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# OPTIMIZED SCHWARZ WAVEFORM RELAXATION METHOD FOR INCOMPRESSIBLE STOKES PROBLEM * 

DUC-QUANG BUI $\ddagger$ CAROLINE JAPHET ${ }^{\ddagger}$ AND PASCAL OMNES ${ }^{\ddagger} \ddagger$


#### Abstract

We propose and analyse the optimized Schwarz waveform relaxation (OSWR) method for the unsteady incompressible Stokes equations. Well-posedness of the local subdomain problems with Robin boundary conditions is proved. Convergence of the velocity is shown through energy estimates; however, pressure converges only up to constant values in the subdomains, and an astute correction technique is proposed to recover these constants from the velocity. The convergence factor of the OSWR algorithm is obtained through a Fourier analysis, and allows to efficiently optimize the space-time Robin transmission conditions involved in the OSWR method. Then, numerical illustrations for the two-dimensional unsteady incompressible Stokes system are presented to illustrate the performance of the OSWR algorithm.


Key words. Unsteady incompressible Stokes system, space-time domain decomposition, optimized Schwarz waveform relaxation, Robin transmission conditions, correction technique for the pressure.

AMS subject classifications. 65M55, 35K45, 76D07, 65M12, 65M22, 65B99.

1. Introduction. The study of physical phenomena, whether natural or industrial, is frequently based on numerical simulations involving an increasing number of degrees of freedom. This growing complexity may require the use of resolution techniques which on the one hand are suitable for parallel computing architectures, and on the other hand allow local space and time stepping adapted to the physics, such as spacetime domain decomposition (DD) methods. In this article we are concerned with such methods, with Robin transmission conditions at the interfaces between subdomains, for solving applications related to incompressible fluid mechanics, that are modelled by the unsteady (Navier)-Stokes system.

The well-posedness of such systems with Robin conditions (without domain decomposition) has been the subject of several works in the steady case, see e.g. [47] for the Stokes problem (where the Robin condition is expressed with the symmetric part of the velocity gradient, instead of the gradient), references [45, 38] for the Oseen and Navier-Stokes systems, and [16] for the Stokes-Darcy Coupling. On the other hand, there are few works in the unsteady case; in [39] existence and uniqueness of a solution with a time-dependent Robin boundary condition of the type curl $\mathbf{u} \times \mathbf{n}=\beta(t) \mathbf{u}$ is addressed. In [29] the Stokes problem with Robin conditions is studied, in the context of a global-in-time DD method applied the coupled nonlinear Stokes and Darcy Flows. The well-posedness is not shown.

In this article we study the well-posedness of the unsteady incompressible Stokes system with Robin boundary conditions of type $\alpha\left(\nu \partial_{\mathbf{n}} \mathbf{u} \cdot \mathbf{n}-p\right)+\mathbf{u} \cdot \mathbf{n}=g(t)$ and $\beta \nu \partial_{\mathbf{n}} \mathbf{u} \times \mathbf{n}+\mathbf{u} \times \mathbf{n}=\xi(t)$, in the context of space-time DD methods.

Concerning the DD approaches with Robin conditions, several studies have been carried out for the incompressible (Navier)-Stokes equations : in [41, 42, 34, 43, 40] the steady Oseen equation (and its application to the non-stationary Navier-Stokes equa-

[^0]tions, using a spatial DD at each time step) is considered. More precisely, in [42, 34, 43] a stabilized finite element approximation is proposed (with non-standard Robin conditions due to the stabilization). The convergence of the DD method is proven for the velocity. For the pressure, the convergence is proven when the original monodomain problem involves Robin boundary conditions on a part of the physical boundary. However, the authors point out that for an Oseen problem with Dirichlet conditions on the whole physical boundary, the pressure of the Robin-Robin DD algorithm will converge up to a constant which can differ for different subdomains. This important observation is also mentioned in [11] for the steady Stokes problem, where the DD method is based on a penalty term on the interface (in that case the Robin conditions are not equivalent to the physical ones). The convergence is shown for a modified pressure in the two-subdomains case. This issue of pressure converging up to a constant that depends on the subdomains is also raised in [33, 23] for the discrete Schwarz algorithm with a DDFV scheme applied to the semi-discrete in time Navier-Stokes system. In $[12,6]$, an optimized Schwarz DD method is studied, and applied at each time step to the semi-discrete in time Navier-Stokes equations. Other transmission conditions (Dirichlet / Neumann) are considered e.g. in [46, 21, 44, 49] for Stokes and Navier-Stokes equations.

In this article we consider global-in-time Schwarz methods which use waveform relaxation techniques, i.e. Schwarz waveform relaxation (SWR). Such iterative methods use computations in the subdomains over the whole time interval, exchanging space-time boundary data through transmission conditions on the space-time interfaces. The main advantage is that space-time discretizations can be chosen independently on each subdomain, and, at the end of each iteration, only a small amount of information is exchanged, which makes the parallelization (in space and time) very efficient.

The space-time boundary data play an important role in the convergence process and can be of Dirichlet [20, 22], absorbing, Robin (or Ventcell) type [19, 35, 4, 25, 24]. The value of the Robin (or Ventcell) parameters can be optimized to improve convergence rates (see [19, 30, 35, 32]), and the corresponding method is called optimized Schwarz waveform relaxation (OSWR). This method is wildly used and analyzed for fluid dynamics, see references above, and e.g. [35, 18, 36, 3, 5, 28, 1, 48].

For the application of the SWR method on the Navier-Stokes equations, we are aware of the article [3] where an OSWR method is proposed for the rotating 3D incompressible hydrostatic Navier-Stokes equations with free surface. However, the hydrostatic nature of the model modifies the structure of the continuity equation which now involves a transport term for the free surface (which plays the same role as the pressure in the momentum equation of the standard Navier-Stokes system), so that the results in [3] cannot apply to the problem considered in the present work. In [12], an SWR method for the Oseen equations is studied; optimal transparent boundary conditions are derived, and local approximations for these nonlocal conditions are proposed. No general convergence analysis of the resulting algorithm (e.g. via energy estimates) is given. A convergence factor is obtained in the idealized case of two half-space subdomains and unbounded time interval, via Laplace-Fourier transforms.

Concerning the compressible Euler and Navier-Stokes equations, in [14, 13] an SWR method is proposed and various numerical experiments are shown.

However, until now, there exists no convergence proof (for SWR or OSWR) for the incompressible Navier-Stokes equations. We contribute to the understanding of the behaviour of the OSWR method by attacking representative, though simpler, model problems. To begin with, we analyze the method on the evolutionary Stokes equations,
a simplified version of the evolutionary Navier-Stokes system in which the convection is simply discarded. The convergence analysis of the velocity iterates involved in the OSWR method, for the Stokes equations, can be performed in a similar manner as for parabolic equations. An extension of this analysis to the evolutionary Oseen equations (a linearization of the Navier-Stokes equations in which the convective velocity field is considered as a given datum) is given in [9]. However, the convergence analysis of the OSWR method has its own obstacle related to the pressure converging only up to constants in the various subdomains, as discussed above. A second purpose of this article is to propose a new technique, in the multidomain case, to recover the pressure from the velocity (at any iteration).

A third purpose of this article is to discuss the choice of the Robin parameters, which play a crucial role in the optimization of the convergence rate. Until recently, the common practice was to derive and optimize a convergence rate in the idealized case of two half-space subdomains and unbounded time interval, via Laplace-Fourier transforms performed on the continuous model (i.e. without taking into account the actual discretization method). We first follow this standard approach in this work, but in a second step modify it to also include the effect of the discretization in the time direction; the Robin parameters obtained with such a modification improve the convergence rate over the standard choice in our numerical tests. Note that studying the influence of the numerical scheme over the OSWR convergence rate is a recent approach, pursued for example in [15, 26, 2].

The remainder of this article is organized as follows. In section 2, we present the model problem and its multidomain form. Since the multi-domain formulation involves local Stokes problems with Robin boundary conditions, we prove the wellposedness of such problems in Section 3. Next, section 4 is dedicated to the algorithm. In section 5 we show that, in general, the pressure calculated by the OSWR algorithm will not converge to the monodomain solution. In section 6 , we obtain a convergence result on the velocity through an energy estimate, and in section 7 , we propose an astute technique to recover the pressure from the velocity. In section 8 , a Fourier analysis is done to get a formulation for the convergence factor of the OSWR algorithm. In section 9, an optimization procedure (based on the convergence factor of the method), that allows to obtain efficient Robin parameters, is given. Then, numerical illustrations for the unsteady Stokes system follow in section 10.
2. Presentation of the model and multidomain formulation. For a bounded domain $\Omega \subseteq \mathbb{R}^{2}$, and for a given viscosity coefficient $\nu>0$ that we suppose constant and uniform, for given initial condition $\mathbf{u}_{0}$ and source term $\mathbf{f}$, we denote respectively by $\mathbf{u}, p$ the velocity and pressure unknowns in the incompressible non-stationary Stokes system:

$$
\begin{array}{rlll}
\partial_{t} \mathbf{u}-\nu \Delta \mathbf{u}+\nabla p & =\mathbf{f} & \text { in } \quad \Omega \times(0, T) \\
\nabla \cdot \mathbf{u} & =0 & \text { in } \quad \Omega \times(0, T) \\
\mathbf{u}(., t=0) & =\mathbf{u}_{0} & \text { in } \quad \Omega  \tag{2.1}\\
\mathbf{u} & =0 & \text { on } \quad \partial \Omega \times(0, T)
\end{array}
$$

This system does not have a unique solution: if $(\mathbf{u}, p)$ is a solution, then $(\mathbf{u}, p+c)$ is also a solution, for any constant $c$. Then, for uniqueness, one needs, for example, the zero-mean condition on the pressure

$$
\begin{equation*}
\int_{\Omega} p=0 . \tag{2.2}
\end{equation*}
$$

Thus, we introduce the notation $L_{0}^{2}(\Omega)=\left\{p \in L^{2}(\Omega), \int_{\Omega} p=0\right\}$.
Next, we shall introduce the following spaces, which are the completions, in $H^{1}(\Omega)$ and in $L^{2}(\Omega)$, respectively, of the set of compactly supported $\mathcal{C}^{\infty}$ functions with vanishing divergence:

$$
\begin{aligned}
V & =\left\{\mathbf{u} \in\left[H_{0}^{1}(\Omega)\right]^{2}, \nabla \cdot \mathbf{u}=0\right\} \\
H & =\left\{\mathbf{u} \in\left[L^{2}(\Omega)\right]^{2}, \nabla \cdot \mathbf{u}=0, \mathbf{u} \cdot \mathbf{n}_{\partial \Omega}=0 \text { on } \partial \Omega\right\}
\end{aligned}
$$

where $\mathbf{n}_{\partial \Omega}$ is the unit, outward pointing, normal vector field on $\partial \Omega$. We denote by $V^{\prime}$ the dual space of $V$. We recall ([7, Proposition IV.5.13]) that, if $\Omega$, $\mathbf{f}$ and $\mathbf{u}_{0}$ regular enough, problem (2.1)-(2.2) has a unique solution such that

$$
\begin{gathered}
\mathbf{u} \in\left(L^{2}((0, T), V) \cap \mathcal{C}^{0}([0, T], H)\right), \partial_{t} \mathbf{u} \in L^{2}\left((0, T), V^{\prime}\right) \\
p \in W^{-1, \infty}\left((0, T), L_{0}^{2}(\Omega)\right)
\end{gathered}
$$

In order to apply a domain-decomposition strategy for this problem, we decompose $\Omega$ into $M$ non-overlapping subdomains $\Omega_{i}$, i.e. $\Omega_{i} \cap \Omega_{j}=\emptyset$ for $i \neq j$, and $\bar{\Omega}=\bigcup_{i=1}^{M} \overline{\Omega_{i}}$. For $i=1,2, \ldots, M$, we denote by $\mathcal{I}_{i}$ the set of indices of the neighbouring subdomain(s) of $\Omega_{i}$ : it holds that $j \in \mathcal{I}_{i}$ if and only if $\left|\partial \Omega_{i} \cap \partial \Omega_{j}\right|>0$, where $|\cdot|$ denotes the one dimensional measure. We denote by $\Gamma_{i j}$ the interface (if it exists) between $\Omega_{i}$ and $\Omega_{j}, \mathbf{n}_{i j}$ the unit normal vector on $\Gamma_{i j}$, directed from $\Omega_{i}$ to $\Omega_{j}$. Note that this implies that $\mathbf{n}_{i j}=-\mathbf{n}_{j i}$.

Denoting by $\mathbf{u}_{i},\left(\mathbf{u}_{0}\right)_{i}, p_{i}$ and $\mathbf{f}_{i}$ the respective restrictions of $\mathbf{u}, \mathbf{u}_{0}, p$ and $\mathbf{f}$ to $\Omega_{i}$, the monodomain problem is equivalent to the following multidomain one

$$
\begin{array}{rlll}
\partial_{t} \mathbf{u}_{i}-\nu \Delta \mathbf{u}_{i}+\nabla p_{i} & =\mathbf{f}_{i} & & \text { in } \\
\nabla \cdot \Omega_{i} \times(0, T), \\
\mathbf{u}_{i} & =0 & & \text { in } \tag{2.3}
\end{array} \Omega_{i} \times(0, T),
$$

for all $i \in \llbracket 1, M \rrbracket$, together with the physical transmission conditions on the space-time interfaces $\Gamma_{i j} \times(0, T), j \in \mathcal{I}_{i}, i \in \llbracket 1, M \rrbracket$,

$$
\begin{align*}
\mathbf{u}_{i j} \cdot \mathbf{n}_{i j} & =-\mathbf{u}_{j i} \cdot \mathbf{n}_{j i}, \\
\mathbf{u}_{j} \times \mathbf{n}_{i j} & =-\mathbf{u}_{j} \times \mathbf{n}_{j i}, \\
\nu \partial_{\mathbf{n}_{i j}} \mathbf{u}_{i} \cdot \mathbf{n}_{i j}-p_{i} & =\nu \partial_{\mathbf{n}_{j i}} \mathbf{u}_{j} \cdot \mathbf{n}_{j i}-p_{j},  \tag{2.4}\\
\nu \partial_{\mathbf{n}_{i j}} \mathbf{u}_{i} \times \mathbf{n}_{i j} & =\nu \partial_{\mathbf{n}_{j i}} \mathbf{u}_{j} \times \mathbf{n}_{j i} .
\end{align*}
$$

For any choice of $\left(\alpha_{i j}, \alpha_{j i}, \beta_{i j}, \beta_{j i}\right) \in\left(\mathbb{R}^{+*}\right)^{4}$, those conditions are equivalent to the following Robin transmission conditions on $\Gamma_{i j} \times(0, T)=\Gamma_{j i} \times(0, T)$ :

$$
\begin{align*}
\alpha_{i j}\left(\nu \partial_{\mathbf{n}_{i j}} \mathbf{u}_{i} \cdot \mathbf{n}_{i j}-p_{i}\right)+\mathbf{u}_{i} \cdot \mathbf{n}_{i j} & =\alpha_{i j}\left(\nu \partial_{\mathbf{n}_{i j}} \mathbf{u}_{j} \cdot \mathbf{n}_{i j}-p_{j}\right)+\mathbf{u}_{j} \cdot \mathbf{n}_{i j}, \\
\alpha_{j i}\left(\nu \partial_{\mathbf{n}_{j i}} \mathbf{u}_{j} \cdot \mathbf{n}_{j i}-p_{j}\right)+\mathbf{u}_{j} \cdot \mathbf{n}_{j i} & =\alpha_{j i}\left(\nu \partial_{\mathbf{n}_{j i}} \mathbf{u}_{i} \cdot \mathbf{n}_{j i}-p_{i}\right)+\mathbf{u}_{i} \cdot \mathbf{n}_{j i}, \\
\beta_{i j} \nu \partial_{\mathbf{n}_{i j}} \mathbf{u}_{i} \times \mathbf{n}_{i j}+\mathbf{u}_{i} \times \mathbf{n}_{i j} & =\beta_{i j} \nu \partial_{\mathbf{n}_{i j}} \mathbf{u}_{j} \times \mathbf{n}_{i j}+\mathbf{u}_{j} \times \mathbf{n}_{i j},  \tag{2.5}\\
\beta_{j i} \nu \partial_{\mathbf{n}_{j i}} \mathbf{u}_{j} \times \mathbf{n}_{j i}+\mathbf{u}_{j} \times \mathbf{n}_{j i} & =\beta_{j i} \nu \partial_{\mathbf{n}_{j i}} \mathbf{u}_{i} \times \mathbf{n}_{j i}+\mathbf{u}_{i} \times \mathbf{n}_{j i} .
\end{align*}
$$

Finally, the zero-mean condition for the pressure is equivalent to

$$
\begin{equation*}
\sum_{i=1}^{M} \int_{\Omega_{i}} p_{i}=0 \tag{2.6}
\end{equation*}
$$

This setting requires that we should study the Stokes system in a domain where Robin boundary conditions are applied on a part of the boundary. This is what is done in the next section.
3. The Stokes problem with Robin boundary conditions. We now consider a domain, still denoted by $\Omega$, for which the boundary is decomposed into two parts: $\partial \Omega=\Gamma_{D} \cup \Gamma_{R}$, with $\left|\Gamma_{R}\right|>0$. Let $\mathbf{n}$ be the outgoing normal vector on $\Gamma_{R}$; we consider the following system, with $\alpha>0$ and $\beta>0$

$$
\begin{array}{rlll}
\partial_{t} \mathbf{u}-\nu \Delta \mathbf{u}+\nabla p & =\mathbf{f} & & \text { in } \\
\nabla \cdot \mathbf{u} & =0 & & \text { in } \\
\mathbf{u}(\cdot, t=0) & =\mathbf{u}_{0} & \text { in } & \Omega, T),  \tag{3.1}\\
\mathbf{u} & =0 & & \text { on } \\
& \Gamma_{D} \times(0, T), \\
\alpha\left(\nu \partial_{\mathbf{n}} \mathbf{u} \cdot \mathbf{n}-p\right)+\mathbf{u} \cdot \mathbf{n} & =g & & \text { on } \\
\Gamma_{R} \times(0, T), \\
\beta \nu \partial_{\mathbf{n}} \mathbf{u} \times \mathbf{n}+\mathbf{u} \times \mathbf{n} & =\xi & & \text { on } \\
& \Gamma_{R} \times(0, T),
\end{array}
$$

where $\mathbf{f}$ is at least in $\left[L^{2}(\Omega \times(0, T))\right]^{2}, g$ and $\xi$ are at least in $\left[L^{2}\left(\Gamma_{R} \times(0, T)\right)\right]$.
In order to set this problem under an appropriate (parabolic) variational form, we multiply the first equation by a divergence-free test function $\mathbf{v}$ (smooth enough) that vanishes on $\Gamma_{D}$ and integrate by parts on $\Omega$. The flux $\left(-\nu \partial_{\mathbf{n}} \mathbf{u}+p \mathbf{n}\right)$ is then decomposed into normal and tangential parts and boundary conditions of (3.1) are used. We obtain then the following parabolic variational problem

$$
\begin{align*}
& \left\langle\partial_{t} \mathbf{u}, \mathbf{v}\right\rangle_{V_{D}^{\prime}, V_{D}}+a(t, \mathbf{u}, \mathbf{v})=c(t, \mathbf{v}), \quad \text { a.e. } t \in(0, T), \forall \mathbf{v} \in V_{D}  \tag{3.2}\\
& \mathbf{u}(0)=\mathbf{u}_{0} \tag{3.3}
\end{align*}
$$

where the spaces are defined as

$$
\begin{aligned}
V_{D} & =\left\{\mathbf{u} \in\left[H^{1}(\Omega)\right]^{2}, \mathbf{u}=0 \text { on } \Gamma_{D}, \nabla \cdot \mathbf{u}=0\right\} \\
H_{D} & =\left\{\mathbf{u} \in\left[L^{2}(\Omega)\right]^{2}, \mathbf{u} \cdot \mathbf{n}=0 \text { on } \Gamma_{D}, \nabla \cdot \mathbf{u}=0\right\}
\end{aligned}
$$

together with their linear and bilinear forms

$$
\begin{align*}
a(\mathbf{u}, \mathbf{v}) & =\nu(\nabla \mathbf{u}, \nabla \mathbf{v})_{\Omega}+\frac{1}{\alpha}(\mathbf{u} \cdot \mathbf{n}, \mathbf{v} \cdot \mathbf{n})_{\Gamma_{R}}+\frac{1}{\beta}(\mathbf{u} \times \mathbf{n}, \mathbf{v} \times \mathbf{n})_{\Gamma_{R}}  \tag{3.4}\\
c(t, \mathbf{v}) & =(\mathbf{f}(t), \mathbf{v})_{\Omega}+\frac{1}{\alpha}(g(t), \mathbf{v} \cdot \mathbf{n})_{\Gamma_{R}}+\frac{1}{\beta}(\xi(t), \mathbf{v} \times \mathbf{n})_{\Gamma_{R}} \tag{3.5}
\end{align*}
$$

Here, $(\cdot, \cdot)_{D}$ denotes, for any set $D$ (whatever the space-dimension of $D$ ) the standard scalar or the matrix-valued scalar $L^{2}$ product on $D$. In the same way, we shall use the notation $\|\cdot\|_{D}$ for the associated $L^{2}(D)$ norm. All terms in the definition of the forms $a$ and $c$ are well-defined for $(\mathbf{u}, \mathbf{v}) \in V_{D} \times V_{D}$.

From these definitions, $V_{D}$ is dense in $H_{D}$ and the embedding $V_{D} \subset H_{D}$ is continuous. We can identify $H_{D}$ with its dual space, and we are in the situation where $V_{D} \subset H_{D} \equiv H_{D}^{\prime} \subset V_{D}^{\prime}$, which is the classical setting for parabolic equations (see e.g. [17, Section 6.1], [8, Page 218]). In this context, we recall the following theorem.

Theorem 3.1. Problem (3.2)-(3.3) admits a unique solution

$$
\mathbf{u} \in\left(L^{2}\left((0, T), V_{D}\right) \cap \mathcal{C}^{0}\left([0, T], H_{D}\right)\right),
$$

with $\partial_{t} \mathbf{u} \in L^{2}\left((0, T), V_{D}^{\prime}\right)$ if the following properties are verified

- $\mathbf{u}_{0} \in H_{D}$ and $c \in L^{2}\left((0, T), V_{D}^{\prime}\right)$,
- The function $t \mapsto a(t, \mathbf{u}, \mathbf{v})$ is measurable for all $(\mathbf{u}, \mathbf{v}) \in V_{D}^{2}$,
- $\exists M \in \mathbb{R}$ such that $|a(t, \mathbf{u}, \mathbf{v})| \leq M\|\mathbf{u}\|_{V_{D}}\|\mathbf{v}\|_{V_{D}}$ for almost every $t$ and for all $(\mathbf{u}, \mathbf{v}) \in V_{D}^{2}$,
- $\exists m>0$ such that $a(t, \mathbf{u}, \mathbf{u}) \geq m\|\mathbf{u}\|_{V_{D}}^{2}$ for almost every $t$ and for all $\mathbf{u} \in V_{D}$.

We shall apply this result to our setting, with the simplification that the bilinear form defined by (3.4) does not depend on time. We obtain the following result:

Theorem 3.2. Assume that $\mathbf{f} \in L^{2}\left((0, T),\left[L^{2}(\Omega)\right]^{2}\right), g, \xi \in L^{2}\left((0, T), L^{2}\left(\Gamma_{R}\right)\right)$, and $\mathbf{u}_{0} \in H_{D}$. Let a and $c$ be defined by (3.4) and (3.5), respectively. Then, problem (3.2)-(3.3) admits a unique solution $\mathbf{u} \in\left(L^{2}\left((0, T), V_{D}\right) \cap \mathcal{C}^{0}\left([0, T], H_{D}\right)\right)$, which is such that $\partial_{t} \mathbf{u} \in L^{2}\left((0, T), V_{D}^{\prime}\right)$.

Proof. We shall show that $a$ and $c$ verify the hypothesis of Theorem 3.1. First, it is well-known that, as soon as $\left|\Gamma_{R}\right|>0$, then

$$
\|\mathbf{u}\|_{V_{D}}:=\left(\|\nabla \mathbf{u}\|_{\Omega}^{2}+\|\mathbf{u}\|_{\Gamma_{R}}^{2}\right)^{\frac{1}{2}}=\left(\|\nabla \mathbf{u}\|_{\Omega}^{2}+\|\mathbf{u} \cdot \mathbf{n}\|_{\Gamma_{R}}^{2}+\|\mathbf{u} \times \mathbf{n}\|_{\Gamma_{R}}^{2}\right)^{\frac{1}{2}}
$$

is a norm equivalent to the $H^{1}$ norm on $V_{D}$, and we shall therefore work with this norm.

Let $M=\max \left(\nu, \frac{1}{\alpha}, \frac{1}{\beta}\right)$. From the Cauchy-Schwarz inequality, we get the continuity of $a(\cdot, \cdot)$ :

$$
|a(\mathbf{u}, \mathbf{v})| \leq M\|\mathbf{u}\|_{V_{D}}\|\mathbf{v}\|_{V_{D}}, \quad \forall \mathbf{u}, \mathbf{v} \in V_{D}
$$

Let $m=\min \left(\nu, \frac{1}{\alpha}, \frac{1}{\beta}\right)>0$. From the definition of $\|\cdot\|_{V_{D}}$, we get the coercivity of $a(\cdot, \cdot)$ :

$$
a(\mathbf{u}, \mathbf{u}) \geq m\|\mathbf{u}\|_{V_{D}}^{2}, \quad \forall \mathbf{u} \in V_{D}
$$

Then, for a.e. $t \in(0, T)$, the continuity of $c(t, \cdot)$ is deduced from the Cauchy-Schwarz inequality and the equivalence between the $H^{1}(\Omega)$-norm and $\|\cdot\|_{V_{D}}$ :

$$
|c(t, \mathbf{v})| \leq\left[C_{1}\|\mathbf{f}(t)\|_{\Omega}+\frac{1}{\alpha}\|g(t)\|_{\Gamma_{R}}+\frac{1}{\beta}\|\xi(t)\|_{\Gamma_{R}}\right]\|\mathbf{v}\|_{V_{D}}
$$

Moreover, thanks to the hypothesis on the time dependence of $\mathbf{f}, g$ and $\xi$, the quantity

$$
C_{1}\|\mathbf{f}(t)\|_{\Omega}+\frac{1}{\alpha}\|g(t)\|_{\Gamma_{R}}+\frac{1}{\beta}\|\xi(t)\|_{\Gamma_{R}}
$$

is square integrable on $(0, T)$, and we can now apply Theorem 3.1, which finishes the proof.

Remark 3.3. Since $V_{D}$ is continuously and densely embedded in $H_{D}$, the fact that $\mathbf{u} \in \mathcal{C}^{0}\left([0, T], H_{D}\right)$ is a consequence of the fact that the space

$$
\mathcal{W}\left(V_{D}, V_{D}^{\prime}\right):=\left\{\mathbf{v}:(0, T) \mapsto V_{D} ; \mathbf{v} \in L^{2}\left((0, T), V_{D}\right) ; \partial_{t} \mathbf{v} \in L^{2}\left((0, T), V_{D}^{\prime}\right)\right\}
$$

is included in $\mathcal{C}^{0}\left([0, T], H_{D}\right)$, as stated, for example, by [17, Lemma 6.2] and [7, Theorem II.5.13].

This has the important implication that it is legitimate to consider $\mathbf{u}(t) \in H_{D}$ for all $t \in[0, T]$. Moreover, the following integral equality holds for all $t \in[0, T]$ and for all $(\mathbf{u}, \mathbf{v}) \in\left[\mathcal{W}\left(V_{D}, V_{D}^{\prime}\right)\right]^{2}($ see $[17$, Lemma 6.3] and [7, Theorem II.5.12]):

$$
\begin{equation*}
\int_{0}^{t}\left(\left\langle\partial_{t} \mathbf{u}(s), \mathbf{v}(s)\right\rangle_{V_{D}^{\prime}, V_{D}}+\left\langle\partial_{t} \mathbf{v}(s), \mathbf{u}(s)\right\rangle_{V_{D}^{\prime}, V_{D}}\right) d s=(\mathbf{u}(t), \mathbf{v}(t))_{\Omega}-(\mathbf{u}(0), \mathbf{v}(0))_{\Omega} \tag{3.6}
\end{equation*}
$$

Now, since we have obtained the velocity $\mathbf{u}$ from the constrained variational problem (3.2)-(3.3), we shall construct the pressure by relaxing the divergence free condition on the velocity test functions, and we shall therefore consider the space

$$
X_{D}=\left\{\mathbf{v} \in\left[H^{1}(\Omega)\right]^{2}, \mathbf{v}=0 \text { on } \Gamma_{D}\right\}
$$

equipped with the above-defined norm $\|\cdot\|_{V_{D}}$. Like often with the Stokes problem, we shall rely on the surjectivity of the divergence operator, and on general properties of surjective mappings in Hilbert spaces. More precisely, we shall use the following results.

Lemma 3.4. The mapping $B$ from $X_{D}$ into $L^{2}(\Omega)$ defined by $B(\mathbf{v})=-\nabla \cdot \mathbf{v}$ is continuous and surjective.

Proof. This is a special case of [17, Lemma 4.9] (with, using the notations of [17], $\partial \Omega_{1}=\Gamma_{D}, \partial \Omega_{2}=\emptyset, \partial \Omega_{3}=\emptyset$ and $\left.\partial \Omega_{4}=\Gamma_{R}\right)$.

Lemma 3.5. Let $L$ be in $\mathcal{L}(E ; F)$ and $L^{T}$ be its adjoint in $\mathcal{L}\left(F^{\prime} ; E^{\prime}\right)$, then if $L$ is surjective in $F$, then $\operatorname{Im} L^{T}$ is closed in $E^{\prime}$.
Before stating the next Lemma, we recall the following definition (see, e.g. [7, Definition IV.2.1]) and properties (see, e.g. [7, Remark IV.2.1])

Definition 3.6. Let $E$ be a Banach space with dual space $E^{\prime}$; then for any subset $A \subset E$, we define $A^{\perp} \subset E^{\prime}$ as follows:

$$
A^{\perp}:=\left\{\phi \in E^{\prime}, \forall x \in A,\langle\phi, x\rangle_{E^{\prime}, E}=0\right\}
$$

Lemma 3.7. If $A \subset C \subset E$, then $C^{\perp} \subset A^{\perp}$.
Lemma 3.8. If $A$ is a linear subspace of $E$, then $\left(A^{\perp}\right)^{\perp}=A$ if and only if $A$ is closed in $E$.
Moreover, we also recall the following general result
Lemma 3.9. Let $L$ be in $\mathcal{L}(E ; F)$, then $\left(\operatorname{Im} L^{T}\right)^{\perp} \subset \operatorname{Ker} L$
Proof. If $f \in\left(\operatorname{Im} L^{T}\right)^{\perp}$, then $\left\langle L^{T} q, f\right\rangle_{E^{\prime}, E}=0, \forall q \in F^{\prime}$. Thus $\langle q, L f\rangle_{F^{\prime}, F}=0$ for all $q \in F^{\prime}$, which means that $L f=0$, and thus $f \in \operatorname{Ker} L$.
From these results, we obtain the following Lemma, which will be useful in the construction of the pressure field:

Lemma 3.10. Let $B^{T}$ be the adjoint operator of $B$, from $L^{2}(\Omega)$ into $X_{D}^{\prime}$. Then for any $\ell$ in $X_{D}^{\prime}$ that vanishes on $V_{D}$, there exists $P \in L^{2}(\Omega)$ such that $\ell=B^{T} P$.

Proof. Since $B$ is in $\mathcal{L}\left(X_{D} ; L^{2}(\Omega)\right)$ and is surjective (Lemma 3.4), then $\left(\operatorname{Im} B^{T}\right)$ is closed in $X_{D}^{\prime}$ (Lemma 3.5), and $\left(\left(\operatorname{Im} B^{T}\right)^{\perp}\right)^{\perp}=\operatorname{Im} B^{T}$ (Lemma 3.8). Now, using

Lemmas 3.9 and 3.7, we get $(\operatorname{Ker} B)^{\perp} \subset\left(\left(\operatorname{Im} B^{T}\right)^{\perp}\right)^{\perp}=\operatorname{Im} B^{T}$. So if $\ell$ in $X_{D}^{\prime}$ vanishes on $V_{D}=\operatorname{Ker} B$, then $\ell$ is in $(\operatorname{Ker} B)^{\perp}$ and so in $\operatorname{Im} B^{T}$, which exactly means that there exists $P \in L^{2}(\Omega)$ such that $\ell=B^{T} P$.

Using this result, we can now state the following theorem.
Theorem 3.11. Assume that $\mathbf{f} \in L^{2}\left((0, T),\left[L^{2}(\Omega)\right]^{2}\right), \xi, g \in L^{2}\left((0, T), L^{2}\left(\Gamma_{R}\right)\right)$ and $\mathbf{u}_{0} \in H_{D}$, then there exists unique $\mathbf{u} \in\left(L^{2}\left((0, T), V_{D}\right) \cap \mathcal{C}^{0}\left([0, T], H_{D}\right)\right)$ and $p \in W^{-1, \infty}\left((0, T), L^{2}(\Omega)\right)$, with $\partial_{t} \mathbf{u} \in L^{2}\left((0, T), V_{D}^{\prime}\right)$ such that $(\mathbf{u}, p)$ verifies problem (3.1) in the sense that

- $\mathbf{u}$ verifies (3.2)-(3.3)
- $p=\partial_{t} P$ with $P \in L^{\infty}\left((0, T), L^{2}(\Omega)\right)$ that satisfies
$\int_{0}^{t} c(s, \mathbf{v}) d s-(\mathbf{u}(t), \mathbf{v})_{\Omega}+\left(\mathbf{u}_{0}, \mathbf{v}\right)_{\Omega}-\int_{0}^{t} a(\mathbf{u}(s), \mathbf{v}) d s=-\int_{\Omega} P(t) \nabla \cdot \mathbf{v}, \quad \forall \mathbf{v} \in X_{D}$.

Proof. Let $\mathbf{u}$ be the solution of (3.2)-(3.3), and consider, for this $\mathbf{u}$, the function $t \mapsto a(\mathbf{u}(t), \mathbf{v})$ and the function $t \mapsto c(t, \mathbf{v})$ where $a$ and $c$ are defined by (3.4) and (3.5). Then their definitions can be straightforwardly extended to consider $\mathbf{v} \in X_{D}$ and, for any $t \in(0, T)$, the following element of $X_{D}^{\prime}$ is well-defined:

$$
b(t, \mathbf{v}):=\int_{0}^{t} c(s, \mathbf{v}) d s-(\mathbf{u}(t), \mathbf{v})_{\Omega}+\left(\mathbf{u}_{0}, \mathbf{v}\right)_{\Omega}-\int_{0}^{t} a(\mathbf{u}(s), \mathbf{v}) d s, \quad \forall \mathbf{v} \in X_{D}
$$

Indeed, one has that

$$
\begin{aligned}
\left|\int_{0}^{t} a(\mathbf{u}(s), \mathbf{v}) d s\right| & \leq \int_{0}^{t} M\|\mathbf{u}(s)\|_{X_{D}}\|\mathbf{v}\|_{X_{D}} d s \\
& \leq M \sqrt{t}\left[\int_{0}^{t}\|\mathbf{u}(s)\|_{X_{D}}^{2} d s\right]^{\frac{1}{2}}\|\mathbf{v}\|_{X_{D}} \\
& \leq M \sqrt{T}\|\mathbf{u}\|_{L^{2}\left((0, T), V_{D}\right)}\|\mathbf{v}\|_{X_{D}}
\end{aligned}
$$

and

$$
\begin{aligned}
\left|\int_{0}^{t} c(s, \mathbf{v}) d s\right| & \leq \int_{0}^{t}\left(C_{1}\|\mathbf{f}(s)\|_{\Omega}+\frac{1}{\alpha}\|g(s)\|_{\Gamma_{R}}+\frac{1}{\beta}\|\xi(s)\|_{\Gamma_{R}}\right)\|\mathbf{v}\|_{X_{D}} \\
& \leq \gamma_{1}\|\mathbf{v}\|_{X_{D}}
\end{aligned}
$$

with

$$
\gamma_{1}=C_{1} \sqrt{T}\|\mathbf{f}\|_{L^{2}\left((0, T),\left[L^{2}(\Omega)\right]^{2}\right)}+\frac{\sqrt{T}}{\alpha}\|g\|_{L^{2}\left((0, T), L^{2}\left(\Gamma_{R}\right)\right)}+\frac{\sqrt{T}}{\beta}\|\xi\|_{L^{2}\left((0, T), L^{2}\left(\Gamma_{R}\right)\right)} .
$$

In addition, since $\mathbf{u}$ belongs to $\mathcal{C}^{0}\left([0, T], H_{D}\right)$, then

$$
\begin{aligned}
\left|-(\mathbf{u}(t), \mathbf{v})_{\Omega}+\left(\mathbf{u}_{0}, \mathbf{v}\right)_{\Omega}\right| & \leq 2\|\mathbf{u}\|_{L^{\infty}\left([0, T],\left[L^{2}(\Omega)\right]^{2}\right)}\|\mathbf{v}\|_{\Omega} \\
& \leq 2 C_{1}\|\mathbf{u}\|_{L^{\infty}\left([0, T],\left[L^{2}(\Omega)\right]^{2}\right)}\|\mathbf{v}\|_{X_{D}}
\end{aligned}
$$

This leads to the fact that

$$
\begin{equation*}
|b(t, \mathbf{v})| \leq C_{2}\|\mathbf{v}\|_{X_{D}}, \quad \forall \mathbf{v} \in V_{D}, \quad \forall t \in(0, T) \tag{3.8}
\end{equation*}
$$

with

$$
C_{2}=2 C_{1}\|\mathbf{u}\|_{L^{\infty}\left([0, T],\left[L^{2}(\Omega)\right]^{2}\right)}+\gamma_{1}+M \sqrt{T}\|\mathbf{u}\|_{L^{2}\left(0, T, V_{D}\right)}
$$

Moreover, from (3.2) and (3.6) (with $\mathbf{v}$ not depending on time), we obtain that $b(t, \mathbf{v})=0$ for all $\mathbf{v} \in V_{D}$, for all $t \in(0, T)$. Thus, using Lemma 3.10, we conclude that, for all $t \in(0, T)$, there exists $P(t) \in L^{2}(\Omega)$ satisfying

$$
\begin{equation*}
b(t, \mathbf{v})=\left\langle B^{T} P(t), \mathbf{v}\right\rangle_{X_{D}^{\prime}, X_{D}}=-(P(t), \nabla \cdot \mathbf{v})_{\Omega}=-\int_{\Omega} P(t) \nabla \cdot \mathbf{v}, \quad \forall \mathbf{v} \in X_{D} \tag{3.9}
\end{equation*}
$$

Moreover, the surjectivity of the divergence mapping leads to the following inf-sup condition: there exists $\gamma_{2}>0$, s.t.

$$
\inf _{q \in L^{2}(\Omega)} \sup _{\mathbf{v} \in X_{D}} \frac{(B \mathbf{v}, q)_{\Omega}}{\|\mathbf{v}\|_{X_{D}}\|q\|_{L^{2}(\Omega)}}=\gamma_{2}>0
$$

which implies, for all $q \in L^{2}(\Omega)$

$$
\begin{equation*}
\gamma_{2}\|q\|_{L^{2}(\Omega)} \leq \sup _{\mathbf{v} \in X_{D}} \frac{(B \mathbf{v}, q)_{\Omega}}{\|\mathbf{v}\|_{X_{D}}} \tag{3.10}
\end{equation*}
$$

In order to use $q=P(t)$ in (3.10), we need to evaluate $(B \mathbf{v}, P(t))_{\Omega}$. From (3.9), we obtain that $(B \mathbf{v}, P(t))_{\Omega}=\left\langle B^{T} P(t), \mathbf{v}\right\rangle_{X_{D}^{\prime}, X_{D}}=b(t, \mathbf{v})$; together with (3.8), we get

$$
\|P(t)\|_{L^{2}(\Omega)} \leq \frac{1}{\gamma_{2}} \sup _{\mathbf{v} \in X_{D}} \frac{b(t, \mathbf{v})}{\|\mathbf{v}\|_{X_{D}}} \leq \frac{C_{2}}{\gamma_{2}}
$$

We conclude that $P(t) \in L^{\infty}\left((0, T), L^{2}(\Omega)\right)$. Then, we define the pressure $p=\partial_{t} P$ and thus $p \in H^{-1, \infty}\left((0, T), L^{2}(\Omega)\right)$.

It remains to show that $p$ is unique. Consider the case $\mathbf{u}_{0}=0$ and $c=0$. Then, we have $\mathbf{u}=0$, and (3.7) leads to $\int_{\Omega} P(t) \nabla \cdot \mathbf{v}=0, \forall \mathbf{v} \in X_{D}$. From the surjectivity of the divergence mapping, one gets that $P(t)=0$ for all $t$, and then $p=0$.
4. Optimized Schwarz Waveform Relaxation Algorithm. The OSWR algorithm for solving the multidomain problem (2.3)-(2.4) is as follows.

```
Algorithm 4.1 (OSWR)
    Choose initial Robin data \(g_{i j}^{0}, \xi_{i j}^{0}\) on \(\Gamma_{i j} \times(0, T), j \in \mathcal{I}_{i}, i=1,2, \ldots, M\)
    for \(\ell=1,2, \ldots\) do
        1. Solve the local space-time Robin problems, for \(i=1,2, \ldots, M\)
```

```
\[
\begin{align*}
& \begin{array}{rlrl}
\partial_{t} \mathbf{u}_{i}^{\ell}-\nu \Delta \mathbf{u}_{i}^{\ell}+\nabla p_{i}^{\ell} & =\mathbf{f}_{i} & & \text { in } \\
\nabla \cdot \mathbf{u}_{i}^{\ell} & =0 & & \text { in } \\
\Omega_{i} \times(0, T) \\
\Omega_{i} \times(0, T)
\end{array} \\
& \mathbf{u}_{i}^{\ell}(., t=0)=\mathbf{u}_{0, i} \quad \text { in } \quad \Omega_{i} \\
& \alpha_{i j}\left(\nu \partial_{\mathbf{n}_{i j}} \mathbf{u}_{i}^{\ell} \cdot \mathbf{n}_{i j}-p_{i}^{\ell}\right)+\mathbf{u}_{i}^{\ell} \cdot \mathbf{n}_{i j}=g_{i j}^{\ell-1} \quad \text { on } \quad \Gamma_{i j} \times(0, T), \quad j \in \mathcal{I}_{i}  \tag{4.1}\\
& \begin{array}{rlrlrl}
\beta_{i j} \nu \partial_{\mathbf{n}_{i j}} \mathbf{u}_{i}^{\ell} \times \mathbf{n}_{i j}+\mathbf{u}_{i}^{\ell} \times \mathbf{n}_{i j} & =\xi_{i j}^{\ell-1} & \text { on } \quad & \Gamma_{i j} \times(0, T), j \in \mathcal{I}_{i} \\
\mathbf{u}_{i}^{\ell} & =0 & \text { on } & & \left(\partial \Omega_{i} \cap \partial \Omega\right) \times(0, T)
\end{array}
\end{align*}
\]
```

2. Update the Robin terms $g_{i j}^{\ell}, \xi_{i j}^{\ell}$ on $\Gamma_{i j} \times(0, T)$, for $j \in \mathcal{I}_{i}, i=1,2, \ldots, M$

$$
\begin{align*}
g_{i j}^{\ell} & =\alpha_{i j}\left(\nu \partial_{\mathbf{n}_{i j}} \mathbf{u}_{j}^{\ell} \cdot \mathbf{n}_{i j}-p_{j}^{\ell}\right)+\mathbf{u}_{j}^{\ell} \cdot \mathbf{n}_{i j},  \tag{4.2a}\\
\xi_{i j}^{\ell} & =\beta_{i j} \nu \partial_{\mathbf{n}_{i j}} \mathbf{u}_{j}^{\ell} \times \mathbf{n}_{i j}+\mathbf{u}_{j}^{\ell} \times \mathbf{n}_{i j} . \tag{4.2b}
\end{align*}
$$

end for

Remark 4.1. Let $i \in \llbracket 1, M \rrbracket, j \in \mathcal{I}_{i}$. Formulas given by (4.2) can be rewritten as

$$
\begin{aligned}
g_{i j}^{\ell} & =\frac{\alpha_{i j}}{\alpha_{j i}}\left(\alpha_{j i}\left(\nu \partial_{\mathbf{n}_{j i}} \mathbf{u}_{j}^{\ell} \cdot \mathbf{n}_{j i}-p_{j}^{\ell}\right)+\mathbf{u}_{j}^{\ell} \cdot \mathbf{n}_{j i}\right)-\frac{\alpha_{i j}}{\alpha_{j i}} \mathbf{u}_{j}^{\ell} \cdot \mathbf{n}_{j i}+\mathbf{u}_{j}^{\ell} \cdot \mathbf{n}_{i j} \\
\xi_{i j}^{\ell} & =\frac{\beta_{i j}}{\beta_{j i}}\left(\beta_{j i} \nu \partial_{\mathbf{n}_{j i}} \mathbf{u}_{j}^{\ell} \times \mathbf{n}_{j i}+\mathbf{u}_{j}^{\ell} \times \mathbf{n}_{j i}\right)-\frac{\beta_{i j}}{\beta_{j i}} \mathbf{u}_{j}^{\ell} \times \mathbf{n}_{j i}+\mathbf{u}_{j}^{\ell} \times \mathbf{n}_{i j}
\end{aligned}
$$

or equivalently, using the Robin transmission conditions in (4.1),

$$
\begin{align*}
g_{i j}^{\ell} & =\frac{\alpha_{i j}}{\alpha_{j i}} g_{j i}^{\ell-1}-\frac{\alpha_{i j}+\alpha_{j i}}{\alpha_{j i}} \mathbf{u}_{j}^{\ell} \cdot \mathbf{n}_{j i},  \tag{4.3a}\\
\xi_{i j}^{\ell} & =\frac{\beta_{i j}}{\beta_{j i}} \xi_{j i}^{\ell-1}-\frac{\beta_{i j}+\beta_{j i}}{\beta_{j i}} \mathbf{u}_{j}^{\ell} \times \mathbf{n}_{j i} . \tag{4.3b}
\end{align*}
$$

One advantage of formula (4.3) is that, if $g_{i j}^{\ell-1}$ and $\xi_{i j}^{\ell-1}$ have $L^{2}\left(\Gamma_{i j}\right)$ regularity, so will $g_{i j}^{\ell}$ and $\xi_{i j}^{\ell}$. Indeed, in (4.3) the regularities of $g_{i j}^{\ell}$ and $\xi_{i j}^{\ell}$ depend only on those of $g_{j i}^{\ell-1}, \xi_{j i}^{\ell-1}$ and $\mathbf{u}_{j}^{\ell}$, whose trace is in $L^{2}\left((0, T), H^{\frac{1}{2}}\left(\Gamma_{i j}\right)\right)$ (recall that we have $\mathbf{u}_{j}^{\ell} \in L^{2}\left((0, T),\left[H^{1}\left(\Omega_{j}\right)\right]^{2}\right)$, see Section 3). On the other hand, formula (4.2) will return new Robin boundary data $g_{i j}^{\ell}$ and $\xi_{i j}^{\ell}$ with a lower regularity, which is not satisfying for an iterative algorithm. Another advantage of formula (4.3) is that it is easier to implement in practice, than formula (4.2).
Now, we may express the iterative algorithm in the following way. We first define

$$
\begin{gathered}
V_{i}=\left\{\mathbf{u} \in\left[H^{1}\left(\Omega_{i}\right)\right]^{2}, \mathbf{u}=0 \text { on } \partial \Omega_{i} \cap \partial \Omega, \nabla \cdot \mathbf{u}=0 \text { in } \Omega_{i}\right\} \\
H_{i}=\left\{\mathbf{u} \in\left[L^{2}\left(\Omega_{i}\right)\right]^{2}, \mathbf{u} \cdot \mathbf{n}_{\partial \Omega_{i}}=0 \text { on } \partial \Omega_{i} \cap \partial \Omega, \nabla \cdot \mathbf{u}=0 \text { in } \Omega_{i}\right\} . \\
X_{i}=\left\{\mathbf{u} \in\left[H^{1}\left(\Omega_{i}\right)\right]^{2}, \mathbf{u}=0 \text { on } \partial \Omega_{i} \cap \partial \Omega\right\}
\end{gathered}
$$

Then, we set, for all $\mathbf{u}, \mathbf{v} \in X_{i}$ and $t \in(0, T)$,

$$
\begin{align*}
a_{i}(\mathbf{u}, \mathbf{v}) & :=\nu(\nabla \mathbf{u}, \nabla \mathbf{v})_{\Omega_{i}}+\sum_{j \in \mathcal{I}_{i}} \frac{1}{\alpha_{i j}}\left(\mathbf{u} \cdot \mathbf{n}_{i j}, \mathbf{v} \cdot \mathbf{n}_{i j}\right)_{\Gamma_{i j}}+\frac{1}{\beta_{i j}}\left(\mathbf{u} \times \mathbf{n}_{i j}, \mathbf{v} \times \mathbf{n}_{i j}\right)_{\Gamma_{i j}}, \\
c_{i}^{\ell}(t, \mathbf{v}) & :=(\mathbf{f}(t), \mathbf{v})_{\Omega_{i}}+\sum_{j \in \mathcal{I}_{i}} \frac{1}{\alpha_{i j}}\left(g_{i j}^{\ell-1}(t), \mathbf{v} \cdot \mathbf{n}_{i j}\right)_{\Gamma_{i j}}+\frac{1}{\beta_{i j}}\left(\xi_{i j}^{\ell-1}(t), \mathbf{v} \times \mathbf{n}_{i j}\right)_{\Gamma_{i j}}, \tag{4.4}
\end{align*}
$$

and the algorithm reads: for all $\ell \geq 1$, given $g_{i j}^{\ell-1}, \xi_{i j}^{\ell-1}$ on each space-time interface $\Gamma_{i j} \times(0, T)$, solve, for each $i=1 \ldots M$ :

$$
\begin{equation*}
\left\langle\partial_{t} \mathbf{u}_{i}^{\ell}, \mathbf{v}\right\rangle_{V_{i}^{\prime}, V_{i}}+a_{i}\left(\mathbf{u}_{i}^{\ell}, \mathbf{v}\right)=c_{i}^{\ell}(t, \mathbf{v}), \quad \text { a.e. } t \in(0, T), \forall \mathbf{v} \in V_{i} \tag{4.5}
\end{equation*}
$$

$$
\mathbf{u}_{i}^{\ell}(0)=\mathbf{u}_{0 . i}
$$

Then we construct $p_{i}^{\ell}=\partial_{t} P_{i}^{\ell}$, where $P_{i}^{\ell}$ is such that

$$
\begin{array}{r}
\left(\mathbf{u}_{i}^{\ell}(t), \mathbf{v}\right)_{\Omega_{i}}-\left(\mathbf{u}_{0, i}, \mathbf{v}\right)_{\Omega_{i}}+\int_{0}^{t} a_{i}\left(\mathbf{u}_{i}^{\ell}(s), \mathbf{v}\right) d s-\left(P_{i}^{\ell}, \nabla \cdot \mathbf{v}\right)_{\Omega_{i}}-\int_{0}^{t} c_{i}^{\ell}(s, \mathbf{v}) d s=0 \\
\forall \mathbf{v} \in X_{i} \tag{4.6}
\end{array}
$$

Finally, the data are updated by using (4.3a)-(4.3b) on the space-time interfaces.
With this formulation, we can state the following result

THEOREM 4.2. Assume that $g_{i j}^{0}, \xi_{i j}^{0} \in L^{2}\left((0, T), L^{2}\left(\Gamma_{i j}\right)\right)$ and $\left.\mathbf{u}_{0}\right|_{\Omega_{i}} \in H_{i}$. Then, the OSWR algorithm is well-defined and for all $\ell, \mathbf{u}_{i}^{\ell} \in L^{2}\left((0, T), V_{i}\right) \cap \mathcal{C}^{0}\left([0, T], H_{i}\right)$, $\partial_{t} \mathbf{u}_{i}^{\ell} \in L^{2}\left((0, T), V_{i}^{\prime}\right), p_{i}^{\ell} \in W^{-1, \infty}\left((0, T), L^{2}\left(\Omega_{i}\right)\right)$ and $g_{i j}^{\ell}, \xi_{i j}^{\ell} \in L^{2}\left((0, T), L^{2}\left(\Gamma_{i j}\right)\right)$.

Proof. By Theorem 3.11, if $g_{i j}^{\ell-1}, \xi_{i j}^{\ell-1} \in L^{2}\left((0, T), L^{2}\left(\Gamma_{i j}\right)\right)$, then one gets $\mathbf{u}_{i}^{\ell}$ verifying (4.5) with $\mathbf{u}_{i}^{\ell} \in L^{2}\left((0, T), V_{i}\right) \cap \mathcal{C}^{0}\left([0, T], H_{i}\right)$ and $\partial_{t} \mathbf{u}_{i}^{\ell} \in L^{2}\left((0, T), V_{i}^{\prime}\right)$. Additionally, Theorem 3.11 tells us that there exists $P_{i}^{\ell}$ verifying (4.6). We take $p_{i}^{\ell}=\partial_{t} P_{i}^{\ell} \in W^{-1, \infty}\left((0, T), L^{2}\left(\Omega_{i}\right)\right)$.

Using the trace theorem, the normal and tangent traces of $\mathbf{u}_{i}^{\ell}$ on $\Gamma_{i j} \times(0, T)$ belong to $L^{2}\left((0, T), L^{2}\left(\Gamma_{i j}\right)\right)$. Hence, using the update formula (4.3), we infer that $g_{i j}^{\ell}, \xi_{i j}^{\ell} \in L^{2}\left((0, T), L^{2}\left(\Gamma_{i j}\right)\right)$.
The proof is then carried out by a simple induction.
Remark 4.3. The OSWR algorithm is constructed without considering the last condition (2.6), hence it may not converge to the monodomain solution. We shall show in the next section that, indeed, the pressure in each subdomain may not converge to the restriction of the monodomain pressure.
5. First observations on the two subdomains case. For the trivial case of a onedimensional problem and two subdomains, one can show that the velocity iterates converge, while the pressure iterates do not converge in general, see [9].

This result generalizes to higher dimensions as follows : let us consider the twosubdomain case, i.e. $M=2$. To simplify notation, we set $\Gamma:=\Gamma_{12}=\Gamma_{21}$, and for any $\phi$ in $(\alpha, g, \mathbf{u})$, we write $\phi_{1}$ and $\phi_{2}$ instead of $\phi_{12}$ and $\phi_{21}$, respectively.

The divergence-free condition of the velocity in each subdomain leads to

$$
\begin{equation*}
\int_{\partial \Omega_{i}} \mathbf{u}_{i}^{\ell} \cdot \mathbf{n}_{\partial \Omega_{i}}=0=\int_{\Gamma} \mathbf{u}_{i}^{\ell} \cdot \mathbf{n}_{i}, \quad i=1,2 \tag{5.1}
\end{equation*}
$$

The update of Robin terms for the normal components can also be written as

$$
g_{i}^{\ell}=\frac{\alpha_{i}}{\alpha_{j}} g_{j}^{\ell-1}-\frac{\alpha_{i}+\alpha_{j}}{\alpha_{j}} \mathbf{u}_{j}^{\ell} \cdot \mathbf{n}_{j}, \quad j=3-i, i=1,2 .
$$

Integrating over $\Gamma$, and taking (5.1) into account, we get

$$
\int_{\Gamma} g_{i}^{\ell}=\frac{\alpha_{i}}{\alpha_{j}} \int_{\Gamma} g_{j}^{\ell-1}=\int_{\Gamma} g_{i}^{\ell-2}, \quad j=3-i, i=1,2
$$

Therefore, a necessary condition for the convergence of the algorithm to the monodomain solution is

$$
\begin{equation*}
\int_{\Gamma} g_{i}^{0}=\int_{\Gamma} g_{i}, \quad i=1,2 \tag{5.2}
\end{equation*}
$$

with $g_{i}=\alpha_{i}\left(\nu \partial_{\mathbf{n}_{i}} \mathbf{u} \cdot \mathbf{n}_{i}-p\right)+\mathbf{u} \cdot \mathbf{n}_{i}, i=1,2$, in which $(\mathbf{u}, p)$ is the monodomain solution of problem (2.1). Condition (5.2) cannot be achieved in practice because the quantity $g_{i}, i=1,2$, is not known.

More precisely, whereas the convergence of the velocity iterates will be proven in Section 6 below, independently of condition (5.2), the pressure iterates will converge only if condition (5.2) is satisfied, and thus will not converge in general. A correction technique to recover the pressure from the velocity will be proposed in Section 7.
6. Convergence of the velocity via energy estimate. In this Section, we suppose additional regularity on $\mathbf{u}_{0}, \mathbf{f}$ and $\Omega$, which leads to regularity properties of the strong solution of problem (2.1)-(2.2). Namely, we recall [31, Theorem 1, Page 86].

THEOREM 6.1. Let $\Omega$ be a bounded domain of $\mathbb{R}^{2}$ with twice continuously differentiable boundary. For any $\mathbf{u}_{0} \in V$ and $\left.\mathbf{f} \in L^{2}\left((0, T), L^{2}(\Omega)\right)^{2}\right)$, problem (2.1)-(2.2) has a unique solution $(\mathbf{u}, p)$ such that

$$
\begin{gathered}
\mathbf{u} \in \mathcal{C}^{0}([0, T], V) \cap L^{2}\left((0, T),\left(H^{2}(\Omega)\right)^{2}\right), \partial_{t} \mathbf{u} \in L^{2}\left((0, T), L^{2}(\Omega)\right)^{2}, \\
p \in L^{2}\left((0, T), H^{1}(\Omega)\right) .
\end{gathered}
$$

Using Theorem 6.1, we prove that, if its hypotheses are satisfied, then the velocity iterates converge to the monodomain velocity.

Theorem 6.2. Assume that the hypotheses of Theorem 6.1 are satisfied. Let $g_{i j}^{0}$ and $\xi_{i j}^{0}$ belong to $L^{2}\left((0, T), L^{2}\left(\Gamma_{i j}\right)\right)$ and let $\mathbf{u}_{i}^{\ell}$ be the velocity component of the solution of Algorithm 4.1 (OSWR). Then, if $\alpha_{i j}=\alpha_{j i}$ and $\beta_{i j}=\beta_{j i}$, the sequence $\mathbf{u}_{i}^{\ell}$ converges to $\mathbf{u}_{i}=\left.\mathbf{u}\right|_{\Omega_{i}}$ in $\mathcal{C}^{0}\left([0, T], H_{i}\right) \cap L^{2}\left(0, T, V_{i}\right)$.

Proof. Denote by $p_{i}=\left.p\right|_{\Omega_{i}}$. Then, thanks to the extra regularity of $(\mathbf{u}, p)$ given by Theorem 6.1, we can define its Robin trace on any space-time interface $\Gamma_{i j} \times(0, T)$

$$
\begin{align*}
g_{i j}^{\ell} & =\frac{\alpha_{i j}}{\alpha_{j i}} g_{j i}^{\ell-1}-\frac{\alpha_{i j}+\alpha_{j i}}{\alpha_{j i}} \mathbf{u}_{j}^{\ell} \cdot \mathbf{n}_{j i},  \tag{6.1a}\\
\xi_{i j}^{\ell} & =\frac{\beta_{i j}}{\beta_{j i}} \xi_{j i}^{\ell-1}-\frac{\beta_{i j}+\beta_{j i}}{\beta_{j i}} \mathbf{u}_{j}^{\ell} \times \mathbf{n}_{j i} . \tag{6.1b}
\end{align*}
$$

and they both belong to $L^{2}\left((0, T), L^{2}\left(\Gamma_{i j}\right)\right)$. Then (2.5) implies

$$
\begin{align*}
g_{i j}^{\ell} & =\frac{\alpha_{i j}}{\alpha_{j i}} g_{j i}^{\ell-1}-\frac{\alpha_{i j}+\alpha_{j i}}{\alpha_{j i}} \mathbf{u}_{j}^{\ell} \cdot \mathbf{n}_{j i},  \tag{6.2a}\\
\xi_{i j}^{\ell} & =\frac{\beta_{i j}}{\beta_{j i}} \xi_{j i}^{\ell-1}-\frac{\beta_{i j}+\beta_{j i}}{\beta_{j i}} \mathbf{u}_{j}^{\ell} \times \mathbf{n}_{j i} . \tag{6.2b}
\end{align*}
$$

Moreover, $\left(\mathbf{u}_{i}, p_{i}\right)$ is the strong solution of each local Robin boundary problem with source term $\mathbf{f}_{i}$, initial condition $\mathbf{u}_{0, i}$ and Robin terms $g_{i j}$ and $\xi_{i j}$ on $\Gamma_{i j}$. We can write these local problems in variational forms similar to (4.4)-(4.5), in which we replace $g_{i j}^{\ell}$ by $g_{i j}$ and $\xi_{i j}^{\ell}$ by $\xi_{i j}$.

We define the errors as the differences between the iterates and the restrictions (to each subdomain) of the monodomain solution and denote by

$$
\begin{equation*}
\mathbf{e}_{i}^{\ell}:=\mathbf{u}_{i}^{\ell}-\mathbf{u}_{i}, h_{i j}^{\ell}=g_{i j}^{\ell}-g_{i j}, \zeta_{i j}^{\ell}=\xi_{i j}^{\ell}-\xi_{i j}, \quad j \in \mathcal{I}_{i}, i \in \llbracket 1, M \rrbracket . \tag{6.3}
\end{equation*}
$$

Then, the errors also verify the following variational problems similar to (4.4)-(4.5): for a.e. $t \in(0, T), \forall \mathbf{v} \in V_{i}$,

$$
\begin{equation*}
\left\langle\partial_{t} \mathbf{e}_{i}^{\ell}, \mathbf{v}\right\rangle_{V_{i}^{\prime}, V_{i}}+a_{i}\left(\mathbf{e}_{i}^{\ell}, \mathbf{v}\right)=\sum_{j \in \mathcal{I}_{i}} \frac{1}{\alpha_{i j}}\left(h_{i j}^{\ell-1}, \mathbf{v} \cdot \mathbf{n}_{i j}\right)_{\Gamma_{i j}}+\sum_{j \in \mathcal{I}_{i}} \frac{1}{\beta_{i j}}\left(\zeta_{i j}^{\ell-1}, \mathbf{v} \times \mathbf{n}_{i j}\right)_{\Gamma_{i j}}, \tag{6.4}
\end{equation*}
$$

with initial condition $\mathbf{e}_{i}^{\ell}(0)=0$. All integrals on $\Gamma_{i j}$ are well defined since $g_{i j}$ and $\xi_{i j}$ are both in $L^{2}\left((0, T), L^{2}\left(\Gamma_{i j}\right)\right)$, and since we have proved that this is also the case for $g_{i j}^{\ell}$ and $\xi_{i j}^{\ell}$ as soon as it is true for $\ell=0$.

With $\alpha_{i j}=\alpha_{j i}$ and $\beta_{i j}=\beta_{j i}$, the update formulas (4.3) and (6.2) for the Robin terms on $\Gamma_{i j} \times(0, T)$ lead to

$$
\begin{equation*}
\mathbf{e}_{i}^{\ell} \cdot \mathbf{n}_{i j}=\frac{1}{2}\left(h_{i j}^{\ell-1}-h_{j i}^{\ell}\right) \quad, \quad \mathbf{e}_{i}^{\ell} \times \mathbf{n}_{i j}=\frac{1}{2}\left(\zeta_{i j}^{\ell-1}-\zeta_{j i}^{\ell}\right) . \tag{6.5}
\end{equation*}
$$

Choosing $\mathbf{e}_{i}^{\ell}$ as test function in (6.4), one gets

$$
\begin{align*}
& \left\langle\partial_{t} \mathbf{e}_{i}^{\ell}, \mathbf{e}_{i}^{\ell}\right\rangle_{V_{i}^{\prime}, V_{i}}+\nu\left(\nabla \mathbf{e}_{i}^{\ell}, \nabla \mathbf{e}_{i}^{\ell}\right)_{\Omega_{i}} \\
& +\sum_{j \in \mathcal{I}_{i}} \frac{1}{\alpha_{i j}}\left(\mathbf{e}_{i}^{\ell} \cdot \mathbf{n}_{i j}, \mathbf{e}_{i}^{\ell} \cdot \mathbf{n}_{i j}\right)_{\Gamma_{i j}}+\sum_{j \in \mathcal{I}_{i}} \frac{1}{\beta_{i j}}\left(\mathbf{e}_{i}^{\ell} \times \mathbf{n}_{i j}, \mathbf{e}_{i}^{\ell} \times \mathbf{n}_{i j}\right)_{\Gamma_{i j}}  \tag{6.6}\\
& =\sum_{j \in \mathcal{I}_{i}} \frac{1}{\alpha_{i j}}\left(h_{i j}^{\ell-1}, \mathbf{e}_{i}^{\ell} \cdot \mathbf{n}_{i j}\right)_{\Gamma_{i j}}+\sum_{j \in \mathcal{I}_{i}} \frac{1}{\beta_{i j}}\left(\zeta_{i j}^{\ell-1}, \mathbf{e}_{i}^{\ell} \times \mathbf{n}_{i j}\right)_{\Gamma_{i j}} .
\end{align*}
$$

On the boundary $\Gamma_{i j}, j \in \mathcal{I}_{i}$, replacing (6.5) into (6.6), one gets

$$
\begin{aligned}
\left\langle\partial_{t} \mathbf{e}_{i}^{\ell}, \mathbf{e}_{i}^{\ell}\right\rangle_{V_{i}^{\prime}, V_{i}} & +\nu\left(\nabla \mathbf{e}_{i}^{\ell}, \nabla \mathbf{e}_{i}^{\ell}\right)_{\Omega_{i}}+\frac{1}{4} \sum_{j \in \mathcal{I}_{i}} \frac{1}{\alpha_{i j}}\left(h_{i j}^{\ell-1}-h_{j i}^{\ell}, h_{i j}^{\ell-1}-h_{j i}^{\ell}\right)_{\Gamma_{i j}} \\
& +\frac{1}{4} \sum_{j \in \mathcal{I}_{i}} \frac{1}{\beta_{i j}}\left(\zeta_{i j}^{\ell-1}-\zeta_{j i}^{\ell}, \zeta_{i j}^{\ell-1}-\zeta_{j i}^{\ell}\right)_{\Gamma_{i j}} \\
& =\frac{1}{2} \sum_{j \in \mathcal{I}_{i}} \frac{1}{\alpha_{i j}}\left(h_{i j}^{\ell-1}, h_{i j}^{\ell-1}-h_{j i}^{\ell}\right)_{\Gamma_{i j}}+\frac{1}{2} \sum_{j \in \mathcal{I}_{i}} \frac{1}{\beta_{i j}}\left(\zeta_{i j}^{\ell-1}, \zeta_{i j}^{\ell-1}-\zeta_{j i}^{\ell}\right)_{\Gamma_{i j}}
\end{aligned}
$$

or equivalently

$$
\begin{array}{r}
\left\langle\partial_{t} \mathbf{e}_{i}^{\ell}, \mathbf{e}_{i}^{\ell}\right\rangle_{V_{i}^{\prime}, V_{i}}+\nu\left\|\nabla \mathbf{e}_{i}^{\ell}\right\|_{\Omega_{i}}^{2}+\frac{1}{4} \sum_{j \in \mathcal{I}_{i}} \frac{1}{\alpha_{i j}}\left\|h_{j i}^{\ell}\right\|_{\Gamma_{i j}}^{2}+\frac{1}{4} \sum_{j \in \mathcal{I}_{i}} \frac{1}{\beta_{i j}}\left\|\zeta_{j i}^{\ell}\right\|_{\Gamma_{i j}}^{2} \\
=\frac{1}{4} \sum_{j \in \mathcal{I}_{i}} \frac{1}{\alpha_{i j}}\left\|h_{i j}^{\ell-1}\right\|_{\Gamma_{i j}}^{2}+\frac{1}{4} \sum_{j \in \mathcal{I}_{i}} \frac{1}{\beta_{i j}}\left\|\zeta_{i j}^{\ell-1}\right\|_{\Gamma_{i j}}^{2} \tag{6.7}
\end{array}
$$

(recall that notation $\|\cdot\|_{D}$ corresponds to the $L^{2}(D)$-norm for any set $D$ ). Adapting (3.6) to $\Omega_{i}$, integrating (6.7) on ( $0, T$ ), and using that $\mathbf{e}_{i}^{\ell}(0)=0$, we get

$$
\begin{align*}
&\left\|\mathbf{e}_{i}^{\ell}(T)\right\|_{\Omega_{i}}^{2}+2 \nu \int_{0}^{T}\left\|\nabla \mathbf{e}_{i}^{\ell}\right\|_{\Omega_{i}}^{2}+\sum_{j \in \mathcal{I}_{i}} \frac{1}{2 \alpha_{i j}} \int_{0}^{T}\left\|h_{j i}^{\ell}\right\|_{\Gamma_{i j}}^{2}+\sum_{j \in \mathcal{I}_{i}} \frac{1}{2 \beta_{i j}} \int_{0}^{T}\left\|\zeta_{j i}^{\ell}\right\|_{\Gamma_{i j}}^{2} \\
&=\sum_{j \in \mathcal{I}_{i}} \frac{1}{2 \alpha_{i j}} \int_{0}^{T}\left\|h_{i j}^{\ell-1}\right\|_{\Gamma_{i j}}^{2}+\sum_{j \in \mathcal{I}_{i}} \int_{0}^{T} \frac{1}{2 \beta_{i j}}\left\|\zeta_{i j}^{\ell-1}\right\|_{\Gamma_{i j}}^{2} . \tag{6.8}
\end{align*}
$$

Then, summing with respect to $i$, from 1 to $M$, we get

$$
\sum_{i=1}^{M}\left\|\mathbf{e}_{i}^{\ell}(., T)\right\|_{\Omega_{i}}^{2}+2 \nu \sum_{i=1}^{M} \int_{0}^{T}\left\|\nabla \mathbf{e}_{i}^{\ell}\right\|_{\Omega_{i}}^{2}+E_{R}^{\ell}=E_{R}^{\ell-1}
$$

where $E_{R}^{\ell}=\sum_{i=1}^{M} \sum_{j \in \mathcal{I}_{i}} \frac{1}{2 \beta_{i j}} \int_{0}^{T}\left\|\zeta_{i j}^{\ell}\right\|_{\Gamma_{i j}}^{2}+\sum_{i=1}^{M} \sum_{j \in \mathcal{I}_{i}} \frac{1}{2 \alpha_{i j}} \int_{0}^{T}\left\|h_{i j}^{\ell}\right\|_{\Gamma_{i j}}^{2}$.
Summing now with respect to $\ell$, from 1 to $L$, we obtain

$$
\sum_{\ell=1}^{L} \sum_{i=1}^{M}\left\|\mathbf{e}_{i}^{\ell}(., T)\right\|_{\Omega_{i}}^{2}+2 \nu \sum_{\ell=1}^{L} \sum_{i=1}^{M} \int_{0}^{T}\left\|\nabla \mathbf{e}_{i}^{\ell}\right\|_{\Omega_{i}}^{2}(t) d t+E_{R}^{L}=E_{R}^{0}
$$

As $E_{R}^{L} \geq 0$ for all $L$, the sums $\sum_{\ell=1}^{L} \sum_{i=1}^{M}\left\|\mathbf{e}_{i}^{\ell}(., T)\right\|_{\Omega_{i}}^{2}$ and $\sum_{\ell=1}^{L} \sum_{i=1}^{M} \int_{0}^{T}\left\|\nabla \mathbf{e}_{i}^{\ell}\right\|_{\Omega_{i}}^{2}$ are bounded; hence $\left\|\mathbf{e}_{i}^{\ell}(T)\right\|_{\Omega_{i}}^{2}$ and $\int_{0}^{T}\left\|\nabla \mathbf{e}_{i}^{\ell}\right\|_{\Omega_{i}}^{2}(t) d t$ tend to 0 when $\ell \rightarrow \infty$.
In addition, in (6.8), we can integrate on $(0, t)$ instead of $(0, T)$, and we get for all $t \in(0, T)$

$$
\sum_{\ell=1}^{L} \sum_{i=1}^{M}\left\|\mathbf{e}_{i}^{\ell}(t)\right\|_{\Omega_{i}}^{2} \leq E_{R}^{0}
$$

This first leads to the convergence of $\left\|\mathbf{e}_{i}^{\ell}(t)\right\|_{\Omega_{i}}$ to 0 for all $t$ and thus to the convergence of $\mathbf{e}_{i}^{\ell}$ to 0 in $\mathcal{C}^{0}\left([0, T], H_{i}\right)$, but also to the fact that, integrating on $(0, T)$, it holds that

$$
\sum_{\ell=1}^{L} \sum_{i=1}^{M} \int_{0}^{T}\left\|\mathbf{e}_{i}^{\ell}(t)\right\|_{\Omega_{i}}^{2} d t \leq T E_{R}^{0}
$$

This implies that $\int_{0}^{T}\left\|\mathbf{e}_{i}^{\ell}(t)\right\|_{\Omega_{i}}^{2} d t$ tends to 0 when $\ell \rightarrow+\infty$. Then, summing with $\int_{0}^{T}\left\|\nabla \mathbf{e}_{i}^{\ell}\right\|_{\Omega_{i}}^{2}(t) d t$ that also tends to 0 , we have that $\int_{0}^{T}\left\|\mathbf{e}_{i}^{\ell}(t)\right\|_{\left[H^{1}\left(\Omega_{i}\right)\right]^{2}}^{2} d t$ tends to 0 , or, in other words, that $\mathbf{e}_{i}^{\ell}$ tends to 0 in $L^{2}\left((0, T), V_{i}\right)$, for $i \in \llbracket 1, M \rrbracket$.
Now, we prove a convergence result for the pressure. We set $P(t)=\int_{0}^{t} p(s) d s$ and $P_{i}=\left.P\right|_{\Omega_{i}}$ and denote the error by $D_{i}^{\ell}(t)=\left(P_{i}^{\ell}-P_{i}\right)(t), i \in \llbracket 1, M \rrbracket$. Then we can state the following result.

Corollary 6.3. Let all hypotheses of Theorem 6.2 be satisfied. Then for all $t \in[0, T]$ it holds that $\left\|D_{i}^{\ell}(t)-\frac{1}{\left|\Omega_{i}\right|} \int_{\Omega_{i}} D_{i}^{\ell}(t)\right\|_{\Omega_{i}} \rightarrow 0$ when $\ell \rightarrow \infty$.

Proof. Let $i \in \llbracket 1, M \rrbracket$. As $\left(\mathbf{u}_{i}, p_{i}\right)$ is the strong solution of the Robin problem with boundary conditions $g_{i j}, \xi_{i j}, j \in \mathcal{I}_{i}$, then $P_{i}$ verifies a variational formulation similar to (4.6) : $\forall \mathbf{v} \in X_{i}$ it holds

$$
\begin{equation*}
\left(\mathbf{u}_{i}(t), \mathbf{v}\right)_{\Omega_{i}}-\left(\mathbf{u}_{0, i}, \mathbf{v}\right)_{\Omega_{i}}+\int_{0}^{t} a_{i}\left(\mathbf{u}_{i}(s), \mathbf{v}\right) d s-\left(P_{i}(t), \nabla \cdot \mathbf{v}\right)_{\Omega_{i}}-\int_{0}^{t} c_{i}(s, \mathbf{v}) d s=0 \tag{6.9}
\end{equation*}
$$

Then, from (4.6) and (6.9), taking the test function $\mathbf{v} \in\left[H_{0}^{1}\left(\Omega_{i}\right)\right]^{2} \subset X_{i}$, the boundary terms in $c_{i}^{\ell}(s, \mathbf{v})$ and $c_{i}(s, \mathbf{v})$ vanish and then $c_{i}^{\ell}(s, \mathbf{v})-c_{i}(s, \mathbf{v})$ also vanishes. Then we get

$$
\left(D_{i}^{\ell}(t), \nabla \cdot \mathbf{v}\right)_{\Omega_{i}}=\left(\mathbf{e}_{i}^{\ell}(t), \mathbf{v}\right)_{\Omega_{i}}+\int_{0}^{t} a_{i}\left(\mathbf{e}_{i}^{\ell}(s), \mathbf{v}\right) d s, \forall \mathbf{v} \in\left[H_{0}^{1}\left(\Omega_{i}\right)\right]^{2} .
$$

As $(c, \nabla \cdot \mathbf{v})_{\Omega_{i}}=0$ for all constants $c$ and $\mathbf{v} \in\left[H_{0}^{1}\left(\Omega_{i}\right)\right]^{2}$, the above formulation implies that $\forall \mathbf{v} \in\left[H_{0}^{1}\left(\Omega_{i}\right)\right]^{2}$

$$
\left(D_{i}^{\ell}(t)-\frac{1}{\left|\Omega_{i}\right|} \int_{\Omega_{i}} D_{i}^{\ell}(t), \nabla \cdot \mathbf{v}\right)_{\Omega_{i}}=\left(\mathbf{e}_{i}^{\ell}(t), \mathbf{v}\right)_{\Omega_{i}}+\int_{0}^{t} a_{i}\left(\mathbf{e}_{i}^{\ell}(s), \mathbf{v}\right) d s
$$

Since $\left(D_{i}^{\ell}-\frac{1}{\left|\Omega_{i}\right|} \int_{\Omega_{i}} D_{i}^{\ell}\right) \in L_{0}^{2}\left(\Omega_{i}\right)=\left\{p \in L^{2}\left(\Omega_{i}\right), \int_{\Omega_{i}} p=0\right\}, i \in \llbracket 1, M \rrbracket$, from the inf-sup condition there exists $\gamma_{3}$ s.t.

$$
\left\|D_{i}^{\ell}-\frac{1}{\left|\Omega_{i}\right|} \int_{\Omega_{i}} D_{i}^{\ell}\right\|_{\Omega_{i}} \leq \frac{1}{\gamma_{3}} \sup _{\mathbf{v} \in\left[H_{0}^{1}\left(\Omega_{i}\right)\right]^{2}} \frac{\left|\left(\mathbf{e}_{i}^{\ell}(t), \mathbf{v}\right)_{\Omega_{i}}+\int_{0}^{t} a_{i}\left(\mathbf{e}_{i}^{\ell}(s), \mathbf{v}\right) d s\right|}{\left.\|\mathbf{v}\|_{\left[H_{0}^{1}\left(\Omega_{i}\right)\right.}\right]^{2}} .
$$

We apply again the continuity of $a_{i}(.,$.

$$
\left|\int_{0}^{t} a_{i}\left(\mathbf{e}_{i}^{\ell}(s), \mathbf{v}\right) d s\right| \leq M_{i} \int_{0}^{t}\left\|\mathbf{e}_{i}^{\ell}(s)\right\|_{X_{i}}\|\mathbf{v}\|_{X_{i}} d s \leq M_{i}\|\mathbf{v}\|_{\left[H_{0}^{1}\left(\Omega_{i}\right)\right]^{2}} \sqrt{T}\left\|\mathbf{e}_{i}^{\ell}\right\|_{L^{2}\left((0, T), X_{i}\right)}
$$

as well as the Cauchy-Schwarz and Poincaré inequalities on $\left(\mathbf{e}_{i}^{\ell}(t), \mathbf{v}\right)_{\Omega_{i}}$, we get

$$
\left\|D_{i}^{\ell}-\frac{1}{\left|\Omega_{i}\right|} \int_{\Omega_{i}} D_{i}^{\ell}\right\|_{\Omega_{i}} \leq \frac{1}{\gamma_{3}}\left[C_{P_{i}}\left\|\mathbf{e}_{i}^{\ell}(t)\right\|_{\Omega_{i}}+M_{i} \sqrt{T}\left\|\mathbf{e}_{i}^{\ell}\right\|_{L^{2}\left((0, T), X_{i}\right)}\right]
$$

with $C_{P_{i}}$ the Poincaré constant of $\Omega_{i}$. From the convergence of the velocity, we get the corollary.

Remark 6.4. Corollary 6.3 tells us that, when $\ell$ grows, the (time primitive of the) pressure converges to 0 , up to constant values in space, possibly depending on the subdomain $\Omega_{i}$ and iteration count $\ell$. And, indeed, numerical results given in Section 10 show that pressure iterates do not converge to the monodomain solution, unless a correction is applied, which is the object of the next Section.
7. Recovering the pressure. Let us introduce the notation $\langle p\rangle_{\mathcal{O}}=\frac{1}{|\mathcal{O}|} \int_{\mathcal{O}} p d x$ for the mean value of a function on a domain $\mathcal{O}$ (whatever the space dimension of $\mathcal{O}$ ).

We set $d_{i}^{\ell}:=p_{i}-p_{i}^{\ell}, i \in \llbracket 1, M \rrbracket$, and recall that $h_{i j}^{\ell}$ is defined in (6.3).
Hypothesis 7.1. In this section, we suppose that, for a.e $t \in(0, T)$
$\bullet\left\|d_{i}^{\ell}-\left\langle d_{i}^{\ell}\right\rangle_{\Omega_{i}}\right\|_{\Omega_{i}} \longrightarrow 0$ for all $i$ when $\ell \longrightarrow+\infty$

- $\left(\left\langle d_{i}^{\ell}\right\rangle_{\Gamma_{i j}}-\left\langle d_{i}^{\ell}\right\rangle_{\Omega_{i}}\right)$ tends to 0 for all $j \in \mathcal{I}_{i}$, for all $i$, when $\ell \longrightarrow+\infty$
- $\left(\left\langle h_{i j}^{\ell}\right\rangle_{\Gamma_{i j}}+\alpha_{i j}\left\langle d_{i}^{\ell}\right\rangle_{\Gamma_{i j}}\right) \longrightarrow 0$ for all $j \in \mathcal{I}_{i}$, for all $i$, when $\ell \longrightarrow+\infty$

Remark 7.2. The above hypothesis can be implied from stronger assumptions on the regularity and convergence of the velocity. Indeed, suppose that $\left(\mathbf{e}_{i}^{\ell}, d_{i}^{\ell}\right)$ is the strong solution of the following Robin problem

$$
\begin{aligned}
\partial_{t} \mathbf{e}_{i}^{\ell}-\nu \Delta \mathbf{e}_{i}^{\ell}+\nabla d_{i}^{\ell} & =0 & & \text { in }
\end{aligned} \quad \Omega_{i} \times(0, T), ~ \begin{array}{rlr}
\nabla \cdot \mathbf{e}_{i}^{\ell} & =0 & \\
\text { in } & \Omega_{i} \times(0, T) \\
\mathbf{e}_{i}^{\ell}(., t=0) & =0 & \\
\text { in } & \Omega_{i} \\
\mathbf{e}_{i}^{\ell} & =0 & \\
\text { on } & & \left(\partial \Omega \cap \partial \Omega_{i}\right) \times(0, T) \\
\alpha_{i j}\left(\nu \partial_{\mathbf{n}_{i j}} \mathbf{e}_{i}^{\ell} \cdot \mathbf{n}_{i j}-d_{i}^{\ell}\right)+\mathbf{e}_{i}^{\ell} \cdot \mathbf{n}_{i j} & =h_{i j}^{\ell} & \\
\text { on } & & \Gamma_{i j} \times(0, T) \\
\beta_{i j} \nu \partial_{\mathbf{n}_{i j}} \mathbf{e}_{i}^{\ell} \times \mathbf{n}+\mathbf{e}_{i}^{\ell} \times \mathbf{n}_{i j} & =\zeta_{i j}^{\ell} & \\
\text { on } & & \Gamma_{i j} \times(0, T)
\end{array}
$$

with the following convergence

$$
\left\|\mathbf{e}_{i}^{\ell}\right\|_{L^{\infty}\left((0, T),\left[H^{2}\left(\Omega_{i}\right)\right]^{2}\right)} \longrightarrow 0, \quad\left\|\partial_{t} \mathbf{e}_{i}^{\ell}\right\|_{L^{\infty}\left((0, T),\left[L^{2}\left(\Omega_{i}\right)\right]^{2}\right)} \longrightarrow 0
$$

From this, we get, for a.e. $t \in(0, T),\left\|\nabla d_{i}^{\ell}(t)\right\|_{\Omega_{i}} \longrightarrow 0$, which implies the first and second items in Hypothesis 7.1. This also implies the convergence of trace of the velocity: for a.e. $t \in(0, T)$, we have $\left\|\alpha_{i j} \nu \partial_{\mathbf{n}_{i j}} \mathbf{e}_{i}^{\ell}(t) \cdot \mathbf{n}_{i j}+\mathbf{e}_{i}^{\ell}(t) \cdot \mathbf{n}_{i j}\right\|_{\Gamma_{i j}} \longrightarrow 0$ that leads to the third item in Hypothesis 7.1.
One can rewrite the three items in Hypothesis 7.1 on the error as follows : when $\ell \longrightarrow+\infty, \forall i \in \llbracket 1, M \rrbracket$,

$$
\begin{array}{r}
\left\|\left(p_{i}^{\ell}-p_{i}\right)-\left(\left\langle p_{i}^{\ell}\right\rangle_{\Omega_{i}}-\left\langle p_{i}\right\rangle_{\Omega_{i}}\right)\right\|_{\Omega_{i}} \longrightarrow 0, \\
\left(\left\langle p_{i}^{\ell}-p_{i}\right\rangle_{\Gamma_{i j}}\right)-\left(\left\langle p_{i}^{\ell}-p_{i}\right\rangle_{\Omega_{i}}\right) \longrightarrow 0, \quad \forall j \in \mathcal{I}_{i}, \\
{\left[\left\langle g_{i j}^{\ell}\right\rangle_{\Gamma_{i j}}-\left\langle g_{i j}\right\rangle_{\Gamma_{i j}}\right]+\alpha_{i j}\left\langle p_{i}^{\ell}-p_{i}\right\rangle_{\Gamma_{i j}} \longrightarrow 0, \quad \forall j \in \mathcal{I}_{i} .} \tag{7.3}
\end{array}
$$

Expression (7.1) shows that $p_{i}^{\ell}(t)$ will tend to $p_{i}(t)$ if and only if the mean-value of $p_{i}^{\ell}(t)$ on $\Omega_{i}$ tends to the mean value of $p_{i}(t)$. However, no constraint was imposed on the mean-value of $p_{i}^{\ell}(t)$ in the algorithm, since, thanks to the Robin boundary conditions, such constraint is not necessary to obtain local well-posed problems at each iteration. In Section 5, we observed cases in which $p_{i}^{\ell}$ does not converge to the monodomain solution $p_{i}$. In this section, we build a modified pressure $\tilde{p}_{i}^{\ell}$ such that $\tilde{p}_{i}^{\ell}(t)$ tends to $p_{i}(t)$ in $L^{2}\left(\Omega_{i}\right), i=1, \ldots, M$.
Let us denote $X_{i}(t):=\left\langle p_{i}(t)\right\rangle_{\Omega_{i}}, \forall i \in \llbracket 1, M \rrbracket$. Then, using this notation, (7.1) reads

$$
\begin{equation*}
\left\|\left(p_{i}^{\ell}(t)-\left\langle p_{i}^{\ell}(t)\right\rangle_{\Omega_{i}}+X_{i}(t)\right)-p_{i}(t)\right\|_{L^{2}\left(\Omega_{i}\right)} \longrightarrow 0 \text { when } \ell \rightarrow \infty \tag{7.4}
\end{equation*}
$$

From (7.4), we see that $\left(p_{i}^{\ell}(t)-\left\langle p_{i}^{\ell}(t)\right\rangle_{\Omega_{i}}+X_{i}(t)\right)$ is the right approximation to calculate at each iteration since it tends to $p_{i}(t)$. However, we do not know how to calculate it because $X_{i}$ is not known. A similar question was raised in the thesis of Lissoni [33, Theorem IV.3.9] at the discrete level, within a Schwarz algorithm applied at each time step of a time marching scheme for the numerical approximation of the incompressible Navier-Stokes equations.
We introduce below a new quantity $Y_{i}^{\ell}(t)$, fully computable at any given iteration $\ell$, that tends to $X_{i}(t)$ when $\ell$ tends to infinity, from which we will define the modified pressure $\tilde{p}_{i}^{\ell}$.

To ease the presentation, we shall set $\left|\Gamma_{i j}\right|=0, \alpha_{i j}=0$ and $g_{i j}^{\ell}=0$ if $j \notin \mathcal{I}_{i}$. Moreover, we introduce the constant matrix

$$
A=\left(a_{i j}\right)_{1 \leq i, j \leq M}, \quad \text { with } \quad a_{i i}=\sum_{j=1, j \neq i}^{M}\left|\Gamma_{i j}\right| \alpha_{i j}, \text { and } a_{i j}=-\left|\Gamma_{j i}\right| \alpha_{j i} \text { if } j \neq i
$$

together with the constant vector $C=\left(\left|\Omega_{1}\right|,\left|\Omega_{2}\right|, \ldots,\left|\Omega_{M}\right|\right)$ and the sequence of vectors $\left(B^{\ell}\right)_{\ell}$, with $B^{\ell}=\left(B_{1}^{\ell}, B_{2}^{\ell}, \ldots, B_{M}^{\ell}\right)^{t}$ defined as

$$
B_{i}^{\ell}=\sum_{j=1}^{M}\left|\Gamma_{i j}\right|\left[\left\langle g_{i j}^{\ell}\right\rangle_{\Gamma_{i j}}+\alpha_{i j}\left\langle p_{i}^{\ell}\right\rangle_{\Omega_{i}}\right]-\sum_{j=1}^{M}\left|\Gamma_{j i}\right|\left[\left\langle g_{j i}^{\ell}\right\rangle_{\Gamma_{j i}}+\alpha_{j i}\left\langle p_{j}^{\ell}\right\rangle_{\Omega_{j}}\right] .
$$

Theorem 7.3. Assume that $\alpha_{i j}=\alpha_{j i}, \forall(i, j)$. We have the following properties
(i) For all $\ell$, the following system

$$
\begin{align*}
& A Y^{\ell}=B^{\ell} \\
& C Y^{\ell}=0 \tag{7.5}
\end{align*}
$$

has a unique solution $Y^{\ell} \in \mathbb{R}^{M}$.
(ii) Moreover, we have $Y^{\ell} \rightarrow X:=\left(X_{1}, X_{2}, \cdots, X_{M}\right)$ in $\mathbb{R}^{M}$, and for all $t$

$$
\begin{equation*}
\left\|\tilde{p}_{i}^{\ell}-p_{i}\right\|_{L^{2}\left(\Omega_{i}\right)} \longrightarrow 0, \text { when } \ell \rightarrow \infty, \text { with } \tilde{p}_{i}^{\ell}(t):=p_{i}^{\ell}(t)-\left\langle p_{i}^{\ell}(t)\right\rangle_{\Omega_{i}}+Y_{i}^{\ell}(t) \tag{7.6}
\end{equation*}
$$

Proof of (i). The proof of Theorem 7.3-(i) relies on two main steps:
(a) Existence of solutions to the system $A Y^{\ell}=B^{\ell}$,
(b) Existence and uniqueness of a solution to system (7.5) thanks to the additional constraint $C Y^{\ell}=0$.

Let us start with (a). Because $\alpha_{i j}=\alpha_{j i}$, it holds that $A$ is symmetric and then existence of at least one solution to the system $A Y^{\ell}=B^{\ell}$ is equivalent to proving that $B^{\ell} \in \operatorname{Im}(A)=(\operatorname{Ker}(A))^{\perp}$. Thus, we start with the determination of $\operatorname{Ker}(A)$.

Let $Y=\left(Y_{1}, Y_{2}, \ldots, Y_{M}\right)^{t} \in \operatorname{Ker}(A)$. Then, we have $\sum_{j=1}^{M} a_{i j} Y_{j}=0, \forall i \in \llbracket 1, M \rrbracket$. As $\alpha_{i j}=\alpha_{j i}$, we have $a_{i i}=-\sum_{j=1, j \neq i}^{M} a_{i j}$, which implies

$$
0=\sum_{j=1}^{M} a_{i j} Y_{j} Y_{i}=\left(\sum_{j=1, j \neq i}^{M} a_{i j} Y_{j} Y_{i}\right)+a_{i i} Y_{i}^{2}=\sum_{j=1, j \neq i}^{M} a_{i j}\left(Y_{j} Y_{i}-Y_{i}^{2}\right)
$$

Summing the above expression in $i$, and using that $a_{i j}=a_{j i}$, we obtain

$$
\sum_{i=1}^{M} \sum_{j=1, j \neq i}^{M} a_{i j}\left(Y_{j} Y_{i}-Y_{i}^{2}\right)=\sum_{i<j} a_{i j}\left(Y_{i}-Y_{j}\right)^{2}=0
$$

As $a_{i j} \leq 0$ for all $(i, j)$ with $i \neq j$, and $a_{i j}<0$ as soon as subdomains $i$ and $j$ are neighbours, this implies that $Y_{i}=Y_{j}$ for any pair of neighbouring subdomains $i$ and $j$. Since $\Omega$ is connected, this finally implies that all $Y_{i}$ are equal i.e. $\operatorname{Ker}(A)=\operatorname{span}(\mathbf{e})$ with $\mathbf{e}=(1,1, \ldots, 1,1)$. Then, $B^{\ell} \in(\operatorname{Ker}(A))^{\perp}$ is equivalent to $B^{\ell} \cdot \mathbf{e}=\sum_{i=1}^{M} B_{i}^{\ell}=0$. This is proved in the following way:

$$
\sum_{i=1}^{M} B_{i}^{\ell}=\sum_{i=1}^{M}\left[\sum_{j=1}^{M}\left|\Gamma_{i j}\right|\left(\left\langle g_{i j}^{\ell}\right\rangle_{\Gamma_{i j}}+\alpha_{i j}\left\langle p_{i}^{\ell}\right\rangle_{\Omega_{i}}\right)-\sum_{j=1}^{M}\left|\Gamma_{j i}\right|\left(\left\langle g_{j i}^{\ell}\right\rangle_{\Gamma_{j i}}+\alpha_{j i}\left\langle p_{j}^{\ell}\right\rangle_{\Omega_{j}}\right)\right]
$$

Denoting $\Delta_{i j}:=\left|\Gamma_{i j}\right|\left(\left\langle g_{i j}^{\ell}\right\rangle_{\Gamma_{i j}}+\alpha_{i j}\left\langle p_{i}^{\ell}\right\rangle_{\Omega_{i}}\right)$, we obtain

$$
\sum_{i=1}^{M} B_{i}^{\ell}=\sum_{i=1}^{M} \sum_{j=1}^{M} \Delta_{i j}-\sum_{i=1}^{M} \sum_{j=1}^{M} \Delta_{j i}=0
$$

Let us now turn to (b). From (a), we know that there exists at least a solution to $A Y=B$; we let $Y^{*}$ be such a solution. All other solutions may be written as $Y=Y^{*}+\mu \mathbf{e}$, with $\mu \in \mathbb{R}$. Existence of a solution to (7.5) follows from the fact that $C \mathbf{e}=|\Omega| \neq 0$ : Choosing $\mu=-\frac{1}{|\Omega|} C Y^{*}$ leads to $C Y=C Y^{*}+\mu C \mathbf{e}=0$ and then $Y$ solves (7.5). As far as uniqueness is concerned, let $Y_{1}$ and $Y_{2}$ be two solutions of (7.5); since $\left(Y_{1}-Y_{2}\right) \in \operatorname{Ker}(A)$, then $\left(Y_{1}-Y_{2}\right)=\tau \mathbf{e}$, with $\tau \in \mathbb{R}$. Since $\tau|\Omega|=\tau C \mathbf{e}=C\left(Y_{1}-Y_{2}\right)=0$ it follows that $\tau=0$ and $Y_{1}=Y_{2}$. This ends the proof of Theorem 7.3-(i).

Proof of Theorem 7.3-(ii). It relies on the two main results:
(c) $B^{\ell} \rightarrow A X$ in $\mathbb{R}^{M}$,
(d) $C X=0$.

Let us prove (c): from the divergence-free property of $\mathbf{u}_{i}$, we have

$$
\begin{equation*}
0=\int_{\Omega_{i}} \nabla \cdot \mathbf{u}_{i}=\int_{\partial \Omega_{i}} \mathbf{u}_{i} \cdot \mathbf{n}_{\partial \Omega_{i}}=\sum_{j \in \mathcal{I}_{i}} \int_{\Gamma_{i j}} \mathbf{u}_{i} \cdot \mathbf{n}_{i j} \tag{7.7}
\end{equation*}
$$

Moreover, from the definition of $g_{i j}$ in (6.1a) and the physical transmission conditions (2.4), we have

$$
\begin{equation*}
\left|\Gamma_{i j}\right|\left\langle g_{i j}\right\rangle_{\Gamma_{i j}}-\left|\Gamma_{j i}\right|\left\langle g_{j i}\right\rangle_{\Gamma_{j i}}=\int_{\Gamma_{i j}}\left(g_{i j}-g_{j i}\right)=2 \int_{\Gamma_{i j}} \mathbf{u}_{i} \cdot \mathbf{n}_{i j} \tag{7.8}
\end{equation*}
$$

Hence, from (7.7) and (7.8) we get

$$
\begin{equation*}
\sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{i j}\right|\left\langle g_{i j}\right\rangle_{\Gamma_{i j}}=\sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{j i}\right|\left\langle g_{j i}\right\rangle_{\Gamma_{j i}} . \tag{7.9}
\end{equation*}
$$

Expression (7.3) is equivalent to

$$
\begin{equation*}
\left\langle g_{i j}^{\ell}\right\rangle_{\Gamma_{i j}}+\alpha_{i j}\left\langle p_{i}^{\ell}-p_{i}\right\rangle_{\Gamma_{i j}} \longrightarrow\left\langle g_{i j}\right\rangle_{\Gamma_{i j}} . \tag{7.10}
\end{equation*}
$$

From (7.2), we may replace $\left\langle p_{i}^{\ell}-p_{i}\right\rangle_{\Gamma_{i j}}$ by $\left\langle p_{i}^{\ell}-p_{i}\right\rangle_{\Omega_{i}}$ in (7.10), then multiply by $\left|\Gamma_{i j}\right|$ and sum over $j \in \mathcal{I}_{i}$ for a given $i$ to obtain

$$
\begin{equation*}
\sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{i j}\right|\left[\left\langle g_{i j}^{\ell}\right\rangle_{\Gamma_{i j}}+\alpha_{i j}\left\langle p_{i}^{\ell}-p_{i}\right\rangle_{\Omega_{i}}\right] \longrightarrow \sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{i j}\right|\left\langle g_{i j}\right\rangle_{\Gamma_{i j}} \tag{7.11}
\end{equation*}
$$

In exactly the same way, we also obtain

$$
\begin{equation*}
\sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{j i}\right|\left[\left\langle g_{j i}^{\ell}\right\rangle_{\Gamma_{j i}}+\alpha_{j i}\left\langle p_{j}^{\ell}-p_{j}\right\rangle_{\Omega_{j}}\right] \longrightarrow \sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{j i}\right|\left\langle g_{j i}\right\rangle_{\Gamma_{j i}} . \tag{7.12}
\end{equation*}
$$

Using (7.11), (7.12) and (7.9), we obtain

$$
\begin{aligned}
& \sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{i j}\right|\left[\left\langle g_{i j}^{\ell}\right\rangle_{\Gamma_{i j}}+\alpha_{i j}\left\langle p_{i}^{\ell}\right\rangle_{\Omega_{i}}-\alpha_{i j}\left\langle p_{i}\right\rangle_{\Omega_{i}}\right] \\
- & \sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{j i}\right|\left[\left\langle g_{j i}^{\ell}\right\rangle_{\Gamma_{j i}}+\alpha_{j i}\left\langle p_{j}^{\ell}\right\rangle_{\Omega_{j}}-\alpha_{j i}\left\langle p_{j}\right\rangle_{\Omega_{j}}\right] \longrightarrow 0,
\end{aligned}
$$

or equivalently

$$
\begin{aligned}
\sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{i j}\right|\left[\left\langle g_{i j}^{\ell}\right\rangle_{\Gamma_{i j}}+\right. & \left.\alpha_{i j}\left\langle p_{i}^{\ell}\right\rangle_{\Omega_{i}}\right]-\sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{j i}\right|\left[\left\langle g_{j i}^{\ell}\right\rangle_{\Gamma_{j i}}+\alpha_{j i}\left\langle p_{j}^{\ell}\right\rangle_{\Omega_{j}}\right] \\
& \longrightarrow \sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{i j}\right| \alpha_{i j}\left\langle p_{i}\right\rangle_{\Omega_{i}}-\sum_{j \in \mathcal{I}_{i}}\left|\Gamma_{j i}\right| \alpha_{j i}\left\langle p_{j}\right\rangle_{\Omega_{j}}
\end{aligned}
$$

This is exactly $B^{\ell} \longrightarrow A X$.
Let us now prove (d): We have

$$
\int_{\Omega} p_{i}=\sum_{i=1}^{M} \int_{\Omega_{i}} p_{i}=\sum_{i=1}^{M}\left|\Omega_{i}\right|\left\langle p_{i}\right\rangle_{\Omega_{i}}=0
$$

i.e. $C X=0$.

We now prove Theorem 7.3-(ii): From the solution $Y^{\ell}$ of (7.5) given by Theorem 7.3-(i), and from (c) and (d), we have $A\left(Y^{\ell}-X\right) \longrightarrow 0$ and $C\left(Y^{\ell}-X\right)=0$. Uniqueness of a solution to $A Z=B$ and $C Z=0$ as soon as $B$ is in $\operatorname{Im}(A)$ and finite dimension now imply that $\left(Y^{\ell}-X\right) \longrightarrow 0$ when $\ell \rightarrow \infty$. Then, from (7.4), with a triangle inequality, we get (7.6).

Remark 7.4. In the general case of $M$ subdomains, the calculation of $\tilde{p}_{i}^{\ell}$ is done only once, at the last OSWR iteration. It involves solving the coarse problem (7.5) when $M>2$, and is given by an explicit formula when $M=2$ (see Corollary 7.6), thus the cost of calculating the modified pressure is negligible.

Remark 7.5. Recovering the correct pressure could also be performed from the fact that $\nabla\left(p_{i}^{\ell}-p_{i}\right)$ tends to zero when $\ell \rightarrow \infty$. Indeed, for a given $\Omega_{i}$, choosing first an arbitrary point $\mathbf{x}_{i} \in \Omega_{i}$, then one may write

$$
p_{i}(\mathbf{x})=p_{i}\left(\mathbf{x}_{i}\right)+\left(\mathbf{x}-\mathbf{x}_{i}\right) \cdot \int_{0}^{1} \nabla p_{i}\left(\mathbf{x}_{i}+t\left(\mathbf{x}-\mathbf{x}_{i}\right)\right) d t \quad, \quad \forall \mathbf{x} \in \Omega_{i} .
$$

Then, one could replace $\nabla p_{i}$ by $\nabla p_{i}^{\ell}$ to obtain approximate values of the pressure at each point $\mathbf{x}$. However, this formula holds on a given subdomain $\Omega_{i}$. In order to relate values of the pressures in $\Omega_{i}$ to those in a neighboring subdomain $\Omega_{j}$ through this kind of formula, one needs to choose a point on the boundary $\Gamma_{i j}$ that will serve as the point $\mathbf{x}_{j}$ in the subdomain $\Omega_{j}$, and so on. At the discrete level, there are several drawbacks to that: this requires further communications between subdomains, the pressure gradient at the boundaries may not be easy to define (e.g. when the pressure is defined as a piecewise constant field like in the Crouzeix-Raviart finite element), and finally there are many ways to go from one cell to another in the mesh, and, due to round-off errors, this may lead to different evaluations of the pressure at a given cell in particular in very large scale computations.

In the two-subdomain case, we use the same notation as in Section 5. Then the calculation of $\tilde{p}_{i}^{\ell}$ can be done by the following explicit formula.

Corollary 7.6. Let $M=2, \alpha=\alpha_{1}=\alpha_{2}$, and define, for $i=1,2$ and $j=3-i$,

$$
\tilde{p}_{i}^{\ell}=p_{i}^{\ell}+\frac{\left|\Omega_{j}\right|}{|\Omega|}\left[\frac{1}{\alpha}\left(\left\langle g_{i}^{\ell}\right\rangle_{\Gamma}-\left\langle g_{j}^{\ell}\right\rangle_{\Gamma}\right)\right]-\frac{\left|\Omega_{i}\right|}{|\Omega|}\left\langle p_{i}^{\ell}\right\rangle_{\Omega_{i}}-\frac{\left|\Omega_{j}\right|}{|\Omega|}\left\langle p_{j}^{\ell}\right\rangle_{\Omega_{j}} .
$$

Then $\tilde{p}_{i}^{\ell}$ tends to $p_{i}$ when $\ell$ tends to infinity, for $i=1,2$.
Proof. For $M=2$ we have

$$
\begin{aligned}
& B_{1}^{\ell}=-B_{2}^{\ell}=|\Gamma|\left[\left\langle g_{1}^{\ell}\right\rangle_{\Gamma}+\alpha\left\langle p_{1}\right\rangle_{\Omega_{1}}\right]-|\Gamma|\left[\left\langle g_{2}^{\ell}\right\rangle_{\Gamma}+\alpha\left\langle p_{2}\right\rangle_{\Omega_{2}}\right], \\
& \mathbf{A}=\left[\begin{array}{cc}
\alpha|\Gamma| & -\alpha|\Gamma| \\
-\alpha|\Gamma| & \alpha|\Gamma|
\end{array}\right], \\
& C=\left[\left|\Omega_{1}\right|\right. \\
&\left.\left|\Omega_{2}\right|\right] .
\end{aligned}
$$

System (7.5) for $M=2$ has a unique solution given by

$$
\begin{aligned}
& Y_{1}^{\ell}=\frac{\left|\Omega_{2}\right|}{|\Omega|}\left[\frac{1}{\alpha}\left(\left\langle g_{1}^{\ell}\right\rangle_{\Gamma}-\left\langle g_{2}^{\ell}\right\rangle_{\Gamma}\right)+\left(\left\langle p_{1}^{\ell}\right\rangle_{\Omega_{1}}-\left\langle p_{2}^{\ell}\right\rangle_{\Omega_{2}}\right)\right], \\
& Y_{2}^{\ell}=\frac{\left|\Omega_{1}\right|}{|\Omega|}\left[\frac{1}{\alpha}\left(\left\langle g_{2}^{\ell}\right\rangle_{\Gamma}-\left\langle g_{1}^{\ell}\right\rangle_{\Gamma}\right)+\left(\left\langle p_{2}^{\ell}\right\rangle_{\Omega_{2}}-\left\langle p_{1}^{\ell}\right\rangle_{\Omega_{1}}\right)\right] .
\end{aligned}
$$

From theorem 7.3, we have

$$
\begin{aligned}
& p_{1}^{\ell}-\left\langle p_{1}^{\ell}\right\rangle_{\Omega_{1}}+Y_{1}^{\ell}=p_{1}^{\ell}+\frac{\left|\Omega_{2}\right|}{|\Omega|}\left[\frac{1}{\alpha}\left(\left\langle g_{1}^{\ell}\right\rangle_{\Gamma}-\left\langle g_{2}^{\ell}\right\rangle_{\Gamma}\right)\right]-\frac{\left|\Omega_{1}\right|}{|\Omega|}\left\langle p_{1}^{\ell}\right\rangle_{\Omega_{1}}-\frac{\left|\Omega_{2}\right|}{|\Omega|}\left\langle p_{2}^{\ell}\right\rangle_{\Omega_{2}} \rightarrow p_{1} \\
& p_{2}^{\ell}-\left\langle p_{2}^{\ell}\right\rangle_{\Omega_{2}}+Y_{2}^{\ell}=p_{2}^{\ell}+\frac{\left|\Omega_{1}\right|}{|\Omega|}\left[\frac{1}{\alpha}\left(\left\langle g_{2}^{\ell}\right\rangle_{\Gamma}-\left\langle g_{1}^{\ell}\right\rangle_{\Gamma}\right)\right]-\frac{\left|\Omega_{1}\right|}{|\Omega|}\left\langle p_{1}^{\ell}\right\rangle_{\Omega_{1}}-\frac{\left|\Omega_{2}\right|}{|\Omega|}\left\langle p_{2}^{\ell}\right\rangle_{\Omega_{2}} \rightarrow p_{2} \square
\end{aligned}
$$

8. Convergence factor via Fourier transform. The aim of this section is to find a way to conveniently choose the parameters $(\alpha, \beta)$ that play an important role in the actual rate of convergence in numerical experiments.
Let $\Omega=\mathbb{R}^{2}$. We consider two subdomains $\Omega_{1}=(-\infty, 0) \times \mathbb{R}$ and $\Omega_{2}=(0,+\infty) \times \mathbb{R}$, as commonly done for the analysis of OSWR methods. To simplify notation, we set $\Gamma:=\Gamma_{12}=\Gamma_{21}=\{x=0\} \times \mathbb{R}$, and denote $\alpha_{12}$ and $\alpha_{21}$ by $\alpha_{1}$ and $\alpha_{2}$, respectively. We denote $\mathbf{u}=(u, v)$ the two components of the velocity and set $\mathbf{f}=\left(f_{x}, f_{y}\right)$. Recall here the Stokes problem

$$
\begin{array}{rlll}
\partial_{t} u-\nu \Delta u+\partial_{x} p & =f_{x} \\
\partial_{t} v-\nu \Delta v+\partial_{y} p & =f_{y} \\
\partial_{x} u+\partial_{y} v & =0 \\
u(., t=0) & =u_{0} \\
v(., t=0) & =v_{0} \\
u, v & \rightarrow 0
\end{array} \quad, \text { in } \quad, \quad \Omega \times(0, T) \quad \Omega, \text { when } \quad|(x, y)| \rightarrow+\infty .
$$

We write the algorithm for the errors using the same notation $(u, v, p)$, which means that, by linearity, we set $f_{x}=f_{y}=0$ and $u_{0}=v_{0}=0$. To avoid additional notation for the Robin terms, we write the OSWR algorithm as follows: starting with $u_{i}^{0}, v_{i}^{0}, p_{i}^{0}$, at step $\ell \geq 1$ and provided $u_{i}^{\ell-1}, v_{i}^{\ell-1}, p_{i}^{\ell-1}$ we solve

$$
\begin{array}{rlll}
\partial_{t} u_{i}^{\ell}-\nu \Delta u_{i}^{\ell}+\partial_{x} p_{i}^{\ell} & =0 \\
\partial_{t} v_{i}^{\ell}-\nu \Delta v_{i}^{\ell}+\partial_{y} p_{i}^{\ell} & =0 \\
\partial_{x} u_{i}^{\ell}+\partial_{y} v_{i}^{\ell} & =0 \\
u_{i}^{\ell}(., t=0) & =0 \\
v_{i}^{\ell}(., t=0) & & =0 & \\
u_{i}^{\ell}, v_{i}^{\ell} & \rightarrow 0 & & \\
& \text { in } & & \Omega_{i} \times(0, T) \\
\Omega_{i} \\
\text { when } & & |(x, y)| \rightarrow+\infty
\end{array}
$$

together with transmission condition on $\Gamma \times(0, T)$, for $i=1,2$ and $j=3-i$ :

$$
\begin{aligned}
\alpha_{i}\left(\nu \partial_{x} u_{i}^{\ell}-p_{i}^{\ell}\right)+(-1)^{i+1} u_{i}^{\ell} & =\alpha_{i}\left(\nu \partial_{x} u_{j}^{\ell-1}-p_{j}^{\ell-1}\right)+(-1)^{i+1} u_{j}^{\ell-1} \\
\nu \beta_{i} \partial_{x} v_{i}^{\ell}+(-1)^{i+1} v_{i}^{\ell} & =\nu \beta_{i} \partial_{x} v_{j}^{\ell-1}+(-1)^{i+1} v_{j}^{\ell-1}
\end{aligned}
$$

Let us consider the system in $\Omega_{1}$, and let $\ell \geq 1$. Taking the Fourier transform in time and in $y$-direction with time frequency $\omega$ and space frequency $k \neq 0$, and, for the sake of simplicity, keeping notation $u, v$ instead of $\hat{u}, \hat{v}$, we get

$$
\begin{align*}
i \omega u_{1}^{\ell}-\nu \partial_{x x} u_{1}^{\ell}+\nu k^{2} u_{1}^{\ell}+\partial_{x} p_{1}^{\ell} & =0  \tag{8.1a}\\
i \omega v_{1}^{\ell}-\nu \partial_{x x} v_{1}^{\ell}+\nu k^{2} v_{1}^{\ell}+i k p_{1}^{\ell} & =0  \tag{8.1b}\\
\partial_{x} u_{1}^{\ell}+i k v_{1}^{\ell} & =0 . \tag{8.1c}
\end{align*}
$$

By differentiating equation (8.1b) with respect to $x$, multiplying (8.1a) by ( $-i k$ ), and summing the resulting equations, and denoting $w_{1}^{\ell}:=\partial_{x} v_{1}^{\ell}-i k u_{1}^{\ell}$ the vorticity, we get the vorticity equation

$$
\begin{equation*}
i \omega w_{1}^{\ell}-\nu \partial_{x x} w_{1}^{\ell}+\nu k^{2} w_{1}^{\ell}=0 \tag{8.2}
\end{equation*}
$$

Denote by $\lambda=\sqrt{k^{2}+\frac{i \omega}{\nu}}$ with positive real part. As $w_{1}$ vanishes at $-\infty$, one gets

$$
\begin{equation*}
w_{1}^{\ell}=E^{\ell} \exp (\lambda x) \tag{8.3}
\end{equation*}
$$

Using the definition of $w_{1}$ and differentiating (8.1c), we get, for $u_{1}$

$$
\begin{equation*}
\partial_{x x} u_{1}^{\ell}-k^{2} u_{1}^{\ell}=-i k w_{1}^{\ell} \tag{8.4}
\end{equation*}
$$

The homogeneous equation associated to (8.4) has characteristic roots $\pm|k|$. As $u_{1}$ and $v_{1}$ vanish at $-\infty$, we only retain the root $|k|$. Given the form (8.3) of the righthand side of (8.4), its solution can be written under the form

$$
u_{1}^{\ell}=A^{\ell} \exp (|k| x)+B^{\ell} \exp (\lambda x)
$$

with $A^{\ell}, B^{\ell} \in \mathbb{C}$. Then, using (8.1c) and (8.1b), we get

$$
\begin{aligned}
v_{1}^{\ell} & =A^{\ell} \frac{i|k|}{k} \exp (|k| x)+B^{\ell} \frac{i \lambda}{k} \exp (\lambda x) \\
p_{1}^{\ell} & =-A^{\ell} \frac{i \omega}{|k|} \exp (|k| x)
\end{aligned}
$$

Similarly, for domain $\Omega_{2}$, there exist $C^{\ell}, D^{\ell} \in \mathbb{C}$ such that

$$
\begin{aligned}
u_{2}^{\ell} & =C^{\ell} \exp (-|k| x)+D^{\ell} \exp (-\lambda x) \\
v_{2}^{\ell} & =-C^{\ell} \frac{i|k|}{k} \exp (-|k| x)-D^{\ell} \frac{i \lambda}{k} \exp (-\lambda x) \\
p_{2}^{\ell} & =C^{\ell} \frac{i \omega}{|k|} \exp (-|k| x)
\end{aligned}
$$

Replacing the above expressions in the transmission conditions, we obtain

$$
\begin{aligned}
\alpha_{1}\left(\nu|k| A^{\ell}+\nu \lambda B^{\ell}+\frac{i \omega}{|k|} A^{\ell}\right)+A^{\ell}+B^{\ell} & = \\
\alpha_{1}\left(-\nu|k| C^{\ell-1}\right. & \left.-\nu \lambda D^{\ell-1}-\frac{i \omega}{|k|} C^{\ell-1}\right)+C^{\ell-1}+D^{\ell-1} \\
\nu \beta_{1}\left(i k A^{\ell}+\frac{i \lambda^{2}}{k} B^{\ell}\right)+\frac{i|k|}{k} A^{\ell}+\frac{i \lambda}{k} B^{\ell} & = \\
\nu \beta_{1}\left(i k C^{\ell-1}\right. & \left.+\frac{i \lambda^{2}}{k} D^{\ell-1}\right)-\frac{i|k|}{k} C^{\ell}-\frac{i \lambda}{k} D^{\ell-1}
\end{aligned}
$$

and

$$
\begin{aligned}
\alpha_{2}\left(-\nu|k| C^{\ell}-\nu \lambda D^{\ell}-\frac{i \omega}{|k|} C^{\ell}\right)-C^{\ell}-D^{\ell} & = \\
\alpha_{2}\left(\nu|k| A^{\ell-1}\right. & \left.+\nu \lambda B^{\ell-1}+\frac{i \omega}{|k|} A^{\ell-1}\right)-A^{\ell-1}-B^{\ell-1} \\
\nu \beta_{2}\left(i k C^{\ell}+\frac{i \lambda^{2}}{k} D^{\ell}\right)+\frac{i|k|}{k} C^{\ell}+\frac{i \lambda}{k} D^{\ell} & = \\
\nu \beta_{2}\left(i k A^{\ell-1}\right. & \left.+\frac{i \lambda^{2}}{k} B^{\ell-1}\right)-\frac{i|k|}{k} A^{\ell}-\frac{i \lambda}{k} B^{\ell-1}
\end{aligned}
$$

These transmission conditions can be written in matrix form as follows :
$\mathcal{M}\left(\alpha_{1}, \beta_{1}\right)\binom{A^{\ell}}{B^{\ell}}=\mathcal{N}\left(\alpha_{1}, \beta_{1}\right)\binom{C^{\ell-1}}{D^{\ell-1}}$ and $\mathcal{M}\left(\alpha_{2}, \beta_{2}\right)\binom{C^{\ell}}{D^{\ell}}=\mathcal{N}\left(\alpha_{2}, \beta_{2}\right)\binom{A^{\ell-1}}{B^{\ell-1}}$
where

$$
\mathcal{M}(\alpha, \beta):=\left[\begin{array}{cc}
1+\frac{\nu \alpha \lambda^{2}}{|k|} & 1+\alpha \nu \lambda  \tag{8.5}\\
\nu \beta k+\frac{|k|}{k} & \frac{\nu \beta \lambda^{2}}{k}+\frac{\lambda}{k}
\end{array}\right], \quad \mathcal{N}(\alpha, \beta):=\left[\begin{array}{cc}
1-\frac{\nu \alpha \lambda^{2}}{|k|} & 1-\alpha \nu \lambda \\
\nu \beta k-\frac{|k|}{k} & \frac{\nu \beta \lambda^{2}}{k}-\frac{\lambda}{k}
\end{array}\right] .
$$

This leads to the following recurrent formulation

$$
\begin{equation*}
\binom{A^{\ell}}{B^{\ell}}=\mathcal{R}\left(\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}\right)\binom{A^{\ell-2}}{B^{\ell-2}}, \quad \forall \ell \geq 2 \tag{8.6}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{R}\left(\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}\right)=\mathcal{M}^{-1}\left(\alpha_{1}, \beta_{1}\right) \mathcal{N}\left(\alpha_{1}, \beta_{1}\right) \mathcal{M}^{-1}\left(\alpha_{2}, \beta_{2}\right) \mathcal{N}\left(\alpha_{2}, \beta_{2}\right) . \tag{8.7}
\end{equation*}
$$

In view of (8.6), the convergence properties of the OSWR algorithm, and in particular its rate, will depend on the spectral radius of the matrix $\mathcal{R}$ defined in (8.7).
$\operatorname{Remark}$ 8.1. If one sets $\tilde{\alpha}:=\nu \alpha$ and $\tilde{\beta}:=\nu \beta$, as well as $\tilde{\omega}:=\frac{\omega}{\nu}$, then matrices $\mathcal{M}$ and $\mathcal{N}$ (defined in (8.5)), depend only on $\tilde{\alpha}, \tilde{\beta}$, on $\tilde{\omega}$ and on $k$. Thus, when $\nu$ varies, the convergence rate remains unchanged if $\tilde{\alpha}$ and $\tilde{\beta}$ are kept constant and if the range in which $\tilde{\omega}$ is considered does not change. As will be seen in Section 9, this is the case if $\nu \Delta t$ and $\nu T$ are kept unchanged. This observation coincides with the fact that the non-dimensional form of the Stokes equation is not modified when $\nu T$ is kept constant.

Remark 8.2. When $k$ tends to 0 , the spectral radius of the matrix $\mathcal{R}$ tends to 1 . This is coherent with what was observed in Section 5 and in Remarks 4.3 and 6.4, which led us to the pressure correction described in Section 7.

Remark 8.3. When $k$ and $\omega$ tend to $+\infty$, the spectral radius of the matrix $\mathcal{R}$ tends to 1 . This implies that analysing the iteration matrix does not help to prove the general convergence (for all frequencies) of the algorithm, and one always needs the energy estimate technique of Section 6 (for another example, see [10]).

Remark 8.4. In practical experiments, all equations are discretized in space and time. As far as space discretization is concerned, the solution of the discrete version of (8.2) remains close to (8.3) if the space discretization parameter is small enough with respect to $\sqrt{\frac{\nu}{\omega}}$; since $\omega$ is in practice bounded by $\frac{\pi}{\Delta t}$, we expect that the above Fourier analysis may remain close to practical experiments if the term $\sqrt{\nu \Delta t}$ is large enough compared to the space discretization parameter. This has indeed recently been observed for the heat equation in [2]. As far as time discretization is concerned, the inclusion of its effect in the convergence analysis of OSWR methods is a current topic of research, and is for example addressed in [15] where a $Z$ - transform is used and in [2], where a discrete-time analysis of the OSWR method is proposed. This issue is also addressed in Section 9.2.
9. Optimized Robin parameters. One can choose $\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}$ to minimize the convergence factor of the continuous OSWR algorithm, defined in the above section. Such parameters are called continuous optimized parameters. However, for the incompressible Stokes problem, we will see in the numerical experiments of Section 10 that better results can be obtained by minimizing the discrete-time counterpart of this convergence factor. The corresponding parameters are then called discrete-time optimized parameters. Both of these optimization procedures are described below.
9.1. Continuous optimized parameters. From Section 8 , the convergence factor is $\varrho\left(\mathcal{R}\left(\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}, k, \omega\right)\right)$, where $\mathcal{R}$ is defined in (8.7), and $\varrho(\mathcal{R})$ denotes the spectral radius of $\mathcal{R}$. While we have $\max _{(k, \omega) \in \mathbb{R}^{2}} \varrho\left(\mathcal{R}\left(\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}, k, \omega\right)\right)=1$, we can use this convergence factor to calculate Robin parameters for numerical computations, for which the frequencies $k$ and $\omega$ are bounded (by frequencies relevant to the global space-time domain and the ones supported by the numerical grid). Thus, we set

$$
\tilde{\rho}_{c}\left(\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}\right):=\max _{\frac{\pi}{L} \leq k \leq \frac{\pi}{h_{\Gamma}}, \frac{\pi}{T} \leq \omega \leq \frac{\pi}{\Delta t}} \varrho\left(\mathcal{R}\left(\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}, k, \omega\right)\right),
$$

where $L$ is a characteristic size of the computational domain and $h_{\Gamma}$ is a measure of the mesh step size on the interface (typically the mean-value of the segment lengths).

Let us consider the one-sided Robin case $\alpha:=\alpha_{1}=\alpha_{2}=\beta_{1}=\beta_{2}$, and set $\rho_{c}(\alpha):=\tilde{\rho}_{c}(\alpha, \alpha, \alpha, \alpha)$. Then, the continuous optimized Robin parameter $\alpha_{c}$ is defined as a solution of the following minimization problem :

$$
\rho_{c}\left(\alpha_{c}\right)=\min _{\alpha>0} \rho_{c}(\alpha) .
$$

9.2. Discrete-time optimized parameters. One can also consider the semi-discrete in time counterpart of the continuous convergence factor to better capture the discretetime frequencies, i.e. replace in the expression of $\mathcal{R}$ the term $i \omega$ by its discrete counterpart using the implicit Euler scheme, that is we replace $i \omega$ by $\frac{1-e^{-i \omega \Delta t}}{\Delta t}$. Equivalently, we replace in the expression of $\mathcal{R}$ (in (8.7)) the term $\omega$ by $\bar{\omega}:=-i\left(\frac{1-e^{-i \omega \Delta t}}{\Delta t}\right)$, and set $\mathcal{R}_{\Delta t}\left(\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}, k, \omega\right):=\mathcal{R}\left(\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}, k, \bar{\omega}\right)$.

Then, as above, we define

$$
\tilde{\rho}\left(\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}\right):=\max _{\frac{\pi}{L} \leq k \leq \frac{\pi}{h}, \frac{\pi}{T} \leq \omega \leq \frac{\pi}{\Delta t}} \varrho\left(\mathcal{R}_{\Delta t}\left(\alpha_{1}, \alpha_{2}, \beta_{1}, \beta_{2}, k, \omega\right)\right) .
$$

Let us consider the one-sided Robin case $\alpha:=\alpha_{1}=\alpha_{2}=\beta_{1}=\beta_{2}$, and define $\rho(\alpha):=\tilde{\rho}(\alpha, \alpha, \alpha, \alpha)$. Then, the Discrete-time $(D T)$ optimized Robin parameter $\alpha^{*}$ is defined as a solution of the following minimization problem :

$$
\rho\left(\alpha^{*}\right)=\min _{\alpha>0} \rho(\alpha)
$$

Remark 9.1. On could also consider optimized Robin-2p parameters $(\alpha, \beta)$ with $\alpha:=\alpha_{1}=\alpha_{2}, \beta:=\beta_{1}=\beta_{2}$, or 2 -sided parameters $(\gamma, \delta)$ with $\gamma:=\alpha_{1}=\beta_{1}$, $\delta:=\alpha_{2}=\beta_{2}$, that optimize the continuous or discrete-time convergence factors as done in [9]. Given their additional complexity, these more general cases will not be considered here, and are the subject of a subsequent article.
10. Numerical results. In this section, we present numerical experiments that illustrate the performances of the OSWR method of Section 4, with Freefem ++ [27]. For the space discretization we use the nonconforming Crouzeix-Raviart Finite Element method in 2D (i.e. piecewise linear elements continuous only at the midpoints of the edges of the mesh for the velocity $\mathbf{u}=\left(u_{x}, u_{y}\right)$, and piecewise constant $\mathbb{P}_{0}$ elements for the pressure $p$ ), and consider the backward Euler method for the time discretization.

In what follows, the term "monodomain solution" will refer to the fully discrete solution obtained on the global mesh without domain decomposition.

We set $\Omega=] 0,1[\times] 0,1[, T=1$, and consider the Stokes problem (2.1), where the value of the diffusion coefficient $\nu$ will be specified in each of the examples below.

From Remark 9.1, only one-sided Robin parameter $\alpha:=\alpha_{1}=\alpha_{2}=\beta_{1}=\beta_{2}$ will be considered. In particular, we will use the theoretical optimized values $\alpha_{c}$ and $\alpha^{*}$ defined in Section 9, which are calculated using the function fminsearch of MATLAB [37]. Random initial Robin data on the space-time interfaces will be used, unless specified.

In Section 10.1 some results are shown on the convergence of the OSWR algorithm, without and with modification of the pressure as in Section 7. In Section 10.2 we illustrate the influence of the Robin parameter on the convergence of the algorithm, and then in Section 10.3 we present results on a more realistic test case.
10.1. Recovering the pressure: a rotating velocity example. The diffusion coefficient is $\nu=1$ and we choose the right-hand side $\mathbf{f}$ and the values of the boundary and initial conditions so that the exact solution is given by

$$
\begin{aligned}
& \mathbf{u}(\mathbf{x}, t)=(-\cos (\pi y) \sin (\pi x) \cos (2 \pi t), \sin (\pi y) \cos (\pi x) \cos (2 \pi t)) \\
& p(\mathbf{x}, t)=\cos (t)\left(x^{2}-y^{2}\right), \quad \forall \mathbf{x} \in \Omega, \quad \forall t \in(0, T)
\end{aligned}
$$

On Figure 1 we show the velocity field $\mathbf{u}$ (on the left), and the pressure $p$ (on the right) at final time $t=1$.


Fig. 1. Example 1: rotating velocity field (left), and pressure (right)
The domain $\Omega$ is decomposed into nine subdomains as in Figure 2, and two meshes will be considered (as shown on Figure 2), with mesh sizes $h=0.0625$ and $h=0.0312$ respectively. To each mesh, the associated time step is $\Delta t=h$.


Fig. 2. Example 1: mesh 1 (left) and mesh 2 (right)

## f






FIG. 3. Example 1: relative errors (for $u_{x}, u_{y}$ and $p$ ) versus iterations with non-modified pressure (top), and modified pressure (bottom), for mesh 1 (left) and mesh 2 (right)

Remark 10.1. Even if we calculate a modified pressure at each iteration, we do not use it in the transmission conditions of Algorithm 4.1, thus this does not change the velocity convergence, as shown on Figure 3.

Remark 10.2. Here and in what follows, the pressure is modified at each iteration to illustrate the convergence of the multidomain solution to the monodomain one. A consequence of Remark 7.4 is that in practice one needs only to modify the pressure at the last OSWR iteration, which makes the cost of the modification negligible.
10.2. Optimized Robin parameters. The domain $\Omega$ is decomposed into two subdomains as in Figure 4, and we consider the three uniform meshes of Figure 4, with
mesh sizes on the interface and associated time steps equal to $h_{\Gamma}=\Delta t=1 / 12$, $h_{\Gamma}=\Delta t=1 / 24$, and $h_{\Gamma}=\Delta t=1 / 48$, respectively. In order to analyze the convergence behavior of the method, we simulate the error equations (i.e. we take homogeneous initial and boundary conditions, and $\mathbf{f}=0$ ). Thus, the OSWR solution converges to zero.


Fig. 4. Example 2: mesh 1 (left), mesh 2 (middle), and mesh 3 (right)
10.2.1. Case with a fixed mesh and different values of $\nu$. We consider mesh 2 (i.e. $h_{\Gamma}=\Delta t=1 / 24$ ). In Figure 5, we plot the evolution of the continuous convergence factor $\rho_{c}$ (on the left) and of the discrete-time convergence factor $\rho$ (on the right), as functions of the Robin parameter $\alpha$, for different values of $\nu: \nu=1$ (solid line), $\nu=0.5$ (dashed line), $\nu=0.1$ (dash-dotted line), $\nu=0.05$ (dotted line). The theoretical optimized values $\alpha_{c}$ (blue circle) and $\alpha^{*}$ (red star), are also shown. We observe that both $\alpha_{c}$ and $\alpha^{*}$ increase when $\nu$ decreases. However, the values of $\alpha_{c}$ and $\alpha^{*}$ are very different, and when $\nu$ decreases, $\alpha^{*}$ increases faster than $\alpha_{c}$, with an associated $\rho\left(\alpha^{*}\right)$ that increases slower than $\rho_{c}\left(\alpha_{c}\right)$.


Fig. 5. Example 2: continuous (left) and discrete-time (right) convergence factors versus $\alpha$, with $\alpha_{c}$ (blue circle) and $\alpha^{*}$ (red star), with $h_{\Gamma}=\Delta t=1 / 24$; for $\nu=1$ (solid line), $\nu=0.5$ (dashed line), $\nu=0.1$ (dash-dotted line), $\nu=0.05$ (dotted line)

In Figure 6, we plot the evolution of the relative errors, of $p, u_{x}$ and $u_{y}$, in the $L^{\infty}\left(0, T ; L^{2}(\Omega)\right)$-norm, in logarithmic scale, after twenty OSWR iterations, as functions of the Robin parameter $\alpha$. We also show the values of the errors obtained with optimized parameter $\alpha=\alpha_{c}$ (blue circle) and DT-optimized parameter $\alpha=\alpha^{*}$


Fig. 6. Example 2: Relative errors after 20 iterations (for $u_{x}, u_{y}$ and $p$ ) versus $\alpha$, with their values at $\alpha_{c}$ (blue circles) and at $\alpha^{*}$ (red stars), with $h_{\Gamma}=\Delta t=1 / 24$; for $\nu=1$ (top left), $\nu=0.5$ (top right), $\nu=0.1$ (bottom left), $\nu=0.05$ (bottom right)
10.2.2. Case with $\nu$ fixed and different space-time meshes. Let us take $\nu=0.1$. In Figure 7, we plot the evolution of the continuous (left) and discrete-time (right) convergence factors, versus $\alpha$, for different space-time meshes with $h_{\Gamma}=\Delta t=1 / 12$ (solid line), $h_{\Gamma}=\Delta t=1 / 24$ (dashed line), and $h_{\Gamma}=\Delta t=1 / 48$ (dash-dotted line). The theoretical optimized values $\alpha_{c}$ (blue circle) and $\alpha^{*}$ (red star) are also shown. We observe that both $\alpha_{c}$ and $\alpha^{*}$ decrease when the space-time mesh is refined. However, the values of $\alpha_{c}$ and $\alpha^{*}$ are again very different.

In Figure 8, we plot the relative errors, of $p, u_{x}$ and $u_{y}$, in the $L^{\infty}\left(0, T ; L^{2}(\Omega)\right)$ norm, after twenty OSWR iterations, versus Robin parameter $\alpha$, for mesh 1 (top left), mesh 2 (top right), and mesh 3 (bottom). We also show the values of the errors obtained with $\alpha=\alpha_{c}$ (blue circle) and $\alpha=\alpha^{*}$ (red star). We observe that $\alpha^{*}$ is close to the numerial Robin value giving the smallest error after the same number of iterations, while $\alpha_{c}$ gives a larger error, for all space-time meshes considered.
10.3. A more realistic test case. In this example we take $\nu=\frac{1}{\mathcal{R} e}$ with $\mathcal{R} e=200$, and $T=5$. The mesh is given on Figure 9, with 22232 mesh elements. The domain is decomposed into two subdomains, with the interface at $y=-0.9$, see Figure 9, where


Fig. 7. Example 2: continuous (left) and discrete-time (right) convergence factors versus $\alpha$, with $\alpha_{c}$ (blue circle) and $\alpha^{*}$ (red star), with $\nu=0.1$; for $h_{\Gamma}=\Delta t=1 / 12$ (solid line), $h_{\Gamma}=\Delta t=$ $1 / 24$ (dashed line), $h_{\Gamma}=\Delta t=1 / 48$ (dash-dotted line)


Fig. 8. Example 2: Relative errors after 20 iterations (for $u_{x}, u_{y}$ and $p$ ) versus $\alpha$, with their values at $\alpha_{c}$ (blue circles) and at $\alpha^{*}$ (red stars), with $\nu=0.1$; for $h_{\Gamma}=\Delta t=1 / 12$ (top left), $h_{\Gamma}=\Delta t=1 / 24$ (top right), $h_{\Gamma}=\Delta t=1 / 48$ (bottom)
domain 1 corresponds to the green and yellow parts, and domain 2 to the black part. The time step is $\Delta t=0.05$.

We set $\Omega_{f}=[-2.625,1.625] \times[-0.9,-0.6]$, represented by the yellow part of the mesh on Figure 9, and which corresponds to the location where the source term $\mathbf{f}$ in the Stokes equations does not vanish. Two different values for this source term will
be used in the numerical tests that follow.


Fig. 9. Example 3: mesh and domain decomposition

In Figure 10, we plot the evolution of the continuous convergence factor $\rho_{c}$ (left) and discrete-time convergence factor $\rho$ (right), as functions of the Robin parameter $\alpha$. The theoretical optimized values $\alpha_{c}$ (blue circle) and $\alpha^{*}$ (red star) are also shown, and their numerical values are $\alpha_{c} \approx 3.2283 \times 10^{-2}$ and $\alpha^{*} \approx 6.6063 \times 10^{-1}$, and differ from about a factor 20.


Fig. 10. Example 3: continuous (left) and discrete-time (right) convergence factors versus $\alpha$, with corresponding theoretical optimized values $\alpha_{c}$ (blue circle) and $\alpha^{*}$ (red star)

In this example we consider two different source terms in $\Omega_{f} \times(0, T)$ : a constant one: $\mathbf{f}=-2$, and then a variable in time one: $\mathbf{f}=-2(\sin (\pi t)+\cos (4 \pi t))$.

In Figures 11 and 12, we plot the pressure $p$ and the velocity field $\left(u_{x}, u_{y}\right)$ respectively, at times $t=1$ and $t=T=5$ (with a fixed color bar for $p$ ), for the case $\mathbf{f}$ constant. We observe that the stationary state is not reached yet.

In Figure 13, we show the evolution of the relative errors, between the OSWR and monodomain solutions, of $u_{x}, u_{y}$, and $p$, in the $L^{\infty}\left(0, T ; L^{2}(\Omega)\right)$-norm, as functions of OSWR iterations, for $\alpha=\alpha_{c}$ (cyan, green and blue curves) and $\alpha=\alpha^{*}$ (magenta, red, and black curves), with zero initial Robin data, with $\mathbf{f}$ constant (left), and $\mathbf{f}$ variable (right). For $\alpha=\alpha^{*}$, the curves of $u_{x}$ and $p$ are quite close, with a faster convergence for $u_{y}$. For $\alpha=\alpha_{c}$, the curves of $u_{x}$ and $u_{y}$ have almost the same speed of convergence, with a slower (resp. faster) convergence for $p$ for the first iterations, for $\mathbf{f}$ constant (resp. variable). Moreover, the convergence is much slower with $\alpha=\alpha_{c}$ than


Fig. 11. Example 3 (f constant): Pressure at $t=1$ (left) and at final time $t=5$ (right)


Fig. 12. Example 3 (f constant): Velocity field at $t=1$ (left) and at final time $t=5$ (right)
with $\alpha=\alpha^{*}$. This illustrates the importance of the effect of the numerical scheme used in the time direction.

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Fig. 13. Example 3 : Relative errors of $p, u_{x}, u_{y}$, versus iterations, with optimized Robin parameters $\alpha_{c}$ (cyan, green and blue curves) and $\alpha^{*}$ (magenta, red, and black curves), with $\mathbf{f}$ constant (left), and $\mathbf{f}$ variable (right).
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