) \text { and }\left(u_{2}, p_{2}\right) \in W_{2} \times L^{2}\left(\Omega_{2}\right) \text { such that } \\
\mu_{1} \int_{\Omega} \nabla u_{1}: \nabla v_{1}-\int_{\Omega} p_{1} \operatorname{div}\left(v_{1}\right) & \forall \int_{\Omega_{1}} f_{1} \cdot v_{1 \mid \Omega_{1}}+\left\langle g, v_{1 \mid \Omega_{2}}\right\rangle_{W_{2}^{\prime}, W_{2}}, \\
\int_{\Omega} q_{1} \operatorname{div}\left(u_{1}\right)=0, & \forall v_{1} \in\left(H_{0}^{1}(\Omega)\right)^{n}, \\
=\int_{\Omega_{2}} f_{2} \cdot v_{2}-\left\langle g, v_{2}\right\rangle_{W_{2}^{\prime}, W_{2}}^{2}(\Omega), \\
\quad+\mu_{1} \int_{\Omega_{2}} \nabla u_{1}: \nabla v_{2}-\int_{\Omega_{2}} p_{1} \operatorname{div}\left(v_{2}\right), & \forall v_{2} \in W_{2},  \tag{30}\\
\mu_{2} \int_{\Omega_{2}} \nabla u_{2}: \nabla v_{2}-\int_{\Omega_{2}} p_{2} \operatorname{div}\left(v_{2}\right) & \forall q_{2} \in L^{2}\left(\Omega_{2}\right),
\end{array}\right.
$$

satisfies the equality

$$
\begin{equation*}
u_{1}=u_{2}, \quad \text { on } \Gamma . \tag{31}
\end{equation*}
$$

For every $f_{1}$ in $\left(L^{2}\left(\Omega_{1}\right)\right)^{n}$, every $f_{2}$ in $\left(L^{2}\left(\Omega_{2}\right)\right)^{n}$ and every $g$ in $W_{2}^{\prime}$, the two subproblems in problem (30) admit a unique solution. They are denoted by $\left(w_{1}^{g}, p_{1}^{g}\right)$ and $\left(w_{2}^{g}, p_{2}^{g}\right)$. These are well-known results on Stokes equations; we refer to [10] for details. In the following theorem, we state the existence of a control $g$ such that the solution of $(30)$ satisfies the condition (31) and explain how to recover the solution of the initial Stokes transmission problem (26).

Theorem 4.1. Let $\Omega$ be a domain that satisfies Assumption $\left(H_{3}\right)$. Consider $f_{1}$ in $\left(L^{2}\left(\Omega_{1}\right)\right)^{n}$ and $f_{2}$ in $\left(L^{2}\left(\Omega_{2}\right)\right)^{n}$. Then, there exists a function $g$ in $W_{2}^{\prime}$ such that the solution $\left(\left(w_{1}^{g}, p_{1}^{g}\right),\left(w_{2}^{g}, p_{2}^{g}\right)\right)$ of (30) satisfies (31). Moreover, we can recover the solution of the Stokes transmission problem (26):

$$
\begin{aligned}
\left(\bar{w}_{1}, \bar{w}_{2}\right) & =\left(w_{1 \mid \Omega_{1}}^{g}, w_{2}^{g}\right) \\
\left(\bar{p}_{1}, \bar{p}_{2}\right) & =\left(p_{1 \mid \Omega_{1}}^{g}-C, p_{2}^{g}-C\right)
\end{aligned}
$$

where

$$
C=\frac{1}{|\Omega|}\left(\int_{\Omega_{1}} p_{1 \mid \Omega_{1}}^{g}+\int_{\Omega_{2}} p_{2}^{g}\right) .
$$

The proof of Theorem 4.1 is similar to the proof of Theorem 2.1. It relies on the construction of extensions for the velocity $\bar{w}_{1}$ and the pressure $\bar{p}_{1}$ in the whole domain $\Omega$. Details are given in Appendix A.

As before, the problem of finding a suitable control $g$ such that the solution of problem (30) satisfies (31), can be formulated as an optimization problem. In practice, it is this minimization problem which is solved in order to obtain a suitable control and recover the solution of the transmission problem (24). The cost function to consider, that we denote by $\tilde{J}$, is now defined from $W_{2}^{\prime}$ to $\mathbb{R}^{+}$with the formula

$$
\begin{equation*}
\tilde{J}(g)=\frac{1}{2} \int_{\Gamma}\left|w_{1}^{g}-w_{2}^{g}\right|^{2}, \tag{32}
\end{equation*}
$$

where $w_{1}^{g}$ and $w_{2}^{g}$ are the velocities of the fluid, solutions of problem (30). Yet, the minimization of this cost function is equivalent to the minimization of the real-valued function

$$
\begin{aligned}
\left(H_{0}^{1}(\Omega)\right)^{n} \times W_{2} & \rightarrow \mathbb{R}^{+} \\
\left(v_{1}, v_{2}\right) & \mapsto \frac{1}{2} \int_{\Gamma}\left|v_{\left.1\right|_{\Gamma}}-v_{\left.2\right|_{\Gamma}}\right|^{2},
\end{aligned}
$$

under the constraint that $v_{1}$ and $v_{2}$ are the velocities that solve the problem (30). Then, to this constrained optimization problem we associate the following Lagrangian function defined from

$$
W_{2}^{\prime} \times\left[\left(\left(H_{0}^{1}(\Omega)\right)^{n} \times L_{0}^{2}(\Omega)\right) \times\left(W_{2} \times L^{2}\left(\Omega_{2}\right)\right)\right] \times\left[\left(\left(H_{0}^{1}(\Omega)\right)^{n} \times L_{0}^{2}(\Omega)\right) \times\left(W_{2} \times L^{2}\left(\Omega_{2}\right)\right)\right]
$$

to $\mathbb{R}$ by

$$
\begin{align*}
\tilde{\mathcal{L}}(g,( & \left.\left.\left(u_{1}, p_{1}\right),\left(u_{2}, p_{2}\right)\right),\left(\left(\lambda_{1}, \pi_{1}\right),\left(\lambda_{2}, \pi_{2}\right)\right)\right)= \\
& \frac{1}{2} \int_{\Gamma}\left|u_{1}-u_{2}\right|^{2}+\mu_{1} \int_{\Omega} \nabla u_{1}: \nabla \lambda_{1}-\mu_{1} \int_{\Omega_{2}} \nabla u_{1}: \nabla \lambda_{2}+\mu_{2} \int_{\Omega_{2}} \nabla u_{2}: \nabla \lambda_{2} \\
& -\int_{\Omega} p_{1} \operatorname{div}\left(\lambda_{1}\right)-\int_{\Omega} \pi_{1} \operatorname{div}\left(u_{1}\right)-\int_{\Omega_{2}} p_{2} \operatorname{div}\left(\lambda_{2}\right)-\int_{\Omega_{2}} \pi_{2} \operatorname{div}\left(u_{2}\right)+\int_{\Omega_{2}} p_{1} \operatorname{div}\left(\lambda_{2}\right) \\
& -\int_{\Omega_{1}} f_{1} \cdot \lambda_{1 \mid \Omega_{1}}-\left\langle g, \lambda_{1 \mid \Omega_{2}}\right\rangle_{W_{2}^{\prime}, W_{2}}-\int_{\Omega_{2}} f_{2} \cdot \lambda_{2}+\left\langle g, \lambda_{2}\right\rangle_{W_{2}^{\prime}, W_{2}} \tag{33}
\end{align*}
$$

Again, this Lagrangian function enables to compute the gradient of $\tilde{J}$ and it is possible to show the equivalence between the minimization of $\tilde{J}$ and the research of a suitable control such that the solution of (30) satisfies the condition (31). For that matter, we introduce the adjoint problem of (30),

$$
\left\{\begin{array}{rlrl}
\text { find }\left(\lambda_{1}, \pi_{1}\right) \in\left(H_{0}^{1}(\Omega)\right)^{n} \times L_{0}^{2}(\Omega) \text { and }\left(\lambda_{2}, \pi_{2}\right) \in W_{2} \times L^{2}\left(\Omega_{2}\right) \text { such that } \\
\mu_{2} \int_{\Omega_{2}} \nabla \lambda_{2}: \nabla v_{2}-\int_{\Omega_{2}} \pi_{2} \operatorname{div}\left(v_{2}\right)= & \int_{\Gamma}\left(w_{1}^{g}-w_{2}^{g}\right) \cdot v_{2}, & \forall v_{2} \in W_{2} \\
\int_{\Omega_{2}} q_{2} \operatorname{div}\left(\lambda_{2}\right)= & 0, & \forall q_{2} \in L^{2}\left(\Omega_{2}\right) \\
\mu_{1} \int_{\Omega} \nabla \lambda_{1}: \nabla v_{1}-\int_{\Omega} \pi_{1} \operatorname{div}\left(v_{1}\right)= & -\int_{\Gamma}\left(w_{1}^{g}-w_{2}^{g}\right) \cdot v_{1} & \\
& +\mu_{1} \int_{\Omega_{2}} \nabla \lambda_{2}: \nabla v_{1}, & \forall v_{1} \in\left(H_{0}^{1}(\Omega)\right)^{n}, \\
\int_{\Omega} q_{1} \operatorname{div}\left(\lambda_{1}\right)= & \int_{\Omega_{2}} q_{1} \operatorname{div}\left(\lambda_{2}\right), & \forall q_{1} \in L_{0}^{2}(\Omega) \tag{34}
\end{array}\right.
$$

Problem (34) admits a unique solution denoted by $\left(\left(\omega_{1}^{g}, \pi_{1}^{g}\right),\left(\omega_{2}^{g}, \pi_{2}^{g}\right)\right)$ (see Appendix A). Then, the existence and the characterization of the gradient of $\tilde{J}$ is given in the following theorem,


Figure 11: Two dimensional representation of the geometry for the Stokes transmission problem.

Theorem 4.2. The mapping $g \in W_{2}^{\prime} \mapsto \tilde{J}(g) \in \mathbb{R}$ is differentiable and its gradient $\nabla \tilde{J}(g) \in$ $W_{2}^{\prime \prime}$ is given by,

$$
\begin{equation*}
\langle\nabla \tilde{J}(g), \delta g\rangle_{W_{2}^{\prime \prime}, W_{2}^{\prime}}=\left\langle\delta g, \omega_{2}^{g}-\omega_{1 \mid \Omega_{2}}^{g}\right\rangle_{W_{2}^{\prime}, W_{2}}, \quad \forall \delta g \in W_{2}^{\prime}, \tag{35}
\end{equation*}
$$

where $\omega_{1}^{g}$ and $\omega_{2}^{g}$ are the unique velocities that verify the adjoint problem (34).
Moreover, the equivalence between the smooth extension problem and its formulation as a minimization problem can now be proved,

Theorem 4.3. A control $g$ in $W_{2}^{\prime}$ is a minimizer of $\tilde{J}$ if and only if the solution of (30) satisfies the condition (31).

For reasons of clarity, the proofs of Theorem 4.2 and Theorem 4.3 are done in Appendix A.

As for the Laplace transmission problem, the formulation of the Stokes transmission problem (24) as a control problem and as a minimization problem on the function $\tilde{J}$ enables us to numerically solve these equations with a fictitious domain approach. Here, we apply this method to the numerical simulation of a two-layer Stokes fluid in the unit square of $\mathbb{R}^{2}$. The geometric settings are represented in Figure 11: $\Omega_{1}$ and $\Omega_{2}$ are separated by a smooth curve which encounters the boundary $\partial \Omega$ at two points, $\left(1, \alpha_{\Gamma}\right)$ and $\left(\beta_{\Gamma}, 1\right)$, where $\alpha_{\Gamma}$ and $\beta_{\Gamma}$ are two positive constants. The interface $\Gamma$ is defined such that a point $(x, y)$ in $\left[\beta_{\Gamma}, 1\right] \times\left[\alpha_{\Gamma}, 1\right]$ belongs to $\Gamma$ if and only if

$$
1-x \leq \beta_{\Gamma} \sqrt{\frac{y-\alpha_{\Gamma}}{y}}
$$

We consider homogeneous Dirichlet boundary conditions on the left and right boundaries and a pressure drop of 1 between the top and the bottom boundaries. No external force are considered and constant viscosities $\mu_{1}=1$ and $\mu_{2}=10$ are chosen. For the interface $\Gamma$, we choose $\alpha_{\Gamma}=0.25$ and $\beta_{\Gamma}=0.5$.

The solution given by the smooth extension method is obtained following Algorithm 1 and using the gradient method with fixed parameter $\rho$. This parameter is chosen such that the gradient method convergences and the stopping criteria is defined in (23). The resolution of the four different Stokes problems appearing in the minimization process is
done using Mini elements. To compare the smooth extension method to the classical finite element method, we use meshes that have about the same numbers of cells: $26 \times 26$.

Results are shown in Figure 12, where both the velocity and the pressure of the fluid are plotted. We compare the solution obtained through the standard finite element method with a mesh which is actually conform with the interface $\Gamma$, to the one produced by the smooth extension method with unstructured meshes. We observe that both solutions are similar at the difference that the Smooth Extension method well represents the physical jump in pressure through the interface $\Gamma$, while the finite element method with Mini elements does not. Moreover, the former also has the advantage to be computed with a non-conformal mesh.

### 4.2 A fluid-structure interaction problem

Now, we are interested in the resolution of a fluid-structure problem where the fluid is modeled by the Stokes equations and the structure by the stationary equations of linear elasticity. The unknowns for these two systems of equations are the fluid velocity and pressure as well as the displacement of the structure from its reference configuration. The fluid problem will be set in Eulerian coordinates, i.e. in the current configuration, whereas the elastic equations will be written in Lagrangian coordinates, i.e. in the reference configuration. This amounts to consider two configurations for both problems, supposing that, at each time, there exists a smooth enough mapping between the two configurations to ensure that all boundaries in the current configuration are sufficiently regular.

Let $n>0$ and $\Omega$ be a domain of $\mathbb{R}^{n}$ that satisfies the following set of hypotheses:
i) Domain $\Omega$ is a bounded connected Lipschitz domain of $\mathbb{R}^{n}$.
ii) Domain $\Omega$ is divided in two connected Lipschitz subdomains, $\Omega_{f}$ the fluid subdomain and $\Omega_{s}$ the solid subdomain.
iii) The interface $\Gamma=\partial \Omega_{f} \cap \partial \Omega_{s}$ is not empty.
iv) The remaining boundaries $\Gamma_{f}=\partial \Omega_{f} \backslash \Gamma$ and $\Gamma_{s}=\partial \Omega_{s} \backslash \Gamma$ are not empty.

We assume that, at each time $t \geq 0$, there exists a deformation $\Phi_{t}$, i.e. a smooth enough injective and orientation-preserving mapping, defined from $\Omega$ to $\mathbb{R}^{n}$, such that the current fluid configuration $\Phi_{t}\left(\Omega_{f}\right)$ and the current solid configuration $\Phi_{t}\left(\Omega_{s}\right)$ also are Lipschitz subdomains of $\Omega$. Moreover, $\Phi_{t}\left(\Omega_{f}\right)$ and $\Phi_{t}\left(\Omega_{s}\right)$ should be connected. We should precise here that the mapping $\Phi_{t}$ directly depends on the displacement of the structure at time $t$. For all $x$ in $\Omega_{s}$, this mapping writes

$$
\Phi_{t}(x)=x+d_{s}(t)(x)
$$

and we can easily extend $\Phi_{t}$ in the whole domain $\Omega$. The existence of a smooth transformation $\Phi_{t}$ is rather complicated to prove and is out of the scope of this study. However, for the numerical resolution of this fluid-structure interaction problem, it is possible to construct $\Phi_{t}$ if the deformation of the structure is reasonable, i.e. if the structure does not enter in contact with itself or the boundary $\partial \Omega \backslash \Gamma_{s}$ and if the mesh which represents the domain $\Omega_{s}$ is admissible (no overlapping cell).

Then, at each time $t \geq 0$, the subdomain $\Phi_{t}\left(\Omega_{f}\right)$ is filled with an incompressible Newtonian fluid whose velocity $u_{f}(t): \Phi_{t}\left(\Omega_{f}\right) \rightarrow \mathbb{R}^{n}$ and pressure $p_{f}(t): \Phi_{t}\left(\Omega_{f}\right) \rightarrow \mathbb{R}$ satisfy the Stokes equations in conservative form

$$
\begin{array}{rlrl}
-\operatorname{div}\left(\sigma_{f}\left(u_{f}(t), p_{f}(t)\right)\right) & =f_{f}(t) & & \text { in } \\
\Phi_{t}\left(\Omega_{f}\right), \\
\operatorname{div}\left(u_{f}(t)\right) & =0 & & \text { in } \\
\Phi_{t}\left(\Omega_{f}\right), \\
u_{f}(t) & =0 & & \text { on } \\
\Phi_{t}\left(\Gamma_{f}\right),
\end{array}
$$



Figure 12: Representation of the solution for the Stokes transmission problem (24) with $\mu_{1}=1, \mu_{2}=10, f_{1}=(1,0), f_{2}=(1,0), \alpha_{\Gamma}=0.25$ and $\beta_{\Gamma}=0.5$. One the top, we compare the magnitude of the fluid velocity obtained with a standard finite element method (a) and with the smooth extension method (b). On the bottom, we compare the pressure of the fluid computed with the classical finite element method (c) and the smooth extension $\operatorname{method}(d)$. The interface $\Gamma$ is highlighted with a white curve.
where $\sigma_{f}$ is the fluid tensor defined for all $u: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ and all $p: \mathbb{R}^{n} \rightarrow \mathbb{R}$ by

$$
\begin{aligned}
\sigma_{f}(u, p) & =2 \mu_{f} D(u)-p I, \\
D(u) & =\frac{1}{2}\left(\nabla u+\nabla u^{T}\right),
\end{aligned}
$$

the constant $\mu_{f}$ is the viscosity of the fluid and $f_{f}(t): \Phi_{t}\left(\Omega_{f}\right) \rightarrow \mathbb{R}^{n}$ is the external force applied to the fluid at time $t$. Inside the fluid lies an elastic medium, whose displacement at time $t, d_{s}(t): \Omega_{s} \rightarrow \mathbb{R}^{n}$, verifies the following equations of linear elasticity written in the reference solid configuration $\Omega_{s}$

$$
\begin{array}{rllll}
-\operatorname{div}\left(\sigma_{s}\left(d_{s}(t)\right)\right) & =f_{s}(t) & & \text { in } & \Omega_{s}, \\
d_{s}(t) & =0 & & \text { on } & \Gamma_{s},
\end{array}
$$

where $\sigma_{s}$ is the solid tensor defined for all $u: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ by

$$
\sigma_{s}(u)=2 \mu_{s} D(u)+\lambda_{s} \operatorname{div}(u) I
$$

the two positive constants $\mu_{s}$ and $\lambda_{s}$ are the Lamé coefficients and $f_{s}(t): \Omega_{s} \rightarrow \mathbb{R}^{n}$ is the external force applied to the structure at time $t$. To complete this system of equations, we consider at each time $t$ the coupling conditions that correspond to the continuity of the velocities and the normal constraints through the fluid-structure interface in the reference configuration $\Gamma$. For that matter, we introduce the fluid velocity and pressure written in the reference fluid configuration, denoted by $w_{f}$ and $q_{f}$, and defined at time $t$ by

$$
w_{f}(t)=u_{f}(t) \circ \Phi_{t}, \quad \text { and } q_{f}(t)=p_{f}(t) \circ \Phi_{t} .
$$

Moreover, we introduce the fluid stress tensor written in the fluid reference configuration, denoted by $\Pi_{f}$, and defined at time $t$ by

$$
\Pi_{f}\left(w_{f}(t), q_{f}(t)\right)=\mu_{f}\left(\nabla w_{f}(t) F\left(d_{s}(t)\right)+\nabla w_{f}(t)^{T} F\left(d_{s}(t)\right)^{T}\right)-q_{f}(t) G\left(d_{s}(t)\right), \quad \text { in } \Omega_{f}
$$

where $F\left(d_{s}(t)\right)$ and $G\left(d_{s}(t)\right)$ are the following matrices:

$$
\begin{align*}
& F\left(d_{s}(t)\right)=\left(\nabla\left(\Phi\left(d_{s}(t)\right)\right)\right)^{-1} \operatorname{cof}\left(\nabla\left(\Phi\left(d_{s}(t)\right)\right)\right),  \tag{36}\\
& G\left(d_{s}(t)\right)=\operatorname{cof}\left(\nabla\left(\Phi\left(d_{s}(t)\right)\right)\right) .
\end{align*}
$$

Thus the transmission conditions write

$$
\begin{array}{rlrl}
\frac{\partial d_{s}}{\partial t}(t) & =u_{f}(t) \circ \Phi_{t} & & \text { on }  \tag{37}\\
\sigma_{s}\left(d_{s}(t)\right) n_{s} & =\Pi_{f}\left(w_{f}(t), q_{f}(t)\right) n_{s} & & \text { on } \\
\Gamma,
\end{array}
$$

where the vector $n_{s}$ denotes the exterior unit normal vector to $\partial \Omega_{s}$. Similarly we denote $n_{d}$ the exterior unit normal vector to $\partial \Omega_{f}$. Furthermore, we suppose that the structure is at rest initially, i.e. that $d_{s}(0)=0$.

For the purpose of the numerical resolution we consider a discretization of $\mathbb{R}^{+}$for the time variable. Let $\delta t>0$ be the step size. We construct a sequence $\left(t_{k}\right)_{k \in \mathbb{R}^{+}}$such that $t_{0}=0$ and $t_{k+1}=t_{k}+\delta t$ for $k>0$. Thus, we define the time discretizations of $u_{f}, p_{f}$ and $d_{s}$ such that, for all $k \geq 0$,

$$
\begin{aligned}
u_{f}^{k} & =u_{f}\left(t_{k}\right), \\
p_{f}^{k} & =p_{f}\left(t_{k}\right), \\
d_{s}^{k} & =d_{s}\left(t_{k}\right) .
\end{aligned}
$$

In addition, we also define $f_{f}^{k}=f_{f}\left(t_{k}\right)$ and $f_{s}^{k}=f_{s}\left(t_{k}\right)$ for all $k \geq 0$. The discretization of the first coupling condition (37) is obtained using the implicit Euler scheme:

$$
\begin{aligned}
d_{s}^{k+1} & =d_{s}^{k}+\delta t u_{f}^{k+1}, \quad \forall k \geq 0 \\
d_{s}^{0} & =\delta t u_{f}^{0}
\end{aligned}
$$

At time $t_{k}$ the current solid domain is given by $\Phi_{t_{k}}\left(\Omega_{s}\right)=\left(i d+d_{s}^{k-1}\right)\left(\Omega_{s}\right)$ and the current fluid domain $\Phi_{t_{k}}\left(\Omega_{f}\right)$ is obtained by extending the mapping $\Phi_{t_{k}}$ in the whole domain $\Omega$. Moreover, because of the homogeneous Dirichlet boundary conditions on the external frontier $\partial \Omega$, it is clear that for all $k \geq 0, \Phi_{t_{k}}(\Omega)=\Omega$. Furthermore, the matrices $F_{t_{k}}$ and $G_{t_{k}}$ only depend on the displacement of the structure at time $t_{k-1}$.

Hence, for all $k \geq 0$, the triplet $\left(u_{f}^{k}, p_{f}^{k}, d_{s}^{k}\right)$ is solution of the following problem,

$$
\left\{\begin{array}{rlrl}
\text { find } u: \Phi_{t_{k}}\left(\Omega_{f}\right) \rightarrow \mathbb{R}^{n}, p: \Phi_{t_{k}}\left(\Omega_{f}\right) \rightarrow \mathbb{R} \text { and } d: \Omega_{s} \rightarrow \mathbb{R}^{n} \text { such that }  \tag{38a}\\
-\operatorname{div}\left(\sigma_{f}(u, p)\right) & =f_{f}^{k}, & & \text { in } \Phi_{t_{k}}\left(\Omega_{f}\right), \\
\operatorname{div}(u) & =0, & & \text { in } \Phi_{t_{k}}\left(\Omega_{f}\right), \\
u & =0, & & \text { in } \Phi_{t_{k}}\left(\Gamma_{f}\right), \\
-\operatorname{div}\left(\sigma_{s}(d)\right) & =f_{s}^{k}, & & \text { in } \Omega_{s}, \\
d & =0, & & \text { on } \Gamma_{s}, \\
d & =d_{s}^{k-1}+\delta t u \circ \Phi_{t_{k}}, & & \text { on } \Gamma, \\
\sigma_{s}(d) n_{s} & =\Pi_{f}\left(u \circ \Phi_{t_{k}}, p \circ \Phi_{t_{k}}\right) n_{s}, & & \text { on } \Gamma,
\end{array}\right.
$$

where by convention $d_{s}^{-1}=0$ for $k=0$. We can define a weak formulation of problem (38). Let $X$ stand for either $\Omega_{f}, \Omega_{s}, \Phi_{t_{k}}\left(\Omega_{f}\right)$ or $\Phi_{t_{k}}\left(\Omega_{s}\right)$. We introduce the following functional spaces:

$$
\begin{aligned}
V_{f}^{k} & =\left\{v \in\left(H^{1}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)\right)^{n} ; v_{\mid \Phi_{t}}\left(\Gamma_{f}\right)=0\right\}, \\
V_{s}^{k} & =\left\{v \in\left(H^{1}\left(\Phi_{t_{k}}\left(\Omega_{s}\right)\right)\right)^{n} ; v_{\mid \Phi_{t_{k}}\left(\Gamma_{s}\right)}=0\right\}, \\
V_{s} & =\left\{v \in\left(H^{1}\left(\Omega_{s}\right)\right)^{n} ; v_{\mid \Gamma_{s}}=0\right\}, \\
W_{u} & =\left\{\left(v_{f}, v_{s}\right) \in V_{f}^{k} \times V_{s} ;\left(v_{f} \circ \Phi_{t_{k}}\right)_{\mid \Gamma}=v_{s \mid \Gamma}\right\}, \\
W_{d} & =\left\{\left(v_{f}, d_{s}\right) \in V_{f}^{k} \times V_{s} ; \delta_{t}\left(v_{f} \circ \Phi_{t_{k}}\right)_{\mid \Gamma}+d_{s}^{k-1}=d_{s \mid \Gamma}\right\}, \\
\left(H_{\operatorname{div}}(X)\right)^{n} & =\left\{\sigma \in\left(L^{2}(X)\right)^{n \times n} ; \operatorname{div}(\sigma) \in\left(L^{2}(X)\right)^{n}\right\} .
\end{aligned}
$$

The weak formulation of problem (38) writes:

$$
\left\{\begin{array}{rlr}
\text { find }(u, d) \in W_{d} \text { and } p \in L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right) \text { such that }  \tag{39}\\
\int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} \sigma_{f}(u, p): \nabla v_{f}+\int_{\Omega_{s}} \sigma_{s}(d): \nabla v_{s} \\
= & \int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} f_{f}^{k} \cdot v_{f}+\int_{\Omega_{s}} f_{s}^{k} \cdot v_{s}, & \forall\left(v_{f}, v_{s}\right) \in W_{u} \\
\int_{\Phi_{t_{k}\left(\Omega_{f}\right)}} q \operatorname{div}(u)=0, & \forall q \in L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right) .
\end{array}\right.
$$

Problem (39) admits a unique solution that we still denote by $\left(u_{f}^{k}, p_{f}^{k}, d_{s}^{k}\right)$ (see Appendix B). Moreover, $\sigma_{f}\left(u_{f}^{k}, p_{f}^{k}\right)$ belongs to $\left(H_{\mathrm{div}}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)\right)^{n}, \sigma_{s}\left(d_{s}^{k}\right)$ belongs to $\left(H_{\mathrm{div}}\left(\Omega_{s}\right)\right)^{n}$ and we can give a weak sense to the second transmission condition in (38c). Let $X$ stands for either $\Omega_{f}$ or $\Omega_{s}$ and let $\eta$ be the exterior normal vector to $X$. Let $\Upsilon=\left(H_{00}^{1 / 2}(\Gamma)\right)^{n}$ be the
image of $\left(H_{\partial X \backslash \Gamma}^{1}(X)\right)^{n}$ by the trace operator on the interface $\Gamma$, i.e. the space of functions in $\left(H^{1 / 2}(\Gamma)\right)^{n}$ whose extension by zero on $\partial X \backslash \Gamma$ belongs to $\left(H^{1 / 2}(\partial X)\right)^{n}$.

Then, for all $\sigma$ in $\left(H_{\mathrm{div}}(X)\right)^{n}$, we have the following Stokes formula:

$$
\begin{equation*}
\int_{X} \sigma: \nabla v+\int_{X} \operatorname{div}(\sigma) \cdot v=\left\langle\gamma_{\eta}(\sigma), v\right\rangle_{\Upsilon^{\prime}, \Upsilon}, \quad \forall v \in\left(H_{\partial X \backslash \Gamma}^{1}(X)\right)^{n} \tag{40}
\end{equation*}
$$

where $\Upsilon^{\prime}$ is the dual space of $\Upsilon$ and $\gamma_{\eta}$ is the trace normal operator on $\Gamma$. Then, the second transmission condition in (38c) is satisfied in the following sense:

$$
\begin{equation*}
\left\langle\gamma_{n_{s}}\left(\sigma_{s}\left(d_{s}^{k}\right)\right), v\right\rangle_{\Upsilon^{\prime}, \Upsilon}=-\left\langle\gamma_{n_{f}}\left(\Pi_{f}\left(u_{f}^{k} \circ \Phi_{t_{k}}, p_{f}^{k} \circ \Phi_{t_{k}}\right)\right), v\right\rangle_{\Upsilon^{\prime}, \Upsilon}, \quad \forall v \in \Upsilon . \tag{41}
\end{equation*}
$$

Similarly, we define $\Upsilon^{k}=\left(H_{00}^{1 / 2}\left(\Phi_{t_{k}}(\Gamma)\right)\right)^{n}$ which enables to also write Stokes formulas for a tensor $\sigma$ in $\Phi_{t_{k}}\left(\Omega_{f}\right)$ and in $\Phi_{t_{k}}\left(\Omega_{s}\right)$.

Now, we present the smooth extension method applied to problem (38). Formally, it writes: find $g$ in $\left(V_{s}^{k}\right)^{\prime}$ such that the solution of the following problem,

$$
\left\{\begin{array}{rlrlrl}
\text { find } u: \Omega \rightarrow \mathbb{R}^{n}, p: \Omega \rightarrow \mathbb{R} & \text { and } d: \Omega_{s} \rightarrow \mathbb{R}^{n} & \text { such that, }  \tag{42a}\\
-\operatorname{div}\left(\sigma_{f}(u, p)\right) & ={\overline{f_{f}^{k}}}^{\Omega}+\bar{g}^{\Omega}, & & \text { in } \Omega, \\
\operatorname{div}(u) & =0, & & \text { in } \Omega \\
u & =0, & & \text { in } \partial \Omega, \\
-\operatorname{div}\left(\sigma_{s}(d)\right) & =f_{s}^{k}, & & \text { in } \Omega_{s} \\
d & =0, & & \text { on } \Gamma_{s} \\
\sigma_{s}(d) n_{s} & =\Pi_{f}\left(u \circ \Phi_{t_{k}}, p \circ \Phi_{t_{k}}\right) n_{s}, & & \text { on } \Gamma
\end{array}\right.
$$

satisfies the equality,

$$
u \circ \Phi_{t_{k}}=\frac{1}{\delta t}\left(d-d_{s}^{k-1}\right), \text { on } \Gamma .
$$

Let $v$ be in $(\mathcal{D}(\Omega))^{n}$ and suppose that $u, p$ and $d$ are sufficiently regular. Moreover, we assume for the moment that $g$ belongs to $\left(L^{2}\left(\Phi_{t_{k}}\left(\Omega_{s}\right)\right)\right)^{n}$ in order to do formal computations. Formally, we multiply the first equation in (42a) by $v$ and the first equation in (42b) by $v_{\mid \Omega_{s}}$, and integrate respectively over $\Omega$ and $\Omega_{s}$. After an integration by part and using the Neumann condition on $\Gamma$, we find

$$
\begin{aligned}
\int_{\Omega} \sigma_{f}(u, p): \nabla v & =\int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} f_{f}^{k} \cdot v+\int_{\Phi_{t_{k}}\left(\Omega_{s}\right)} g \cdot v \\
\int_{\Omega_{s}} \sigma_{s}(d): \nabla v_{\mid \Omega_{s}} & =\int_{\Omega_{s}} f_{s}^{k} \cdot v_{\mid \Omega_{s}}+\int_{\Gamma}\left(\Pi_{f}\left(u \circ \Phi_{t_{k}}, p \circ \Phi_{t_{k}}\right) n_{s}\right) \cdot v
\end{aligned}
$$

Moreover, we remark that

$$
\begin{aligned}
\int_{\Gamma}\left(\Pi_{f}\right. & \left.\left(u \circ \Phi_{t_{k}}, p \circ \Phi_{t_{k}}\right) n_{s}\right) \cdot v \\
& =\int_{\Omega_{s}} \operatorname{div}\left(\Pi_{f}\left(u \circ \Phi_{t_{k}}, p \circ \Phi_{t_{k}}\right)\right) \cdot v_{\mid \Omega_{s}}+\int_{\Omega_{s}} \Pi_{f}\left(u \circ \Phi_{t_{k}}, p \circ \Phi_{t_{k}}\right): \nabla v_{\mid \Omega_{s}}, \\
& =-\int_{\Omega_{s}} \operatorname{det}\left(\nabla \Phi_{t_{k}}\right)\left(g \circ \Phi_{t_{k}}\right) \cdot v_{\mid \Omega_{s}}+\int_{\Omega_{s}} \Pi_{f}\left(u \circ \Phi_{t_{k}}, p \circ \Phi_{t_{k}}\right): \nabla v_{\mid \Omega_{s}}, \\
& =-\int_{\Phi_{t_{k}}\left(\Omega_{s}\right)} g \cdot\left(v_{\mid \Omega_{s}} \circ \Phi_{t_{k}}^{-1}\right)+\int_{\Omega_{s}} \Pi_{f}\left(u \circ \Phi_{t_{k}}, p \circ \Phi_{t_{k}}\right): \nabla v_{\mid \Omega_{s}}
\end{aligned}
$$

Thus, we define the weak formulation of the smooth extension problem (42), which makes sense for $g$ in $\left(V_{s}^{k}\right)^{\prime}$, as the problem of finding a suitable control $g$ in $\left(V_{s}^{k}\right)^{\prime}$, such that the solution of the following problem,

$$
\left\{\begin{array}{rlrl}
\text { find } u \in H_{0}^{1}(\Omega), p \in L_{0}^{2}(\Omega) \text { and } d \in V_{s} \text { such that } & &  \tag{43}\\
\int_{\Omega} \sigma_{f}(u, p): \nabla v_{f}= & \int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} f_{f}^{k} \cdot v_{f \mid \Phi_{t_{k}}\left(\Omega_{f}\right)} & & \\
\int_{\Omega} q \operatorname{div}(u)= & 0, & & \forall v_{f} \in H_{0}^{1}(\Omega), \\
\int_{\Omega_{s}} \sigma_{s}(d): \nabla v_{s}= & \left.\int_{\Omega_{k}} f_{s}\right\rangle_{s}^{k} \cdot v_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}},\left\langle g, v_{s} \circ \Phi_{t_{k}}^{-1}\right\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}} & & \\
& +\int_{\Omega_{s}} \Pi_{f}\left(u \circ \Phi_{t_{k}}, p \circ \Phi_{t_{k}}\right): \nabla v_{s}, & & \forall v_{s} \in V_{s},
\end{array}\right.
$$

verifies the equality

$$
\begin{equation*}
u \circ \Phi_{t_{k}}=\frac{1}{\delta t}\left(d-d_{s}^{k-1}\right), \text { on } \Gamma \tag{44}
\end{equation*}
$$

For every $f_{f}^{k} \in\left(L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)\right)^{n}$ and every $g \in\left(V_{s}^{k}\right)^{\prime}$, there exists a unique solution to the Stokes problem appearing in problem (43), denoted by ( $u^{g}, p^{g}$ ) (see [10]). On the other hand, for all $f_{s}^{k} \in\left(L^{2}\left(\Omega_{s}\right)\right)^{n}$, the weak problem of linear elasticity that appears in (43) also admits a unique solution in $V_{s}$, denoted by $d^{g}$ (see [13]). In the following theorem, we state the existence of a control $g$ such that the solution of (43) satisfies the equality (44) and explain how to recover the solution of the initial fluid-structure problem (39).

Theorem 4.4. Let $\Omega$ be a domain that satisfies Assumption $\left(H_{4}\right)$. Consider $f_{f}^{k}$ in $\left(L_{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)\right)^{n}$ and $f_{s}^{k}$ in $\left(L_{2}\left(\Omega_{s}\right)\right)^{n}$.Then, there exists a function $g$ in $\left(V_{s}^{k}\right)^{\prime}$ such that the solution $\left(u^{g}, p^{g}, d^{g}\right)$ of problem (43) satisfies (44). Moreover, we can recover the solution of the fluid-structure problem (39):

$$
\begin{aligned}
\left(u_{f}^{k}, p_{f}^{k}\right) & =\left(u_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, p_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}\right) \\
d_{s}^{k} & =d^{g}
\end{aligned}
$$

The proof of Theorem 4.4 relies on the construction of extensions for the functions $u_{f}^{k}$ and $p_{f}^{k}$ to the whole domain $\Omega$, which is what we have done for the Stokes transmission problem. Hence, the proof of Theorem 4.4 directly follows from the one of Theorem 4.1. This is detailed in Appendix B.

As before, the problem of finding a control such that the solution of (42) verifies (44) can be formulated as an optimization problem on the following cost function, defined for any $k \geq 0$,

$$
\begin{align*}
J_{k}: V_{s, k}^{\prime} & \rightarrow \mathbb{R}^{+} \\
g & \mapsto \frac{1}{2} \int_{\Gamma}\left|u^{g} \circ \Phi_{t_{k}}-\frac{1}{\delta t}\left(d^{g}-d_{s}^{k-1}\right)\right|^{2}, \tag{45}
\end{align*}
$$

where $u^{g}$ and $d^{g}$ are the velocity of the fluid and the displacement of the structure that solve problem (43) at time $t_{k}$. The function $d_{s}^{k-1}$ is the displacement of the structure that solves problem (38) at time $t_{k-1}$. Yet, the minimization of this cost function is equivalent to the minimization of the real-valued function,

$$
\begin{aligned}
\left(H_{0}^{1}(\Omega)\right)^{n} \times V_{s} & \rightarrow \mathbb{R}^{+} \\
(u, d) & \mapsto \frac{1}{2} \int_{\Gamma}\left|u \circ \Phi_{t_{k}}-\frac{1}{\delta t}\left(d-d_{s}^{k-1}\right)\right|^{2}
\end{aligned}
$$

under the constraint that $u$ and $d$ are the velocity of the fluid and the displacement of the structure that solve problem (43). Then, to this constrained optimization problem, we can associate the following Lagrangian function, defined from

$$
\left(V_{s}^{k}\right)^{\prime} \times\left(\left(H_{0}^{1}(\Omega)\right)^{n} \times L_{0}^{2}(\Omega) \times V_{s}\right) \times\left(\left(H_{0}^{1}(\Omega)\right)^{n} \times L_{0}^{2}(\Omega) \times V_{s}\right),
$$

to $\mathbb{R}$ by,

$$
\begin{align*}
& \mathcal{L}_{k}\left(g,(u, p, d),\left(\lambda_{f}, \pi, \nu_{s}\right)\right)= \\
& \quad \frac{1}{2} \int_{\Gamma}\left|u \circ \Phi_{t_{k}}-\frac{1}{\delta t}\left(d-d_{s}^{k-1}\right)\right|^{2}+\int_{\Omega} \sigma_{f}(u, p): \nabla \lambda_{f}-\int_{\Omega} \pi \operatorname{div}(u) \\
& \quad+\int_{\Omega_{s}} \sigma_{s}(d): \nabla \nu_{s}-\int_{\Omega_{s}} \Pi_{f}\left(u \circ \Phi_{t_{k}}, p \circ \Phi_{t_{k}}\right): \nabla \nu_{s}  \tag{46}\\
& \quad-\int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} f_{f}^{k} \cdot \lambda_{f}-\left\langle g, \lambda_{f}\right\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}}-\int_{\Omega_{s}} f_{s}^{k} \cdot \nu_{s}+\left\langle g, \nu_{s} \circ \Phi_{t_{k}}^{-1}\right\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}} .
\end{align*}
$$

Again, this Lagrangian function can be used to compute the gradient of $J_{k}$ and show that the minimization of $J_{k}$ is equivalent to the problem of finding a suitable control $g$ such that the solution of (43) satisfies (44). For that purpose, we introduce the adjoint problem of (43),

$$
\left\{\begin{array}{rlrlrl}
\text { find } \nu_{s} \in V_{s}, \lambda_{f} \in\left(H_{0}^{1}(\Omega)\right)^{n} \text { and } \pi \in L_{0}^{2}(\Omega) \text { such that, } & \\
\int_{\Omega_{s}} \sigma_{s}\left(\nu_{s}\right): \nabla v_{s} & =\frac{1}{\delta t} \int_{\Gamma}\left(u^{g} \circ \Phi_{t_{k}}-\frac{1}{\delta t}\left(d^{g}-d_{s}^{k-1}\right) \cdot v_{s},\right. & & \forall v_{s} \in V_{s}, \\
\int_{\Omega} \sigma_{f}\left(\lambda_{f}, \pi\right): \nabla v_{f} & =2 \mu_{f} \int_{\Phi_{t_{k}}\left(\Omega_{s}\right)} D\left(\nu_{s} \circ \Phi_{t_{k}}^{-1}\right): D\left(v_{f}\right) & & \\
& -\int_{\Gamma}\left(u^{g} \circ \Phi_{t_{k}}-\frac{1}{\delta t}\left(d^{g}-d_{s}^{k-1}\right)\right) \cdot v_{f} \circ \Phi_{t_{k}}, & & \forall v_{f} \in\left(H_{0}^{1}(\Omega)\right)^{n}, \\
\int_{\Omega} q \operatorname{div}\left(\lambda_{f}\right) & =\int_{\Phi_{t_{k}\left(\Omega_{s}\right)}} q \operatorname{div}\left(\nu_{s} \circ \Phi_{t_{k}}^{-1}\right), & & \forall q \in L_{0}^{2}(\Omega) . \tag{47}
\end{array}\right.
$$

Problem (47) admits a unique solution that we denote $\left(\nu^{g}, \lambda^{g}, \pi^{g}\right)$. Then, the existence and the characterization of the gradient of $J_{k}$ is given in the following theorem,

Theorem 4.5. The mapping $g \in\left(V_{s}^{k}\right)^{\prime} \mapsto J_{k}(g) \in \mathbb{R}^{+}$is differentiable and its gradient $\nabla J_{k}(g)$ in $\left(V_{s}^{k}\right)^{\prime \prime}$ is given by,

$$
\begin{equation*}
\left\langle\nabla J_{k}(g), \delta g\right\rangle_{\left(V_{s}^{k}\right)^{\prime \prime},\left(V_{s}^{k}\right)^{\prime}}=\left\langle\delta g, \nu^{g} \circ \Phi_{t_{k}}^{-1}-\lambda_{\mid \Phi_{t_{k}}\left(\Omega_{s}\right)}^{g}\right\rangle_{\left(V_{s}^{k}\right)^{\prime},\left(V_{s}^{k}\right)}, \quad \forall \delta g \in V_{s}^{\prime}, \tag{48}
\end{equation*}
$$

where $\lambda^{g}$ and $\nu^{g}$ satisfy the adjoint problem (47).
Moreover, the equivalence between the smooth extension problem and its formulation as a minimization problem can also be stated in the case of a fluid-structure interaction problem.

Theorem 4.6. A control $g$ in $\left(V_{s}^{k}\right)^{\prime}$ is a minimizer of $J_{k}$ if and only if the solution of (43) satisfies (44).

Proofs of Theorem 4.5 and Theorem 4.6 can be easily adapted from the ones already done for the Stokes transmission problem. However, the change of domains between the structure in reference configuration and the fluid in current configuration can be confusing. For that matter, all proofs are detailed in Appendix B.


Figure 13: Two dimensional representation of the initial geometry for the fluid-structure problem.

The formulation of the fluid-structure interaction problem as a control problem and as a minimization problem enables us to numerically solve these equations with a fictitious domain approach. Here, we apply this method to the numerical simulation of the bending of an elastic beam in a viscous fluid subjected to shear boundary condition. The initial geometry of the problem is represented in Figure 13: $\Omega_{f}$ is the rectangle $[0,2] \times[0,1]$ in $\mathbb{R}^{2}$ and $\Omega_{s}$ is a rectangular beam of length $L_{c}$ and of radius $r_{c}$. This beam is anchored at the bottom of $\Omega_{f}$, at positions $\left(x_{c}-r_{c}, 0\right)$ and $\left(x_{c}+r_{c}, 0\right)$. We consider periodic boundary conditions on the left and right boundaries of the fluid domain and imposed a shear condition on the top, given by

$$
u_{b c}(x)=(3,0), \quad \forall x \in[0,1] .
$$

On the bottom boundaries of the fluid and solid domains, homogeneous Dirichlet boundary conditions are considered. No external force is considered and constant values are chosen for the fluid viscosity, $\mu_{f}=1$, the Young's modulus of the solid material, $E_{s}=10^{5}$, and its Poisson's ratio, $\nu_{s}=0.49$. The Lamé coefficients $\mu_{s}$ and $\lambda_{s}$ are then given by

$$
\begin{aligned}
\mu_{s} & =\frac{E_{s}}{2\left(1+\nu_{s}\right)} \\
\lambda_{s} & =\frac{\nu_{s} E_{s}}{\left(\left(1+\nu_{s}\right)\left(1-2 \nu_{s}\right)\right)}
\end{aligned}
$$

For the parameters of the beam, we choose $x_{c}=1, r_{c}=0.05$ and $L_{c}=0.65$. The final time of the simulation is set to $T=0.5$.

The solution given by the Smooth Extension method is obtained following Algorithm 1 and using the L-BFGS algorithm (see [28]). The resolution of the two different Stokes problems appearing in the minimization process is done using Mini elements, while the resolution of the two elasticity problems is done using $P 1$ elements. Fluid problems are solved on a fixed mesh representing the whole domain $\Omega$ and elasticity problems are solve on a fixed mesh representing $\Omega_{s}$, i.e. in the solid reference configuration. Moreover, the fluid mesh does not conform with the solid boundary. We compare this solution to the one obtained with conformal meshes and using a Lagrangian multiplier to ensure the continuity of the fluid and solid velocities through the interface $\Gamma$. With this method, an Arbitrary


Figure 14: Comparison of the stationary states obtained through the Smooth Extension method and the ALE method (a). The darkest grey mesh represents the equilibrium position of the beam obtained with the ALE method. The lightest grey mesh represents the initial position of the beam. The other three meshes represent the equilibrium positions of the beam obtained with the Smooth Extension method for different coarsening ratios ( $0.99,0.71,0.46$ ): the lighter is the colour the coarser is the fluid mesh. Zoom on the tips of the beams (b).

Lagrangian-Eulerian (ALE) methodology (see [9]) is used to ensure the mesh conformity when the beam bends and prevent cells from overlapping within the fluid mesh. However, if the quality of the fluid mesh is too poor, one needs to remesh it (this is done using the Mmg remeshing software [15]). In the following, this method will be referred as the ALE method.

In this test case, the fluid-structure system attains a stationary state, where the beam is at equilibrium in a deformed configuration, which is well catched by both the ALE and the Smooth Extension methods (see Figure 14). To study the robustness of the Smooth Extension method, we coarsen the fluid mesh and observe the consequences on the dynamic of the system. To do so we define the coarsening ratio of the fluid mesh as the ratio of the number of nodes in the reference fluid mesh used in the ALE method divided by the number of nodes of the fluid mesh used in the Smooth Extension method. For example, a coarsening ratio of 1 means that the two fluid meshes have the same number of nodes, whereas a ratio of 0.5 means that the fluid mesh used in the Smooth Extension method has half as many nodes than the reference fluid mesh used for the ALE method. Then, in Figure 14, we observe that the coarser is the fluid mesh, the farther is the stationary state from the equilibrium state obtained with the ALE method. In order to quantify this error, we consider the $L^{1}(\Omega)$ distance between the position of the beam obtained with the ALE method and the one obtained with the Smooth Extension method at time $t$, defined by

$$
d_{S E M}^{A L E}(t)=\frac{\int_{\Omega}\left|\chi_{\Omega_{s}^{S E M}(t)}-\chi_{\Omega_{s}^{A L E}(t)}\right|}{\int_{\Omega}\left|\chi_{\Omega_{s}^{S E M}(t)}\right|+\int_{\Omega}\left|\chi_{\Omega_{s}^{A L E}(t)}\right|},
$$

where $\chi_{\Omega_{s}^{S E M}(t)}$ is a $L^{2}(\Omega)$ function which has value 1 inside the solid current domain $\Omega_{s}^{S E M}(t)$ obtained with the Smooth Extension method, and 0 otherwise. Similarly, $\chi_{\Omega_{s}^{A L E}(t)}$ is a $L^{2}(\Omega)$ function which has value 1 inside the solid current domain $\Omega_{s}^{A L E}(t)$ obtained


Figure 15: Distance $d_{S E M}^{A L E}$ between the two solid configuration $\Omega_{s}^{S E M}$ and $\Omega_{s}^{A L E}$ in function of the time (a). Global norm in time $\left\|d_{S E M}^{A L E}\right\|$ in function of the coarsening ratio of the fluid mesh in the Smooth Extension method (b).
with the ALE method, and 0 otherwise. Consequently, this distance is 0 if the two domains are the same and 1 if they do not overlap. Moreover, for a final time $T>0$, we also consider the global norm in time of $d_{S E M}^{A L E}$ defined by

$$
\left\|d_{S E M}^{A L E}\right\|=\frac{1}{T} \int_{0}^{T}\left|d_{S E M}^{A L E}(t)\right| d t
$$

which is 0 if the two domains $\Omega_{s}^{A L E}(t)$ and $\Omega_{s}^{S E M}(t)$ are identical for all $t$ in $[0, T]$ and 1 if they never overlap. Then, we compute the distance $d_{S E M}^{A L E}$ in function of the time for different coarsening ratios to study their influence on the dynamic of the system when using the Smooth Extension method. This is represented in Figure 15a, where we observe that, for all coarsening ratios, the distance $d_{S E M}^{A L E}$ increases in time to reach a constant value when the stationary state is attained. This result corroborates and quantifies what we observed on Figure 14, i.e. that the error on the equilibrium position of the beam seems to increase when the coarsening ratio decreases, but stays relatively low considering the coarsening of the fluid mesh. To go further, we plot in Figure 15b the global norm in time of $d_{S E M}^{A L E}$ in function of the coarsening ratio. In addition to the already mentioned fact that the error tends to increase when the coarsening ratio decreases, we remark that the error in time is just above 0.1 for a coarsening ratio of 0.12 , which implies that the coarsening of the fluid mesh in the Smooth Extension method does not drastically change the time dynamic of the bending of the beam in this test case.

All these results suggest that the Smooth Extension method is well suited for time dependent problems involving a moving structure in a viscous fluid, where the fluid mesh is fixed, possibly Cartesian and coarser that the structure mesh.

## 5 Conclusion

In this chapter we have presented a numerical strategy for the resolution of transmission problems with non-conformal meshes and which preserves optimal rates of convergence in
space. It is based on a control formulation of the transmission problem, namely the Smooth Extension formulation, whose numerical resolution can be done by minimizing a particular objective function. This method allows the use of standard finite element functional spaces along with fixed structured or unstructured meshes and pre-existing finite element solvers and optimization algorithms.

This smooth extension method has been derived in the particular case of the transmission Laplace problem with only two subdomains and considering Dirichlet boundary conditions. Other boundary conditions could also be considered with no additional difficulty, provided that the initial transmission problem is well-posed. The same methodology should also work for transmission problems with more than two subdomains. In addition, we have shown that the Smooth Extension method can be applied to a wide variety of transmission problems, even the ones with totally different operators, such as the fluidstructure interaction problem studied in subsection 4.2.

The Smooth Extension method has been compared to a standard numerical method for different test cases and has shown to give good approximations of the solutions for each one of them. However, some numerical aspects, such as the use of fast solvers on Cartesian meshes or the choice of the optimization algorithm appearing in the minimization process, should be explored further.

## A Proofs of theorems related to the Stokes transmission problem

This appendix is dedicated to the proofs of all results stated in subsection 4.1. In particular, we are interested in showing the existence of the control $g$, in giving an explicit formula for the gradient of $\tilde{J}$ and, finally, in proving the equivalence between the minimisation of $\tilde{J}$ and the resolution of the Stokes transmission problem.

We follow the order of the previous enumeration and start with the proof of Theorem 4.1.

Proof of Theorem 4.1. We can construct two extension operators $E_{u}$ and $E_{p}$ that extend $\bar{w}_{1}$ into the whole space $\left(H_{0}^{1}(\Omega)\right)^{n}$ (and such that $E_{u} \bar{w}_{1}$ is divergence-free) and $\bar{p}_{1}$ into the whole space $L_{0}^{2}(\Omega)$. Indeed, consider the operators defined by

$$
\begin{aligned}
E_{u} \bar{w}_{1} & =\left\{\begin{array}{lll}
\bar{w}_{1} & \text { in } & \Omega_{1} \\
\bar{w}_{2} & \text { in } & \Omega_{2} \\
E_{p} \bar{p}_{1} & =\left\{\begin{array}{lll}
\bar{p}_{1} & \text { in } & \Omega_{1} \\
\bar{p}_{2} & \text { in } & \Omega_{2}
\end{array} .\right.
\end{array} . \begin{array}{l}
\text { an }
\end{array}\right)
\end{aligned}
$$

Since $\left(\bar{w}_{1}, \bar{w}_{2}\right)$ belongs to $\mathcal{W}$, the function $E_{u} \bar{w}_{1}$ is an extension of $\bar{w}_{1}$ which belongs to $\left(H_{0}^{1}(\Omega)\right)^{n}$. Moreover, it satisfies $\operatorname{div}\left(E_{u} \bar{w}_{1}\right)=0$ because both $\bar{w}_{1}$ and $\bar{w}_{2}$ are divergencefree. Similarly, since ( $\bar{p}_{1}, \bar{p}_{2}$ ) belongs to $Q$, the function $E_{p} \bar{p}_{1}$ is an extension of $\bar{p}_{1}$ which belongs to $L_{0}^{2}(\Omega)$. Furthermore, $\mu_{1} \nabla\left(E_{u} \bar{w}_{1}\right)-\left(E_{p} \bar{p}_{1}\right) I$ belongs to $\left(H_{\text {div }}(\Omega)\right)^{n}$.

Then, we construct a suitable control $g$ in $W_{2}^{\prime}$ such that

$$
\begin{align*}
\langle g, v\rangle_{W_{2}^{\prime}, W_{2}}= & \mu_{1} \int_{\Omega_{2}} \nabla\left(E_{u} \bar{w}_{1}\right)_{\Omega_{2}}: \nabla v-\int_{\Omega_{2}} \bar{p}_{1 \mid \Omega_{2}} \operatorname{div}(v)  \tag{49}\\
& +\left\langle\gamma_{n_{1}}\left(\sigma_{E}\right), v\right\rangle_{\Upsilon^{\prime}, \Upsilon},
\end{align*} \forall v \in W_{2},
$$

where we define $\sigma_{E}=\mu_{1} \nabla\left(E_{u} \bar{w}_{1}\right)-\left(E_{p} \bar{p}_{1}\right) I$. Using the Stokes formula (27) and the
definition (49), it follows that the extensions $E_{u} \bar{w}_{1}$ and $E_{p} \bar{p}_{1}$ satisfy

$$
\begin{array}{rlrl}
\mu_{1} \int_{\Omega} \nabla\left(E_{u} \bar{w}_{1}\right): \nabla v_{1}-\int_{\Omega}\left(E_{p} \bar{p}_{1}\right) \operatorname{div}\left(v_{1}\right) & & \\
& =\int_{\Omega_{1}} f_{1} \cdot v_{1 \mid \Omega_{1}}+\left\langle g, v_{1 \mid \Omega_{2}}\right\rangle_{W_{2}^{\prime}, W_{2}}, & & \forall v_{1} \in H_{0}^{1}(\Omega) \\
\int_{\Omega} q_{1} \operatorname{div}\left(E_{u} \bar{w}_{1}\right) & =0, & \forall q_{1} \in L_{0}^{2}(\Omega) .
\end{array}
$$

Similarly, using (27), (49) and the weak transmission condition (28), ( $\bar{w}_{2}, \bar{p}_{2}$ ) satisfies

$$
\begin{gathered}
\mu_{2} \int_{\Omega_{2}} \nabla \bar{w}_{2}: \nabla v_{2}-\int_{\Omega_{2}} \bar{p}_{2} \operatorname{div}\left(v_{2}\right) \\
=\int_{\Omega_{2}} f_{2} \cdot v_{2}+\mu_{1} \int_{\Omega_{2}} \nabla\left(E_{u} u_{1}\right)_{\Omega_{2}}: \nabla v_{2}-\left\langle g, v_{2}\right\rangle_{W_{2}^{\prime}, W_{2}}-\int_{\Omega_{2}}\left(E_{p} \bar{p}_{1 \mid \Omega_{2}}\right) \operatorname{div}\left(v_{2}\right), \forall v_{2} \in W_{2},
\end{gathered}
$$

Finally, we conclude that $\left(\left(E_{u} \bar{w}_{1}, E_{p} \bar{p}_{1}\right),\left(\bar{w}_{2}, \bar{p}_{2}\right)\right)$ is the solution of problem (30). Thus,

$$
\left(\left(w_{1}^{g}, p_{1}^{g}\right),\left(w_{2}^{g}, p_{2}^{g}\right)\right)=\left(\left(E_{u} \bar{w}_{1}, E_{p} \bar{p}_{1}\right),\left(\bar{w}_{2}, \bar{p}_{2}\right)\right)
$$

and, by construction, the equality (31) is satisfied. This proves the first part of the theorem.
Now, suppose that $g$ is a control such that the equality (31) is satisfied. In particular, the equality $w_{1 \mid \Gamma}^{g}=w_{2 \mid \Gamma}^{g}$ implies that $\left(w_{1 \mid \Omega_{1}}^{g}, w_{2}^{g}\right)$ belongs to the space $\mathcal{W}$. Let us define the following Hilbert space:

$$
\tilde{\mathcal{W}}=\left\{(v 1, v 2) \in\left(H_{0}^{1}(\Omega)\right)^{n} \times W_{2} ; v_{1 \mid \Omega_{1}}=v_{2}\right\} .
$$

As the unique solution of problem (30), the couples $\left(w_{1}^{g}, p_{1}^{g}\right)$ and $\left(w_{2}^{g}, p_{2}^{g}\right)$ satisfy, in particular, for all $\left(v_{1}, v_{2}\right)$ in $\tilde{\mathcal{W}}$ and for all $\left(q_{1}, q_{2}\right)$ in $L_{0}^{2}(\Omega) \times L^{2}\left(\Omega_{2}\right)$, the equations

$$
\begin{align*}
\mu_{1} \int_{\Omega} \nabla w_{1}^{g}: \nabla v_{1}-\int_{\Omega_{1}} p_{1}^{g} \operatorname{div}\left(v_{1}\right) & =\int f_{1} \cdot v_{1 \mid \Omega_{1}}+\left\langle g, v_{1 \mid \Omega_{2}}\right\rangle_{W_{2}^{\prime}, W_{2}}, \\
\int_{\Omega} q \operatorname{div}\left(w_{1}^{g}\right)= & 0, \\
\mu_{2} \int_{\Omega_{2}} \nabla w_{2}^{g}: \nabla v_{2}-\int_{\Omega_{2}} p_{2}^{g} \operatorname{div}\left(v_{2}\right)= & \int f_{2} \cdot v_{2}-\left\langle g, v_{2}\right\rangle_{W_{2}^{\prime}, W_{2}}  \tag{50}\\
& +\mu_{1} \int_{\Omega_{2}} \nabla w_{1}^{g}: \nabla v_{2}-\int_{\Omega_{2}} p_{1}^{g} \operatorname{div}\left(v_{2}\right), \\
\int_{\Omega_{2}} q_{\mid \Omega_{2}} \operatorname{div}\left(w_{2}^{g}\right)= & 0 .
\end{align*}
$$

Then, summing equations in (50), it follows that $\left(\left(w_{1 \mid \Omega_{1}}^{g}, w_{2}^{g}\right),\left(p_{1 \mid \Omega_{1}}^{g}, p_{2}^{g}\right)\right)$ satisfies

$$
\begin{aligned}
& \mu_{1} \int_{\Omega_{1}} \nabla w_{1 \mid \Omega_{1}}^{g}: \nabla v_{1 \mid \Omega_{1}}-\int_{\Omega_{1}} p_{1 \mid \Omega_{1}}^{g} \operatorname{div}\left(v_{1 \mid \Omega_{1}}\right)+\mu_{2} \int_{\Omega_{2}} \nabla w_{2}^{g}: \nabla v_{2}-\int_{\Omega_{2}} p_{2}^{g} \operatorname{div}\left(v_{2}\right) \\
& =\int_{\Omega_{1}} f_{1} \cdot v_{1}+\int_{\Omega_{2}} f_{2} \cdot v_{2}, \quad \forall\left(v_{1}, v_{2}\right) \in \tilde{\mathcal{W}}, \\
& \int_{\Omega_{1}} q_{1 \mid \Omega_{1}} \operatorname{div}\left(w_{1 \mid \Omega_{1}}^{g}\right)+\int_{\Omega_{2}} q_{2} \operatorname{div}\left(w_{2}^{g}\right)=0, \quad \forall\left(q_{1}, q_{2}\right) \in L_{0}^{2}(\Omega) \times L^{2}\left(\Omega_{2}\right) .
\end{aligned}
$$

Because the test function $v_{1}$ belongs to $H_{0}^{1}(\Omega)$ and because $v_{2}=v_{1 \mid \Omega_{1}}$, we can redefine the pressures $p_{1}^{g}$ and $p_{2}^{g}$ up to a constant, by

$$
\begin{aligned}
& \tilde{p}_{1}^{g}=p_{1}^{g}-\frac{1}{|\Omega|}\left(\int_{\Omega_{2}} p_{2}^{g}+\int_{\Omega_{1}} p_{1 \mid \Omega_{1}}^{g}\right), \\
& \tilde{p}_{2}^{g}=p_{2}^{g}-\frac{1}{|\Omega|}\left(\int_{\Omega_{2}} p_{2}^{g}+\int_{\Omega_{1}} p_{1 \mid \Omega_{1}}^{g}\right),
\end{aligned}
$$

such that $\left(\tilde{p}_{1 \mid \Omega_{1}}^{g}, \tilde{p}_{2}^{g}\right)$ belongs to the space $Q$. Moreover, for all $\left(v_{1}, v_{2}\right)$ in $\mathcal{W}$ we can extend $v_{1}$ in the whole space $H_{0}^{1}(\Omega)$ using $v_{2}$, as we did it for $\bar{w}_{1}$ and $\bar{w}_{2}$. We still denote $E_{u} v_{1}$ this extension. Similarly, for all $\left(q_{1}, q_{2}\right)$ in $Q$, we can extend $q_{1}$ in the whole space $L_{0}^{2}(\Omega)$ and we still denote $E_{p} q_{1}$ this extension. Then, the couple ( $E_{u} v_{1}, v_{2}$ ) belongs to $\tilde{\mathcal{W}}$, the couple $\left(E_{p} q_{1}, q_{2}\right)$ belongs to $L_{0}^{2}(\Omega) \times L^{2}\left(\Omega_{2}\right)$ and, finally, $\left(\left(w_{1 \mid \Omega_{1}}^{g}, w_{2}^{g}\right),\left(\tilde{p}_{1 \mid \Omega_{1}}^{g}, \tilde{p}_{2}^{g}\right)\right)$, satisfies the equations

$$
\begin{array}{rlrl}
\mu_{1} \int_{\Omega_{1}} \nabla w_{1 \mid \Omega_{1}}^{g}: \nabla v_{1}-\int_{\Omega_{1}} \tilde{p}_{1 \mid \Omega_{1}}^{g} \operatorname{div}\left(v_{1}\right) & +\mu_{2} \int_{\Omega_{2}} \nabla w_{2}^{g}: \nabla v_{2}-\int_{\Omega_{2}} \tilde{p}_{2}^{g} \operatorname{div}\left(v_{2}\right) \\
& =\int_{\Omega_{1}} f_{1} \cdot v_{1}+\int_{\Omega_{2}} f_{2} \cdot v_{2}, & \forall\left(v_{1}, v_{2}\right) \in \mathcal{W}, \\
\int_{\Omega_{1}} q_{1 \mid \Omega_{1}} \operatorname{div}\left(w_{1 \mid \Omega_{1}}^{g}\right)+\int_{\Omega_{2}} q_{2} \operatorname{div}\left(w_{2}^{g}\right) & =0, & \forall\left(q_{1}, q_{2}\right) \in Q .
\end{array}
$$

Thus, $\left(\left(w_{1 \mid \Omega_{1}}^{g}, w_{2}^{g}\right),\left(\tilde{p}_{1 \mid \Omega_{1}}^{g}, \tilde{p}_{2}^{g}\right)\right)$ is the unique solution of the initial Stokes transmission problem (26), which proves that we can recover the solution of the Stokes transmission problem from the solution of the smooth extension problem.

We continue with the proof of Theorem 4.2, i.e. we show the existence of the gradient of $\tilde{J}$ and give an explicit formula to compute it.

Proof of Theorem 4.2. We start by computing the derivatives of the Lagrangian function $\tilde{\mathcal{L}}$ with respect to $u_{1}, p_{1}, u_{2}$ and $p_{2}$, to justify the adjoint equations written in (34). These computations are similar to the ones made in the proof of Lemma 2.1 and we have

$$
\begin{aligned}
&\left\langle\frac{\partial \tilde{\mathcal{L}}}{\partial u_{1}}, \delta u_{1}\right\rangle_{\left(H^{-1}(\Omega)\right)^{n},\left(H_{0}^{1}(\Omega)\right)^{n}}=\int_{\Gamma}\left(u_{1}-u_{2}\right) \cdot \delta u_{1}+\mu_{1} \int_{\Omega} \nabla \lambda_{1}: \nabla \delta u_{1} \\
&-\int_{\Omega} \pi_{1} \operatorname{div}\left(\delta u_{1}\right)-\mu_{1} \int_{\Omega_{2}} \nabla \lambda_{2}: \nabla \delta u_{1}, \forall \delta u_{1} \in\left(H_{0}^{1}(\Omega)\right)^{n}, \\
&\left\langle\frac{\partial \tilde{\mathcal{L}}}{\partial u_{2}}, \delta u_{2}\right\rangle_{W_{2}^{\prime}, W_{2}}=-\int_{\Gamma}\left(u_{1}-u_{2}\right) \cdot \delta u_{2}+\mu_{2} \int_{\Omega_{2}} \nabla \lambda_{2}: \nabla \delta u_{2} \\
&-\int_{\Omega_{2}} \pi_{2} \operatorname{div}\left(\delta u_{2}\right), \\
& \begin{aligned}
&\left\langle\frac{\partial \tilde{\mathcal{L}}}{\partial p_{1}}, \delta p_{1}\right\rangle_{L^{2}(\Omega), L^{2}(\Omega)}-\int_{\Omega} \delta p_{1} \operatorname{div}\left(\lambda_{1}\right)+\int_{\Omega_{2}} \delta p_{1} \operatorname{div}\left(\lambda_{2}\right), \quad \forall \delta p_{1} \in L_{0}^{2}(\Omega), \\
& \forall \delta u_{2} \in W_{2}, \\
&\left\langle\frac{\partial \tilde{\mathcal{L}}}{\partial p_{2}}, \delta p_{2}\right\rangle_{L^{2}\left(\Omega_{2}\right), L^{2}(\Omega)}=-\int_{\Omega_{2}} \delta p_{2} \operatorname{div}\left(\lambda_{2}\right), \quad \forall \delta p_{2} \in L^{2}\left(\Omega_{2}\right) .
\end{aligned}
\end{aligned}
$$

Then, taking $\left(u_{1}, p_{1}\right)=\left(w_{1}^{g}, p_{1}^{g}\right)$ and $\left(u_{2}, p_{2}\right)=\left(w_{2}^{g}, p_{2}^{g}\right)$, we can deduce that the adjoint
problem associated to the direct problem (30) is,

$$
\begin{cases}\text { find }\left(\lambda_{1}, \pi_{1}\right) \in\left(H_{0}^{1}(\Omega)\right)^{n} \times L_{0}^{2}(\Omega) \text { and }\left(\lambda_{2}, \pi_{2}\right) \in W_{2} \times L^{2}(\Omega) \text { such that }  \tag{51}\\ \left\langle\frac{\partial \tilde{\mathcal{L}}}{\partial u_{1}}, v_{1}\right\rangle_{\left(H^{-1}(\Omega)\right)^{n},\left(H_{0}^{1}(\Omega)\right)^{n}} & =0, \quad \forall v_{1} \in\left(H_{0}^{1}(\Omega)\right)^{n}, \\ \left\langle\frac{\partial \tilde{\mathcal{L}}}{\partial p_{1}}, q_{1}\right\rangle_{L^{2}(\Omega), L^{2}(\Omega)} & =0, \quad \forall q_{1} \in L_{0}^{2}(\Omega), \\ \left\langle\frac{\partial \tilde{\mathcal{L}}}{\partial u_{2}}, v_{2}\right\rangle_{W_{2}^{\prime}, W_{2}} & 0, \quad \forall v_{2} \in W_{2}, \\ \left\langle\frac{\partial \tilde{\mathcal{L}}}{\partial p_{2}}, q_{2}\right\rangle_{L^{2}\left(\Omega_{2}\right), L^{2}(\Omega)} & =0, \quad \forall q_{2} \in L^{2}\left(\Omega_{2}\right) .\end{cases}
$$

Thus, the adjoint problem for the Stokes transmission problem is indeed the weak problem written in (34). Problem (34) consists in two Stokes problems, whose well-posedness derives from well-known results about the Stokes equations (see [10]). We denote $\left(\left(\omega_{1}^{g}, \pi_{1}^{g}\right),\left(\omega_{2}^{g}, \pi_{2}^{g}\right)\right)$ its unique solution in $\left(\left(H_{0}^{1}(\Omega)\right)^{n} \times L_{0}^{2}(\Omega)\right) \times\left(W_{2} \times L^{2}\left(\Omega_{2}\right)\right)$.

The differentiability of $\tilde{J}$ relies on the same arguments that the ones used in the proof of Theorem 2.3. Likewise, taking $\left(u_{1}, p_{1}\right)=\left(w_{1}^{g}, p_{1}^{g}\right)$ and $\left(u_{2}, p_{2}\right)=\left(w_{2}^{g}, p_{2}^{g}\right)$, the Lagragian (33) reduces to

$$
\tilde{\mathcal{L}}\left(g,\left(\left(w_{1}^{g}, p_{1}^{g}\right),\left(w_{2}^{g}, p_{2}^{g}\right)\right),\left(\left(\lambda_{1}, \pi_{1}\right),\left(\lambda_{2}, \pi_{2}\right)\right)\right)=\tilde{J}(g), \quad \forall g \in W_{2}^{\prime}
$$

Differentiating this previous inequality with respect to $g$ using the chain rule and taking

$$
\left(\lambda_{1}, \pi_{1}\right)=\left(\omega_{1}^{g}, \pi_{1}^{g}\right) \text { and }\left(\lambda_{2}, \pi_{2}\right)=\left(\omega_{2}^{g}, \pi_{2}^{g}\right)
$$

we find that, for all $\delta g$ in $W_{2}^{\prime}$,

$$
\langle\nabla \tilde{J}, \delta g\rangle_{W_{2}^{\prime \prime}, W_{2}}=\left\langle\frac{\partial \tilde{\mathcal{L}}}{\partial g}\left(g,\left(\left(w_{1}^{g}, p_{1}^{g}\right),\left(w_{2}^{g}, p_{2}^{g}\right)\right),\left(\left(\omega_{1}^{g}, \pi_{1}^{g}\right),\left(\omega_{2}^{g}, \pi_{2}^{g}\right)\right)\right), \delta g\right\rangle_{W_{2}^{\prime \prime}, W_{2}^{\prime}}
$$

Moreover, the differentiate of $\tilde{\mathcal{L}}$ with respect to $g, \frac{\partial \tilde{\mathcal{L}}}{\partial g} \in W_{2}^{\prime \prime}$, is defined such that, for all $\delta g$ in $W_{2}^{\prime}$,

$$
\left\langle\frac{\partial \tilde{\mathcal{L}}}{\partial g}\left(g,\left(\left(u_{1}, p_{1}\right),\left(u_{2}, p_{2}\right)\right),\left(\left(\lambda_{1}, \pi_{1}\right),\left(\lambda_{2}, \pi_{2}\right)\right)\right), \delta g\right\rangle_{W_{2}^{\prime \prime}, W_{2}^{\prime}}=\left\langle\delta g, \lambda_{2}-\lambda_{1 \mid \Omega_{2}}\right\rangle_{W_{2}^{\prime}, W_{2}}
$$

Finally, the gradient of $\tilde{J}$ is given by,

$$
\langle\nabla \tilde{J}, \delta g\rangle_{W_{2}^{\prime \prime}, W_{2}}=\left\langle\delta g, \lambda_{2}-\lambda_{1 \mid \Omega_{2}}\right\rangle_{W_{2}^{\prime}, W_{2}}, \quad \forall \delta g \in W_{2}^{\prime}
$$

Now, let us show the equivalence between the minimisation of $\tilde{J}$ and the research of a suitable control such that the solution of (30) satisfies the conditions (31).

Proof of Theorem 4.3. The reciprocal statement is straightforward. If, for a given $g$, the solution of (30) satisfies the condition (31) then, $\tilde{J}(g)=0$.

Now, let $g$ be a minimiser of $\tilde{J}$. It follows that one has

$$
\left\langle\delta g, \omega_{2}^{g}-\omega_{1 \mid \Omega_{2}}^{g}\right\rangle_{W_{2}^{\prime}, W_{2}}=0, \quad \forall \delta g \in W_{2}^{\prime},
$$

which means that $\omega_{1 \mid \Omega_{2}}^{g}=\omega_{2}^{g}$ and the couple $\left(\omega_{1 \mid \Omega_{1}}^{g}, \omega_{2}^{g}\right)$ belongs to $\mathcal{W}$. Yet, the functions $\omega_{1}^{g}, \omega_{2}^{g}, \pi_{1}^{g}$ and $\pi_{2}^{g}$ satisfy the equations in problem (34). In particular, for all ( $v_{1}, v_{2}$ ) in $\tilde{\mathcal{W}}$ and for all $\left(q_{1}, q_{2}\right)$ in $L_{0}^{2}(\Omega) \times L^{2}\left(\Omega_{2}\right)$, where

$$
\tilde{\mathcal{W}}=\left\{\left(v_{1}, v_{2}\right) \in\left(H_{0}^{1}(\Omega)\right)^{n} \times W_{2} ; v_{1 \mid \Omega_{2}}=v_{2}\right\},
$$

these equations write,

$$
\begin{align*}
\mu_{2} \int_{\Omega_{2}} \nabla \omega_{2}^{g}: \nabla v_{2}-\int_{\Omega_{2}} \pi_{2}^{g} \operatorname{div}\left(v_{2}\right) & =\int_{\Gamma}\left(w_{1}^{g}-w_{2}^{g}\right) \cdot v_{1}, \\
\int_{\Omega_{2}} q_{2} \operatorname{div}\left(\omega_{2}^{g}\right) & =0, \\
\mu_{1} \int_{\Omega} \nabla \omega_{1}^{g}: \nabla v_{1}-\int_{\Omega} \pi_{1}^{g} \operatorname{div}\left(v_{1}\right) & =-\int_{\Gamma}\left(w_{1}^{g}-w_{2}^{g}\right) \cdot v_{1}+\mu_{1} \int_{\Omega_{2}} \nabla \omega_{2}^{g}: \nabla v_{1 \mid \Omega_{2}}, \\
\int_{\Omega} q_{1} \operatorname{div}\left(\omega_{1}^{g}\right) & =\int_{\Omega_{2}} q_{1 \mid \Omega_{2}} \operatorname{div}\left(\omega_{2}^{g}\right) . \tag{52}
\end{align*}
$$

Summing the first equation in (52) with the third one and summing the second equation with the fourth one and using the fact that $\omega_{1 \mid \Omega_{2}}^{g}=\omega_{2}^{g}$, we find that

$$
\begin{aligned}
\mu_{1} \int_{\Omega_{1}} \nabla \omega_{1 \mid \Omega_{1}}^{g}: \nabla v_{1 \mid \Omega_{1}}-\int_{\Omega_{1}} \pi_{1 \mid \Omega_{1}}^{g} \operatorname{div}\left(v_{1 \mid \Omega_{1}}\right) & +\mu_{2} \quad \int_{\Omega_{2}} \nabla \omega_{2}^{g}: \nabla v_{2}-\int_{\Omega_{2}}\left(\pi_{2}^{g}+\pi_{1 \mid \Omega_{2}}^{g}\right) \operatorname{div}\left(v_{2}\right) \\
& =0, \quad \forall\left(v_{1}, v_{2}\right) \in \tilde{\mathcal{W}}, \\
\int_{\Omega_{1}} q_{1 \mid \Omega_{1}} \operatorname{div}\left(\omega_{1 \mid \Omega_{1}}^{g}\right)+\int_{\Omega_{2}} q_{2} \operatorname{div}\left(\omega_{2}^{g}\right)=0, & \forall\left(q_{1}, q_{2}\right) \in L_{0}^{2}(\Omega) \times L^{2}\left(\Omega_{2}\right) .
\end{aligned}
$$

Moreover, because the test function $v_{1}$ belongs to $\left(H_{0}^{1}(\Omega)\right)^{n}$, we can redefine the pressure $\pi_{1}^{g}$ up to a constant by

$$
\tilde{\pi}_{1}^{g}=\pi_{1}^{g}-\frac{1}{|\Omega|} \int_{\Omega_{2}} \pi_{2}^{g}
$$

such that ( $\tilde{\pi}_{1 \mid \Omega_{1}}^{g}, \pi_{2}^{g}+\tilde{\pi}_{1 \mid \Omega_{2}}^{g}$ ) belongs to the space $Q$. Moreover, let $\left(v_{1}, v_{2}\right)$ be in $\mathcal{W}$. We can construt an extension of $v_{1}$ in the whole space $\left(H_{0}^{1}(\Omega)\right)^{n}$ using $v_{2}$, as in the proof of Theorem 4.1 and we still denote $E_{u} v_{1}$ this extension. Likewise, for all $\left(q_{1}, q_{2}\right)$ in $Q$, we can extend $q_{1}$ in the whole space $L_{0}^{2}(\Omega)$ and we still denote $E_{p} q_{1}$ this extension. Then, the couple ( $E_{u} v_{1}, v_{2}$ ) belongs to $\tilde{\mathcal{W}}$, the couple $\left(E_{p} q_{1}, q_{2}\right)$ belongs to $L_{0}^{2}(\Omega) \times L^{2}\left(\Omega_{2}\right)$ and it follows that $\left(\left(\omega_{1 \mid \Omega_{1}}^{g}, \omega_{2}^{g}\right),\left(\tilde{\pi}_{1 \mid \Omega_{1}}^{g}, \pi_{2}^{g}+\tilde{\pi}_{1 \mid \Omega_{2}}^{g}\right)\right)$ satisfies the equations

$$
\begin{aligned}
\mu_{1} \int_{\Omega_{1}} \nabla \omega_{1 \mid \Omega_{1}}^{g}: \nabla v_{1}-\int_{\Omega_{1}} \tilde{\pi}_{1 \mid \Omega_{1}}^{g} \operatorname{div}\left(v_{1}\right) & +\mu_{2} \quad \int_{\Omega_{2}} \nabla \omega_{2}^{g}: \nabla v_{2}-\int_{\Omega_{2}}\left(\pi_{2}^{g}+\tilde{\pi}_{1 \mid \Omega_{2}}^{g}\right) \operatorname{div}\left(v_{2}\right) \\
& =0, \quad \forall\left(v_{1}, v_{2}\right) \in \mathcal{W}, \\
\int_{\Omega_{1}} q_{1} \operatorname{div}\left(\omega_{1 \mid \Omega_{1}}^{g}\right)+\int_{\Omega_{2}} q_{2} \operatorname{div}\left(\omega_{2}^{g}\right) & =0, \quad \forall\left(q_{1}, q_{2}\right) \in Q .
\end{aligned}
$$

We conclude that $\left(\left(\omega_{1 \mid \Omega_{1}}^{g}, \omega_{2}^{g}\right),\left(\tilde{\pi}_{1 \mid \Omega_{1}}^{g}, \pi_{2}^{g}+\tilde{\pi}_{1 \mid \Omega_{2}}^{g}\right)\right)$ is solution of a Stokes problem similar to (26) but with no external force and, thus, is the zero of $\mathcal{W} \times Q$. Then, the first equation in (34) becomes

$$
\int_{\Gamma}\left(w_{1}^{g}-w_{2}^{g}\right) \cdot v_{2}=0, \quad \forall v_{2} \in \tilde{W}_{2}
$$

where

$$
\tilde{W}_{2}=\left\{v \in W_{2} ; \operatorname{div}(v)=0\right\} .
$$

Besides, $w_{1 \mid \Gamma}^{g}$ and $w_{2 \mid \Gamma}^{g}$ belongs to the space

$$
W_{\Gamma}=\left\{v \in H_{00}^{1 / 2}(\Gamma) ; \int_{\Gamma} v=0\right\}
$$

because $\operatorname{div}\left(w_{1 \mid \Omega_{1}}^{g}\right)=0$ and $\operatorname{div}\left(w_{2}^{g}\right)=0$, and for all $v_{2}$ in $W_{\Gamma}$ we can construct an extension of $v_{2}$ in the whole space $\tilde{W}_{2}$, according to Bogovskii's result in [8]. In particular, taking $v_{2}=w_{1 \mid \Omega_{1}}^{g}-w_{2}^{g}$ in $W_{\Gamma}$, it follows that

$$
\int_{\Gamma}\left|w_{1}^{g}-w_{2}^{g}\right|^{2}=0
$$

which means that $w_{1 \mid \Gamma}^{g}=w_{2 \mid \Gamma}^{g}$ and the equality (31) is satisfied.

## B Proofs of theorems related to the fluid-structure interaction problem

This appendix is dedicated to the proofs of all results stated in subsection 4.2. In particular, we are interested in proving the well-posedness of Problem (39), in showing the existence of the control $g$, in giving an explicit formula for the gradient of $J_{k}$ and, finally, in proving the equivalence between the minimisation of $J_{k}$ and the resolution of the fluid-structure interaction problem.

We follow the order of the previous enumeration and start with the well-posedness of Problem (39). To study Problem (39), we do a change in variable on the displacement $d_{s}^{k}$, in order to work on a velocity-velocity formulation of the fluid-structure problem. We introduce the velocity of the structure at time $t_{k}$,

$$
u_{s}^{k}=\frac{1}{\delta t}\left(d_{s}^{k}-d_{s}^{k-1}\right)
$$

Because Problem (39) is linear, it is completely equivalent to the problem where $d_{s}^{k}$ has been replaced with $\delta u_{s}^{k}+d_{s}^{k-1}$ :

$$
\left\{\begin{array}{cl}
\text { find }\left(u_{f}, u_{s}\right) \in W_{u} \text { and } p \in L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right) \text { such that }  \tag{53}\\
a\left(\left(u_{f}, u_{s}\right),\left(v_{f}, v_{s}\right)\right)-\left(B\left(v_{f}, v_{s}\right), p\right)_{L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)}=L\left(v_{f}, v_{s}\right), & \forall\left(v_{f}, v_{s}\right) \in W_{u}, \\
\left(B\left(u_{f}, u_{s}\right), q\right)_{L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)}=0, & \forall q \in L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right),
\end{array}\right.
$$

where $(\cdot, \cdot)_{L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)}$ denotes the scalar product in $L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)$ and $a, L$ and $B$ are defined by

$$
\begin{align*}
a\left(\left(u_{f}, u_{s}\right),\left(v_{f}, v_{s}\right)\right)= & 2 \mu_{f} \int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} D\left(u_{f}\right): D\left(v_{f}\right)+2 \delta_{t} \mu_{s} \int_{\Omega_{s}} D\left(u_{s}\right): D\left(v_{s}\right) \\
& +\delta_{t} \lambda_{s} \int_{\Omega_{s}} \operatorname{div}\left(u_{s}\right) \operatorname{div}\left(v_{s}\right), \quad \forall\left(u_{f}, u_{s}\right),\left(v_{f}, f_{s}\right) \in W_{u}  \tag{54}\\
L\left(v_{f}, v_{s}\right)= & -2 \mu_{s} \int_{\Omega_{s}} D\left(d_{s}^{k-1}\right): D\left(v_{s}\right)-\lambda_{s} \int_{\Omega_{s}} \operatorname{div}\left(d_{s}^{k-1}\right) \operatorname{div}\left(v_{s}\right) \\
& +\int_{\Phi_{t_{k}\left(\Omega_{f}\right)}} f_{f} \cdot v_{f}+\int_{\Omega_{s}} f_{s} \cdot v_{s}, \\
B\left(v_{f}, v_{s}\right)= & \forall\left(v_{f}, v_{s}\right) \in W_{u} \\
\operatorname{div}\left(v_{f}\right), & \forall\left(v_{f}, v_{s}\right) \in W_{u} .
\end{align*}
$$

This is a saddle-point problem, whose well-posedness can be proved using results from [11]. Using the well-known Cauchy-Schwartz and Korn inequalities, we show that the bilinear form $a$ is continuous and coercive and that the linear form $L$ is continuous. The operator $B$ is also linear and continuous. It remains to show that $B$ is surjective form $W_{u}$ to $L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)$. Indeed, let $q$ be in $L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)$. We can easily extend $q$ in the whole space $L_{0}^{2}(\Omega)$; it is sufficient to take

$$
E_{p} q=\left\{\begin{array}{rll}
q & \text { in } & \Phi_{t_{k}}\left(\Omega_{f}\right)  \tag{55}\\
\frac{-1}{\left|\Phi_{t_{k}}\left(\Omega_{s}\right)\right|} \int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} q & \text { in } & \Phi_{t_{k}}\left(\Omega_{s}\right) .
\end{array}\right.
$$

Then, $E_{p} q$ belongs to $L_{0}^{2}(\Omega)$ and Bogovskii's result in [8] implies that there exists $\tilde{u}$ in $H_{0}^{1}(\Omega)$ such that $\operatorname{div}(\tilde{u})=E_{p} q$. Finally, we define $v_{f}=\tilde{u}_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}$ and $v_{s}=\tilde{u}_{\mid \Phi_{t_{k}}\left(\Omega_{s}\right)}$ and it follows that the couple $\left(v_{f}, v_{s} \circ \Phi_{t_{k}}^{-1}\right)$ belongs to $W_{u}$ and satisfies $B\left(v_{f}, v_{s} \circ \Phi_{t_{k}}^{-1}\right)=$ $\operatorname{div}\left(v_{f}\right)=q$. This proves the surjectivity of the operator $B$. Thus, problem (53), and consequently problem (39), are well-posed.

We now give a proof of Theorem 4.4.
Proof of Theorem 4.4. We can construct two extension operators, still denoted $E_{u}$ and $E_{p}$, that extend $u_{f}^{k}$ into the whole space $\left(H_{0}^{1}(\Omega)\right)^{n}$ (and stays divergence-free) and $p_{f}^{k}$ into the whole space $L_{0}^{2}(\Omega)$. Indeed, for the pressure, consider the operator defined in (55). Then, the extension of the pressure $E_{p} p_{f}^{k}$ belongs to $L_{0}^{2}(\Omega)$. Now, because the fluid is incompressible, we have that

$$
0=\int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} \operatorname{div}\left(u_{f}^{k}\right)=\int_{\Phi_{t_{k}}(\Gamma)} u_{f}^{k} \cdot n_{f}^{k}
$$

Thus, using one more time Bogovskii's result, there exists $\tilde{u}_{f}$ in $H^{1}\left(\Phi_{t_{k}}\left(\Omega_{s}\right)\right)$ such that

$$
\operatorname{div}\left(\tilde{u}_{f}\right)=0, \quad \tilde{u}_{f \mid \Phi_{t_{k}}(\Gamma)}=u_{f \mid \Phi_{t_{k}}(\Gamma)}^{k} \text { and } \tilde{u}_{f \mid \Phi_{t_{k}}\left(\Gamma_{s}\right)}=0
$$

Then, the extension of the fluid velocity, $E_{u} u_{f}^{k}$, is defined such that

$$
E_{u} u_{f}^{k}=\left\{\begin{array}{ccc}
u_{f}^{k} & \text { in } & \Phi_{t_{k}}\left(\Omega_{f}\right) \\
\tilde{u}_{f} & \text { in } & \Phi_{t_{k}}\left(\Omega_{s}\right)
\end{array} .\right.
$$

The extension $E_{u} u_{f}^{k}$ belongs to $\left(H_{0}^{1}(\Omega)\right)^{n}$, is divergence-free and satisfies the equality

$$
E_{u} u_{f}^{k} \circ \Phi_{t_{k}}=\frac{1}{\delta t}\left(d_{s}^{k}-d_{s}^{k-1}\right) \text { on } \Gamma
$$

Furthermore, $\sigma_{f}\left(E_{u} u_{f}^{k}, E_{p} p_{f}^{k}\right)$ belongs to $\left(L^{2}(\Omega)\right)^{n \times n}$.
Now, we construct a suitable control $g$ in $\left(V_{s}^{k}\right)^{\prime}$ such that

$$
\begin{align*}
\langle g, v\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}}= & \int_{\Phi_{t_{k}}\left(\Omega_{s}\right)} \sigma_{f}\left(E_{u} u_{f}^{k}, E_{p} p_{f}^{k}\right): \nabla v  \tag{56}\\
& +\left\langle\gamma_{n_{f}^{k}}\left(\sigma_{f}\left(E_{u} u_{f}^{k}, E_{p} p_{f}^{k}\right)\right), v\right\rangle_{\left(\Upsilon^{k}\right)^{\prime}, \Upsilon^{k}}, \quad \forall v \in V_{s}^{k}
\end{align*}
$$

where $n_{f}^{k}$ is the unit exterior normal vector to $\partial \Phi_{t_{k}}\left(\Omega_{f}\right)$. Using the Stokes formula (40) and the definition (56), it follows that the extensions $E_{u} u_{f}^{k}$ and $E_{p} p_{f}^{k}$ satisfy

$$
\begin{aligned}
\int_{\Omega} \sigma_{f}\left(E_{u} u_{f}^{k}, E_{p} p_{f}^{k}\right): \nabla v_{f}= & \int_{\Phi_{t_{k}\left(\Omega_{f}\right)}} f_{f}^{k} \cdot v_{f \mid \Phi_{t_{k}}\left(\Omega_{f}\right)} & \\
& +\left\langle g, v_{f \mid \Phi_{t_{k}}\left(\Omega_{s}\right)}\right\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}}, & \forall v_{f} \in\left(H_{0}^{1}(\Omega)\right)^{n}, \\
\int_{\Omega} q \operatorname{div}\left(E_{u} u_{f}^{k}\right)= & 0, & \forall q \in L_{0}^{2}(\Omega)
\end{aligned}
$$

Similarly, using (40), (56) and the weak transmission condition (41), $d_{s}^{k}$ satisfies

$$
\begin{aligned}
\int_{\Omega_{s}} \sigma_{s}\left(d_{s}^{k}\right): \nabla v_{s}= & \int_{\Omega_{s}} f_{s}^{k} \cdot v_{s}-\left\langle g, v_{s} \circ \Phi_{t_{k}-1}^{-1}\right\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}} \\
& +\int_{\Omega_{s}} \Pi_{f}\left(E_{u} u_{f}^{k} \circ \Phi_{t_{k}}, E_{p} p_{f}^{k} \circ \Phi_{t_{k}}\right): \nabla v_{s}, \quad \forall v_{s} \in V_{s} .
\end{aligned}
$$

Finally, we conclude that $\left(E_{u} u_{f}^{k}, E_{p} p_{f}^{k}, d_{s}^{k}\right)$ is the solution of problem (43). Thus,

$$
\left(u^{g}, p^{g}, d^{g}\right)=\left(E_{u} u_{f}^{k}, E_{p} p_{f}^{k}, d_{s}^{k}\right)
$$

and condition (44) is satisfied by construction. This proves the first part of the theorem.
Now, suppose that $g$ is a control such that the equality (44) is satisfied. In particular, this implies that the couple $\left(u_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, d^{g}\right)$ belongs to the space $W_{d}$. Let us define the following Hilbert space:

$$
\tilde{W}_{u}=\left\{\left(v_{f}, v_{s}\right) \in\left(H_{0}^{1}(\Omega)\right)^{n} \times V_{s} ;\left(v_{f} \circ \Phi_{t_{k}}\right)_{\Omega_{s}}=v_{s}\right\} .
$$

As the unique solution of problem (43), the triplet $\left(u^{g}, p^{g}, d^{g}\right)$ satisfies, in particular for all $\left(v_{f}, v_{s}\right)$ in $\tilde{W}_{u}$ and for all $q$ in $L_{0}^{2}(\Omega)$, the equations

$$
\begin{align*}
\int_{\Omega} \sigma_{f}\left(u^{g}, p^{g}\right): \nabla v_{f} & =\int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} f_{f} \cdot v_{f}+\left\langle g, v_{f}\right\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}}, \\
\int_{\Omega} q \operatorname{div}\left(u^{g}\right) & =0,  \tag{57}\\
\int_{\Omega_{s}} \sigma_{s}\left(d^{g}\right): \nabla v_{s} & =\int_{\Omega_{s}} f_{s} \cdot v_{s}-\left\langle g, v_{f}\right\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}}+\int_{\Phi_{t_{k}\left(\Omega_{s}\right)}} \sigma_{f}\left(u^{g}, p^{g}\right): \nabla v_{f} .
\end{align*}
$$

Then adding the first and the third equations in (57) and using the fact that $u^{g}$ is divergence-free, it follows that the triplet $\left(u_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, p_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, d^{g}\right)$ satisfies

$$
\begin{array}{rlrl}
\int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} \sigma_{f}\left(u^{g}, p^{g}\right): & \nabla v_{f}+\int_{\Omega_{s}} \sigma_{s}\left(d^{g}\right): \nabla v_{s} & \\
& =\int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} f_{f} \cdot v_{f}+\int_{\Omega_{s}} f_{s} \cdot v_{s}, \quad \forall\left(v_{f}, v_{s}\right) \in \tilde{W}_{u}, \\
\int_{\Phi_{t_{k}\left(\Omega_{f}\right)}} q \operatorname{div}\left(u^{g}\right) & =0, & \forall q \in L_{0}^{2}(\Omega) .
\end{array}
$$

Moreover, for all $\left(v_{f}, v_{s}\right)$ in $W_{u}$, we can extend $v_{f}$ in the whole space $\left(H_{0}^{1}(\Omega)\right)^{n}$ as we did it for $u_{f}^{k}$ and we denote $E_{u} v_{f}$ this extension. Similarly, for all $q$ in $L_{0}^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right)$, we extend $q$ in the whole space $L_{0}^{2}(\Omega)$ and we denote $E_{p} q$ this extension. Then, the couple $\left(E_{u} v_{f}, v_{s}\right)$ belongs to $\tilde{W}_{u}, E_{p} q$ belongs to $L_{0}^{2}(\Omega)$ and, finally, the triplet $\left(u_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, p_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, d^{g}\right)$ satisfies the equations

$$
\begin{array}{rlrl}
\int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} \sigma_{f}\left(u_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, p_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}\right): \nabla v_{f}+\int_{\Omega_{s}} \sigma_{s}\left(d^{g}\right): \nabla v_{s} & \\
& =\int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} f_{f} \cdot v_{f}+\int_{\Omega_{s}} f_{s} \cdot v_{s}, & \forall\left(v_{f}, v_{s}\right) \in W_{u} \\
\int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} q \operatorname{div}\left(u_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}\right) & =0, & \forall q \in L^{2}\left(\Phi_{t_{k}}\left(\Omega_{f}\right)\right) .
\end{array}
$$

Thus, the triplet $\left(u_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, p_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, d^{g}\right)$ is the unique solution of the initial fluid-structure problem (39), which proves that we can recover the solution of the fluid-structure interaction problem from the solution of its smooth extension formulation.

Since there exists at least one suitable control $g$ for the smooth extension problem, we can hope to obtain it by a minimisation process on the function $J_{k}$. This is actually possible according to Theorem 4.5, which states that $J_{k}$ is differentiable with respect to $g$ and gives a characterisation of its gradient. Here, we prove this result.

Proof of Theorem 4.5. On the first hand, the differentiation of the Lagrangian function $\mathcal{L}_{k}$, defined in (46), with respect to $u, p, d$ and $g$ follows the same process that the one explained for the Laplace and Stokes transmission problems. Here, we only give the expressions of these different derivatives:

$$
\begin{aligned}
\left\langle\frac{\mathcal{L}_{k}}{\partial u}, \delta u\right\rangle_{H^{-1}(\Omega), H_{0}^{1}(\Omega)}= & \int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} \sigma_{f}\left(\lambda_{f}, \pi\right): \nabla \delta u \\
& -2 \mu_{f} \int_{\Phi_{t_{k}}\left(\Omega_{s}\right)} D\left(\nu_{s} \circ \Phi_{t_{k}}^{-1}\right): D(\delta u) \\
& +\int_{\Gamma}\left(u \circ \Phi_{t_{k}}-\frac{1}{\delta t}\left(d-d_{d}^{k-1}\right)\right) \cdot\left(\delta u \circ \Phi_{t_{k}}\right), \quad \forall \delta u \in\left(H_{0}^{1}(\Omega)\right)^{n}, \\
= & -\int_{\Omega} \delta p \operatorname{div}\left(\lambda_{f}\right)+\int_{\Omega_{s}} \delta p \operatorname{div}\left(\nu_{s} \circ \Phi_{t_{k}}^{-1}\right), \quad \forall \delta p \in L_{0}^{2}(\Omega), \\
\left\langle\frac{\mathcal{L}_{k}}{\partial p}, \delta p\right\rangle_{L^{2}(\Omega), L^{2}(\Omega)}= & \int_{\Omega_{s}} \sigma_{s}\left(\nu_{s}\right): \nabla \delta d \\
\left\langle\frac{\mathcal{L}_{k}}{\partial d}, \delta d\right\rangle_{V_{s}^{\prime}, V_{s}}= & -\frac{1}{\delta t} \int_{\Gamma}\left(u \circ \Phi_{t_{k}}-\frac{1}{\delta t}\left(d-d_{s}^{k-1}\right)\right) \cdot \delta d, \quad \forall \delta d \in V_{s}, \\
\left\langle\frac{\mathcal{L}_{k}}{\partial g}, \delta g\right\rangle_{V_{s}^{\prime \prime}, k}, V_{s, k}^{\prime}= & \left\langle\delta g, \nu_{s} \circ \Phi_{t_{k}}^{-1}-\lambda_{f\left|\Phi_{t}\left(\Omega_{s}\right)\right\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}}, \quad \forall \delta g \in V_{s, k}^{\prime} .}=\right.
\end{aligned}
$$

They enable us, in particular, to recover the adjoint equations, written in (47). Problem (47) consists in a linear elasticity problem and a Stokes problem, whose well-posedness derives from the same arguments that we already used. We denote by $\left(\nu^{g}, \lambda^{g}, \pi^{g}\right)$ its unique solution.

On the other hand, the differentiability of $J_{k}$ relies on the same arguments that the ones used in the proof of Theorem 2.3 and the fact that the transformation $\Phi_{t_{k}}$ is sufficiently regular. Then, replacing $(u, p, d)$ in the Lagrangian function (46), by the solution ( $u^{g}, p^{g}, d^{g}$ ) of the smooth extension problem (43), the Lagrangian (46) reduces to

$$
\mathcal{L}_{k}\left(g,\left(u^{g}, p^{g}, d^{g}\right),\left(\lambda_{f}, \pi, \nu_{s}\right)\right)=J_{k}(g), \quad \forall g \in\left(V_{s}^{k}\right)^{\prime} .
$$

Differentiating with respect to $g$ using the chain rule and replacing the triplet $\left(\lambda_{f}, \pi, \nu_{s}\right)$ by the solution ( $\lambda^{g}, \pi^{g}, \nu^{g}$ ) of the adjoint problem (47), the gradient of $J_{k}$ is finally given by,

$$
\left\langle\nabla J_{k}(g), \delta g\right\rangle_{\left(V_{s}^{k}\right)^{\prime \prime},\left(V_{s}^{k}\right)^{\prime}}=\left\langle\delta g, \nu_{k}^{g} \circ \Phi_{t_{k}}^{-1}-\lambda_{k \mid \Phi_{t_{k}}\left(\Omega_{s}\right)}^{g}\right\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}}, \quad \forall \delta g \in\left(V_{s}^{k}\right)^{\prime} .
$$

With this explicit expression of the gradient of $J_{k}$, we can now state the equivalence between the research of a suitable control such that the solution of the smooth extension problem (43) satisfies the condition (44) and the minimisation of $J_{k}$. This is the result of Theorem 4.6, that we prove in the following.

Proof of Theorem 4.6. On the one hand, if for a given $g$ in $\left(V_{s}^{k}\right)^{\prime}$, the solution $\left(u^{g}, p^{g}, d^{g}\right)$ of the smooth extension problem (43) satisfies the condition (44), then

$$
J_{k}(g)=0 .
$$

On the other hand, suppose that $g$ is a minimiser of $J_{k}$. Then, the adjoints $\lambda^{g}$ and $\nu^{g}$ satisfy the equality,

$$
\left\langle\delta g, \nu^{g} \circ \Phi_{t_{k}}^{-1}-\lambda_{\mid \Phi_{t_{k}}\left(\Omega_{s}\right)}^{g}\right\rangle_{\left(V_{s}^{k}\right)^{\prime}, V_{s}^{k}}=0, \quad \forall \delta g \in\left(V_{s}^{k}\right)^{\prime},
$$

which corresponds to the fact that $g$ is a zero for the gradient of $J_{k}$. It follows, in particular, that,

$$
\begin{equation*}
\nu^{g} \circ \Phi_{t_{k}}^{-1}=\lambda_{\mid \Phi_{t_{k}}\left(\Omega_{s}\right)}^{g}, \tag{58}
\end{equation*}
$$

and the couple $\left(\lambda_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, \nu^{g}\right)$ belongs to the space $W_{u}$. Yet, $\lambda^{g}, \pi^{g}$ and $\nu^{g}$ satisfy the adjoint equations in (47). In particular, for all $\left(v_{f}, v_{s}\right)$ in $\hat{W}_{u}$ and for all $q$ in $L_{0}^{2}(\Omega)$, where

$$
\hat{W}_{u}=\left\{\left(v_{f}, v_{s}\right) \in\left(H_{0}^{1}(\Omega)\right)^{n} \times V_{s} ; \operatorname{div}\left(v_{f}\right)=0,\left(v_{f} \circ \Phi_{t_{k}}\right)_{\mid \Omega_{s}}=v_{s}\right\},
$$

these equations write

$$
\begin{gather*}
\mu_{s} \int_{\Omega_{s}} \sigma_{s}\left(\nu^{g}\right): \nabla v_{s}=\frac{1}{\delta t} \int_{\Gamma}\left(u^{g} \circ \Phi_{t_{k}}-\frac{1}{\delta t}\left(d^{g}-d_{s}^{k-1}\right)\right) \cdot v_{s} \\
2 \mu_{f} \int_{\Phi_{t_{k}}\left(\Omega_{f}\right)} D\left(\lambda^{g}\right): D\left(v_{f}\right)-\int_{\Omega} \pi^{g} \operatorname{div}\left(v_{f}\right)=-\int_{\Gamma}\left(u^{g} \circ \Phi_{t_{k}}-\frac{1}{\delta t}\left(d^{g}-d_{s}^{k-1}\right)\right) \cdot v_{s}  \tag{59}\\
\int_{\Phi_{t_{k}\left(\Omega_{f}\right)}} q \operatorname{div}\left(\lambda^{g}\right)=0
\end{gather*}
$$

Multiplying the first equation in (59) by $\delta t$ and summing it with the second one, we obtain that,

$$
\begin{equation*}
2 \mu_{f} \int_{\Phi_{t_{k}\left(\Omega_{f}\right)}} D\left(\lambda_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}\right): D\left(v_{f}\right)+\delta t \mu_{s} \int_{\Omega_{s}} \sigma_{s}\left(\nu^{g}\right): \nabla v_{s}=0, \quad \forall\left(v_{f}, v_{s}\right) \in \hat{W}_{u} . \tag{60}
\end{equation*}
$$

Moreover, let $\left(v_{f}, v_{s}\right)$ be in $W_{u}$. We can construct an extension of $v_{f}$ in the whole space $\left(H_{0}^{1}(\Omega)\right)^{n}$, denoted $E_{u} v_{f}$, such that $\operatorname{div}\left(E_{u} v_{f}\right)=0$ and $\left(\left(E_{u} v_{f}\right) \circ \Phi_{t_{k}}\right)_{\mid \Gamma}=v_{s \mid \Gamma}$, as we did it in the proof of Theorem 4.4. Then, the couple ( $E_{u} v_{f}, v_{s}$ ) belongs to $\hat{W}_{u}$ and it follows that the couple $\left(\lambda_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, \nu^{g}\right)$ satisfies the equation

$$
2 \mu_{f} \int_{\Phi_{t_{k}\left(\Omega_{f}\right)}} D\left(\lambda_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}\right): D\left(v_{f}\right)+\delta t \mu_{s} \int_{\Omega_{s}} \sigma_{s}\left(\nu^{g}\right): \nabla v_{s}=0, \quad \forall\left(v_{f}, v_{s}\right) \in W_{u} .
$$

We conclude that $\left(\lambda_{\mid \Phi_{t_{k}}\left(\Omega_{f}\right)}^{g}, \nu^{g}\right)$ is solution of the problem

$$
\left\{\begin{array}{c}
\text { find }(\lambda, \nu) \text { in } W_{u} \text { such that, }  \tag{61}\\
a\left((\lambda, \nu),\left(v_{f}, v_{s}\right)\right)=0, \quad \forall\left(v_{f}, v_{s}\right) \in W_{u},
\end{array}\right.
$$

where the bilinear continuous and coercive form $a$ has been defined in (54), which admits the zero of $W_{u}$ as unique solution. Then, the first equation in (47) becomes

$$
\int_{\Gamma}\left(u^{g} \circ \Phi_{t_{k}}-\frac{1}{\delta_{t}}\left(d^{g}-d_{s}^{k-1}\right)\right) \cdot v_{s}=0, \quad \forall v_{s} \in V_{s}
$$

Moreover, for all $v_{s}$ in $\left(H_{00}^{1 / 2}(\Gamma)\right)^{n}$ we can construct an extension of $v_{s}$ in the whole space $V_{s}$. In particular, taking $v_{s}=\left(u^{g} \circ \Phi_{t_{k}}\right)_{\mid \Gamma}-\frac{1}{\delta_{t}}\left(d_{\mid \Gamma}^{g}-d_{s \mid \Gamma}^{k-1}\right)$, it follows that

$$
\int_{\Gamma}\left|\left(u^{g} \circ \Phi_{t_{k}}\right)_{\mid \Gamma}-\frac{1}{\delta_{t}}\left(d_{\mid \Gamma}^{g}-d_{s \mid \Gamma}^{k-1}\right)\right|^{2}=0,
$$

which means that $\left(u^{g} \circ \Phi_{t_{k}}\right)_{\mid \Gamma}=\frac{1}{\delta_{t}}\left(d_{\mid \Gamma}^{g}-d_{s \mid \Gamma}^{k-1}\right)$ and the equality (44) is satisfied.

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