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# METHODS OF REFLECTIONS: RELATIONS WITH SCHWARZ METHODS AND CLASSICAL STATIONARY ITERATIONS, SCALABILITY AND PRECONDITIONING. 

GABRIELE CIARAMELLA, MARTIN J. GANDER, LAURENCE HALPERN, AND JULIEN SALOMON


#### Abstract

The basic idea of the method of reflections appeared almost two hundred years ago; it is a method of successive approximations for the interaction of particles within a fluid, and it seems intuitively related to the Schwarz domain decomposition methods, the subdomains being the complements of the particle domains. We show in this paper that indeed there is a direct correspondence between the methods of reflections and Schwarz methods in the two particle/subdomain case. This allows us to give a new convergence analysis based on maximum principle techniques with precise convergence estimates that one could not obtain otherwise. We then show however also that in the case of more than two particles/subdomains, the methods of reflections and the Schwarz methods are really different methods, with different convergence properties. Using substructuring techniques from domain decomposition, we then show that the methods of reflections are classical block Jacobi and block Gauss-Seidel methods for the interface traces, and we derive new, relaxed versions of the methods of reflections with better convergence properties. We finally also introduce for the first time coarse corrections for the methods of reflections to make them scalable in the case when the number of particles becomes large. The substructured formulations permit the easy use of the methods of reflections as preconditioners for Krylov methods, and we illustrate scalability and preconditioning properties with numerical experiments.


## 1. Introduction

We start by briefly tracing the invention of the method of reflections and its historical development, for a more detailed treatment, see [10]. The fundamental idea for the method of reflections can already be found in the book of Murphy [39, page 93] from 1833 under the name "principle of successive influences", which clearly indicates that the method of reflections is a method of successive approximations. Lamb used in 1906 a similar approach for Laplace's equation [33, page 122]. A related method in the work of Lorentz [37, page 29] was later called "Spiegelungsmethode" (method of reflections) in [25, page 928]. The method of reflections itself was then presented in concrete mathematical notation by Smoluchowski in 1911 with the goal to understand how the motion of a sphere in a viscous fluid is influenced by the presence or motion of another sphere [42]. In the case of the Stokes equation, the analytical solution for one sphere is already a series expansion, and thus Smoluchowski assumes that the radii of the two spheres are small

[^0]compared to their distance, and then uses the method of reflections to compute a series expansion of the coupled solution up to some order in the inverse distance of the spheres. The method is thus a direct method to obtain a series solution in the inverse of the distance, up to some order. Smoluchowski then generalizes the method of reflections to the case of more than two spheres, leading to similar series approximations. In 1934, Golusin introduced a parallel method of reflections for Laplace's equation for $J$ objects [22, 21], and derived a condition for its convergence, which indicates that the parallel method of reflections in the case of more than two objects converges only under certain additional conditions. Golusin conjectured that these conditions depend on the distance between the objects and their radius. In 1942, Burgers [3, 4] investigated the influence of the concentration of spherical particles on the sedimentation velocity for the Stokes equation, mentioning the work of Smoluchowski, but without describing precisely an algorithm, and using to a large extend physical intuition. In 1959, Kynch presented for the Stokes equation the idea of simply summing two solutions corresponding to two particles alone in [32, page 197], under the assumption that the distance between their centers is again large. This could be interpreted as a parallel method of reflections, where the separate contributions are also summed, but again, no general algorithm is given. Happel and Brenner explained in 1983 a different parallel version of the method of reflections which alternates between one fixed object and the group of all the others treated in parallel, see [24], with the goal to increase the order of approximation of the expansion of the solution in a neighborhood of a given object. Their method has to be applied (independently) for each particle.

Luke gave then in 1989 a first convergence analysis for the alternating method of reflections applied to the Stokes equation [38], using a variational characterization of the method based on projections, similar in style of one of the classical convergence analyses of the Schwarz method given by Lions in the first of his three seminal papers [35]. Kim and Karrila dedicated an entire chapter in their book from 1991 to the parallel method of reflections for the Stokes problem [30]. The method is first motivated like in [32] by just summing two one-particle solutions, and only the first terms in the series expansion are obtained. Dhont also dedicated a special section to the alternating method of reflections for the Stokes equation in 1996 [12, Section 5.12]. The case of two objects is first treated, and then an extension to three objects is given, where Dhont goes cyclically through the three object in the algorithm. Balabane proved in 1997 convergence of the alternating method of reflections for the Helmholtz equation in unbounded domains in [2], and generalized his results to the parallel method of reflections in [1]. These convergence results are valid however only in low frequency regimes. In 2001, Ichiki and Brady presented the parallel method of reflections [27] for Stokes type problems. They started with the two particle case, and then gave a generalization just summing all contributions that were computed in parallel. They presented this iterative approach also in matrix form, relating it to a stationary iteration based on a matrix splitting. By numerical experiments, they showed that the method does not converge for three particles, if the separation distance of the particles is not large enough. They thus proposed to use the method as a preconditioner for a Krylov method. Traytak posed in 2006 in a short note directly the parallel method of reflections for $N$ objects, written in PDE form for Laplace's equation [44, Section 2], and then used a theorem proved by Golusin [21] to derive sufficient conditions for the convergence based on
the distances between the objects. More recently, Höfer and Velázquez used the parallel method of reflections also as an analytic tool to prove homogenization results [26] (see also [28]), and they modified the usual parallel method by adding different weighting coefficients. Since they were interested in the theoretical case of an infinite number of objects, they could not use an alternating method. Laurent, Legendre and Salomon studied the alternating and parallel methods of reflections in [34] for various types of boundary conditions, introducing also an averaged version of the parallel method. They proved convergence based on the alternating projection method in Hilbert spaces, see for example [40], and also using techniques like in $[2,1]$.

So there are two main variants of the method of reflections: the alternating one and the parallel one. There are also two different approaches to analyze the convergence of the method of reflections: first, people worked on direct estimates performed on the single/double layer formulation of the boundary value problems involved in the iterations, see $[21,44,2,1]$. There is however also the interpretation of the method as alternating projections in Hilbert spaces, see [38, 34]. In the case of orthogonal projections this interpretation leads to convergence estimates.

When we started studying methods of reflections more than three years ago, we thought immediately that the methods must be intimately related to the Schwarz domain decomposition methods. This intuition was confirmed when we studied the literature and found that analysis techniques based on projections were already used to study methods of reflections like for Schwarz methods, and such a possible relation was even mentioned in the literature ("This paper considers a reflection method in the spirit of Schwarz's alternating procedure" [38]). We will show here that the methods of reflections can be indeed identified with Schwarz domain decomposition methods in the case of two particles/subdomains, which leads to a new convergence proof with sharp convergence estimates using maximum principle techniques. In the case of many particles/subdomains however the methods of reflections are fundamentally different from Schwarz domain decomposition methods. The main difference is that in domain decomposition, the interface data is to be determined by the subdomain iteration, whereas in the method of reflections, the interface data is given on the boundary of the particles by the problem. Substructuring techniques from domain decomposition allow us however to reformulate the methods of reflections in the form of standard block Gauss-Seidel and block Jacobi iterations in the traces for integral formulations. This reformulation leads to a new understanding of the method of reflections also for many particles, and reveals shortcomings, especially in the parallel variant. It also allows us to develop new coarse space corrections, again using domain decomposition techniques, in order to obtain scalable methods of reflections when the number of particles grows.

Our paper is structured as follows: we start by presenting in Section 2 the class of Laplace problems on perforated domains we will use to understand the methods of reflections. We then present the different forms of the methods of reflections: in Section 3 the alternating method of reflections, in Section 4 the parallel method of reflections, in Section 5 some other variants of the method of reflections. We give in each case reformulations of the methods using domain decomposition techniques, which allows us to study their convergence properties and give new convergence proofs and convergence estimates. In Section 6, we then give a scalability analysis, and also propose for the first time for the methods of reflections a coarse space


Figure 1. Left: example of a domain $\Omega$ (a disk) with three objects $O_{j}$ (holes). Right: example of the functions $\varphi_{j}^{a}$ and $\varphi_{j}^{b}$ for $J=3$.
correction to make them scalable in the case when the number of particles becomes large. We show numerical experiments in Section 7, followed by a brief conclusion in Section 8.

## 2. The Laplace problem in perforated domains

We consider a Laplace problem where $J$ objects $O_{j}$ (holes) are present in a simply connected domain $\Omega \subset \mathbb{R}^{n}$; for an example, see Figure 1 on the left. We assume that $\Omega$ and $O_{j}$ are simply connected with sufficiently smooth boundaries $\partial \Omega$ and $\partial O_{j}$, e.g., of class $C^{1}$, and that

$$
\begin{equation*}
\bar{O}_{j} \cap \bar{O}_{k}=\emptyset \text { for any } j \neq k \tag{1}
\end{equation*}
$$

Our goal is finding the (weak) solution $u \in H_{0}^{1}(\Omega)$ to the problem

$$
\begin{align*}
\Delta u & =0 \text { in } \Omega \backslash \cup_{j} \overline{O_{j}} \\
u & =0 \text { on } \partial \Omega  \tag{2}\\
u & =g_{j} \text { on } \partial O_{j} \text { for } j=1, \ldots, J,
\end{align*}
$$

where we assume that the functions $g_{j}$ are bounded on $\partial O_{j}$ and in $H^{1 / 2}\left(\partial O_{j}\right)$ for $j=1, \ldots, J$. This problem is well-posed and uniquely solved by $u \in H_{0}^{1}$ ( $\Omega \backslash$ $\left.\cup_{j} \overline{O_{j}}\right) \cap C^{2}(\Omega)$; see, e.g., $[20,31]$. In order to introduce the method of reflections, we modify (2) by extending harmonically the solution $u$ into the objects $O_{j}$. We define $\widetilde{\Omega}:=\Omega \backslash \cup_{j=1}^{J} \partial O_{j}$ and then (2) becomes

$$
\begin{align*}
\Delta u & =0 \text { in } \widetilde{\Omega} \\
u & =0 \text { on } \partial \Omega  \tag{3}\\
u & =g_{j} \text { on } \partial O_{j} \text { for } j=1, \ldots, J .
\end{align*}
$$

The solution $u$ to (3) belongs to the Hilbert space

$$
\begin{equation*}
H:=\left\{v \in H_{0}^{1}(\Omega): \Delta v=0 \text { in } \widetilde{\Omega}\right\} \tag{4}
\end{equation*}
$$

endowed with the $H_{0}^{1}$ scalar product $\langle\cdot, \cdot\rangle_{H_{0}^{1}}:=\langle\nabla(\cdot), \nabla(\cdot)\rangle_{L^{2}}$, and by Weyl's theorem [20], we have that $H \subset C^{\infty}(\widetilde{\Omega})$.

We introduce the subspaces $H_{j}$ of $H$ defined by

$$
\begin{equation*}
H_{j}:=\left\{v \in H: \Delta v=0 \text { in } \Omega \backslash \partial O_{j}\right\} . \tag{5}
\end{equation*}
$$

If $\Omega$ is unbounded, one needs to consider the completion of $H_{0}^{1}$ for the semi-norm induced by $\langle\cdot, \cdot\rangle_{H_{0}^{1}}$, for instance weighted Sobolev spaces, see [11, 23], but similar results could be obtained.

It has been proved in [34] that $H=\sum_{j=1}^{J} H_{j}$. For our purposes, we need to refine this result. We have the following lemma.
Lemma 1. Under the Assumption (1), we have that $H=\bigoplus_{j=1}^{J} H_{j}$, which means that the subspaces $H_{j}$ are linearly independent, that is $\left(H_{1}+\cdots+H_{j-1}\right) \cap H_{j}=\{0\}$ for any $j \geq 2$.

Proof. Consider any function $v$ in $H_{j} \cap H_{k}$ for any $j \neq k$. Then $v$ is harmonic everywhere and zero on $\partial \Omega$, thus $v=0$.

This and the fact that $H=\sum_{j=1}^{J} H_{j}$ proved in [34] imply the claim.
To explain Lemma 1 we consider a one-dimensional example, and then show that if Assumption (1) is not satisfied, then the result does not hold in general. The domain $\Omega$ is the interval $(0,1)$ and the objects are subintervals $O_{j}=\left(a_{j}, b_{j}\right)$ for $1 \leq j \leq J$ such that $0<a_{1}<b_{1}<a_{2}<\cdots<a_{J}<b_{J}<1$. In this case, the space $H_{j}$ is spanned by the hat-functions

$$
\varphi_{j}^{a}(x):=\left\{\begin{array}{ll}
\frac{x}{a_{j}} & \text { if } x \in\left[0, a_{j}\right], \\
\frac{b_{j}-x}{b_{j}-a_{j}} & \text { if } x \in\left[a_{j}, b_{j}\right], \\
0 & \text { if } x \in\left[b_{j}, 1\right],
\end{array} \quad \varphi_{j}^{b}(x):= \begin{cases}0 & \text { if } x \in\left[0, a_{j}\right] \\
\frac{x-a_{j}}{b_{j}-a_{j}} & \text { if } x \in\left[a_{j}, b_{j}\right] \\
\frac{1-x}{1-b_{j}} & \text { if } x \in\left[b_{j}, 1\right]\end{cases}\right.
$$

that have value 1 at $a_{j}$ and $b_{j}$ and are supported in $\left[0, b_{j}\right]$ and $\left[a_{j}, 1\right]$, for an example, see Figure 1 (right). Notice that $\operatorname{dim} H_{j}=2$. Therefore, since all the points $a_{j}$ and $b_{j}$ are distinct, we have that all the functions $\varphi_{j}^{a}$ and $\varphi_{j}^{b}$ are linearly independent, $H=\operatorname{span}\left\{\varphi_{1}^{a}, \varphi_{1}^{b}, \ldots, \varphi_{J}^{a}, \varphi_{J}^{b}\right\}$, and $\operatorname{dim} H=2 J$. The hypothesis $a_{j}<b_{j}<a_{j+1}<$ $b_{j+1}$ then clearly implies that $\left(H_{1}+\cdots+H_{k-1}\right) \cap H_{k}=\{0\}$ for $2 \leq k \leq J$, which is the result of Lemma 1. On the other hand, if we assume that two objects are not distinct, that is, e.g., $b_{j}=a_{j+1}$, then one can verify that $c_{0} \varphi_{j+1}^{a}+c_{1} \varphi_{j+1}^{b}+c_{2} \varphi_{j}^{a}+$ $c_{3} \varphi_{j}^{b}=0$, with $c_{0}=1, c_{1}=\frac{1-b_{j+1}}{1-b_{j}}, c_{2}=-\frac{a_{j}}{a_{j+1}}$, and $c_{3}=-1$. Hence the functions $\varphi_{j}^{a}, \varphi_{j}^{b}, \varphi_{j+1}^{a}$, and $\varphi_{j+1}^{b}$ are not linearly independent and $\operatorname{dim} H=2 J-1<2 J=$ $\sum_{j=1}^{J} \operatorname{dim} H_{j}$. Hence $H$ cannot be written as the direct sum of the subspaces $H_{j}$.

Iterative methods suitable for solving problem (2)-(3) sequentially or in parallel are the methods of reflections (MR), which we will present and study in the following sections.

## 3. The alternating method of reflections

In this section, we study the alternating method of reflections (AltMR). We first recall in Section 3.1 the classical AltMR formulation present in the literature. In Section 3.2, we introduce a new, equivalent (volume) formulation of the AltMR. This is then used in Section 3.3 where we present the AltMR in substructured form, and study its relations with the classical block Gauss-Seidel iterative method. Finally, in Section 3.4 we study the relations between the AltMR and the alternating Schwarz method.
3.1. The classical AltMR formulations. The alternating method of reflections was invented for two objects by Smoluchowski in 1911 [42], and then extended to the general case of $J$ objects by Luke [38]. To understand this method, we begin by quoting a sentence from Luke to apply it to our problem: "The strategy of the method of reflections is to repeatedly correct the boundary values on the various
objects. The reflection procedure is formalized as follows": it starts with $u^{0} \in H$, that is

$$
\begin{equation*}
\Delta u^{0}=0 \text { in } \Omega \backslash \cup_{j} \partial O_{j}, u^{0}=0 \text { on } \partial \Omega \tag{6}
\end{equation*}
$$

Note that $u^{0}$ does not satisfy the boundary conditions of (2) on the holes, which would require

$$
u^{0}=g_{j} \text { on } \partial O_{j} \text { for } j=1, \ldots, J
$$

The sequence of approximate solutions $\left\{u^{k}\right\}_{k \in \mathbb{N}}$ is defined as follows: given $u^{0}$, one computes for iteration index $k=1,2, \ldots$ from $u^{k-1}$

$$
\begin{aligned}
u^{k-1+\frac{1}{J}} & =u^{k-1}+d_{1}^{k} \\
u^{k-1+\frac{2}{J}} & =u^{k-1}+d_{1}^{k}+d_{2}^{k} \\
u^{k-1+\frac{3}{J}} & =u^{k-1}+d_{1}^{k}+d_{2}^{k}+d_{3}^{k}
\end{aligned}
$$

and then obtains the new approximation

$$
\begin{equation*}
u^{k}=u^{k-1}+\sum_{j=1}^{J} d_{j}^{k} \tag{7}
\end{equation*}
$$

where $d_{j}^{k} \in H_{j}$ is computed in such a way that, when added to $u^{k-1+\frac{j-1}{J}}$, the value on the boundary $\partial O_{j}$ is corrected to its given boundary value $g_{j}$, that is

$$
u^{k-1+\frac{j}{J}}=u^{k-1+\frac{j-1}{J}}+d_{j}^{k}=g_{j} \text { on } \partial O_{j} .
$$

This means that $d_{j}^{k}$, for $j=1, \ldots, J$, must be the solution to

$$
\begin{aligned}
\Delta d_{j}^{k} & =0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{k}=0 \text { on } \partial \Omega \\
d_{j}^{k} & =g_{j}-u^{k-1+\frac{j-1}{J}}=g_{j}-u^{k-1}-\sum_{\ell=1}^{j-1} d_{\ell}^{k} \text { on } \partial O_{j} .
\end{aligned}
$$

A simple and intuitive explanation of the iteration described in (7)-(8) can be found in [10].

In $[34,2,10]$ an equivalent form of (8) is presented, which we derive now for completeness: we consider the boundary condition of (8) on $\partial O_{j}$, for $j=1, \ldots, J$, and manipulate it as follows:

$$
\begin{aligned}
d_{j}^{k} & =g_{j}-u^{k-1}-\sum_{\ell=1}^{j-1} d_{\ell}^{k}=g_{j}-u^{k-2}-\sum_{\ell=1}^{J} d_{\ell}^{k-1}-\sum_{\ell=1}^{j-1} d_{\ell}^{k} \\
& =g_{j}-u^{k-2}-\sum_{\ell=1}^{j} d_{\ell}^{k-1}-\sum_{\ell=j+1}^{J} d_{\ell}^{k-1}-\sum_{\ell=1}^{j-1} d_{\ell}^{k}
\end{aligned}
$$

Now, we notice that $g_{j}-u^{k-2}-\sum_{\ell=1}^{j} d_{\ell}^{k-1}=0$, since $d_{j}^{k-1}$ solves at iteration $k-1$ problem (8), that is

$$
\begin{aligned}
\Delta d_{j}^{k-1} & =0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{k-1}=0 \text { on } \partial \Omega \\
d_{j}^{k-1} & =g_{j}-u^{k-2}-\sum_{\ell=1}^{j-1} d_{\ell}^{k-1} \text { on } \partial O_{j}
\end{aligned}
$$

Therefore, we have obtained that $d_{j}^{k}$ can be expressed on $\partial O_{j}$ as combinations of other differences only,

$$
\begin{equation*}
d_{j}^{k}=-\sum_{\ell=j+1}^{J} d_{\ell}^{k-1}-\sum_{\ell=1}^{j-1} d_{\ell}^{k} \tag{9}
\end{equation*}
$$

and the explicit dependence on $g_{j}$ and $u^{k-1}$ disappeared. Hence, iteration (8) becomes

$$
\begin{align*}
\Delta d_{j}^{k} & =0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{k}=0 \text { on } \partial \Omega \\
d_{j}^{k} & =-\sum_{\ell=1}^{j-1} d_{\ell}^{k}-\sum_{\ell=j+1}^{J} d_{\ell}^{k-1} \text { on } \partial O_{j} \tag{10}
\end{align*}
$$

for $j=1, \ldots, J$, that is the form of the sequential method of reflections presented in $[34,2]$. Now the sequences $\left\{d_{j}^{k}\right\}_{k \in \mathbb{N}^{+}}$have to be initialized for all $j$. To this end, it is sufficient, for example, to consider (8) for $k=1$ :

$$
\begin{align*}
\Delta d_{j}^{1} & =0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{1}=0 \text { on } \partial \Omega \\
d_{j}^{1} & =g_{j}-\sum_{\ell=1}^{j-1} d_{\ell}^{1}-u^{0} \text { on } \partial O_{j} \tag{11}
\end{align*}
$$

for $j=1, \ldots, J$.
Notice that the AltMR iteration in the form (8) cannot be initialized with an arbitrary function: the initial guess $u^{0}$ must be in $H$ (recall (6)). A choice of a function $u^{0} \notin H$ will produce a sequence $\left\{u^{k}\right\}_{k}$ that does not necessarily converge to the solution $u$ to (2). In fact, (8) produces corrections $d_{j}^{k}$ which are harmonic in $\Omega \backslash \cup_{j=1}^{J} \partial O_{j}$. Hence, if $u^{0}$ can be decomposed into the sum of harmonic and non-harmonic components, namely $u^{0}=u_{\text {harm }}^{0}+u_{\text {non-harm }}^{0}$, then only $u_{\text {harm }}^{0}$ is corrected by the MR procedure. The AltMR as a program is given in Algorithm 1.

```
Algorithm 1 Alternating Method of Reflections (AltMR)
Input: \(K\) (maximum number of iterations), tol (tolerance).
    Set \(u^{0} \in H\) and \(k=1\).
    for \(j=1: J\) do
        Initialize \(d_{j}^{1}\) solving problem (11).
    end for
    Compute \(u^{1}=u^{0}+\sum_{j=1}^{J} d_{j}^{1}\).
    while \(k<K\) and \(\left\|u^{k}-u^{k-1}\right\|>\) tol do
        Update \(k=k+1\).
        for \(j=1: J\) do
            Compute \(d_{j}^{k}\) solving problem (10).
        end for
        Compute the approximation \(u^{k}=u^{k-1}+\sum_{j=1}^{J} d_{j}^{k}\).
    end while
```

3.2. A new formulation of the AltMR. In this section, we formulate the AltMR in a third equivalent way. To do so, we recall the $d_{j}^{k}$ that solve (8), define $J$ sequences $\left\{v_{j}^{k}\right\}_{k \in \mathbb{N}}$ for any $j \in\{1, \ldots, J\}$ as

$$
\begin{equation*}
v_{j}^{k}:=v_{j}^{k-1}+d_{j}^{k} \text { with } v_{j}^{0} \in H_{j}, \tag{12}
\end{equation*}
$$

and $u_{0}:=\sum_{j=1}^{J} v_{j}^{0}$ (recall Lemma 1). Notice that since $d_{j}^{k} \in H_{j}, v_{j}^{k} \in H_{j}$ as well. Now, we notice that (7) implies that $u^{k}=u_{0}+\sum_{j=1}^{J} \sum_{n=0}^{k} d_{j}^{n}$, which allows us to compute

$$
u^{k}=u_{0}+\sum_{j=1}^{J} \sum_{n=1}^{k} d_{j}^{n}=u_{0}+\sum_{j=1}^{J} \sum_{n=0}^{k}\left(v_{j}^{n}-v_{j}^{n-1}\right)=u_{0}+\sum_{j=1}^{J}\left(v_{j}^{k}-v_{j}^{0}\right)
$$

and recalling that $\sum_{j=1}^{J} v_{j}^{0}=u_{0}$ we obtain

$$
\begin{equation*}
u^{k}=\sum_{j=1}^{J} v_{j}^{k} \tag{13}
\end{equation*}
$$

Notice that equation (13) allows us to express the approximate solution at the $k$ th iteration as a decomposition (unique according to Lemma 1) into $J$ components, each of them being associated with an object and belonging to one of the spaces $H_{j}$. Starting from (12), using (13) and the boundary condition in (8), that is

$$
d_{j}^{k}=g_{j}-u^{k-1}-\sum_{\ell=1}^{j-1} d_{\ell}^{k} \text { on } \partial O_{j}
$$

we have on $\partial O_{j}$ that

$$
\begin{aligned}
v_{j}^{k} & =v_{j}^{k-1}+d_{j}^{k}=v_{j}^{k-1}+g_{j}-u^{k-1}-\sum_{\ell=1}^{j-1}\left(v_{\ell}^{k}-v_{\ell}^{k-1}\right) \\
& =g_{j}+v_{j}^{k-1}-\sum_{\ell=1}^{J} v_{\ell}^{k-1}-\sum_{\ell=1}^{j-1} v_{\ell}^{k}+\sum_{\ell=1}^{j-1} v_{\ell}^{k-1} \\
& =g_{j}-\sum_{\ell=j+1}^{J} v_{\ell}^{k-1}-\sum_{\ell=1}^{j-1} v_{\ell}^{k}
\end{aligned}
$$

Recalling that $v_{j}^{k} \in H_{j}$, we obtain that $v_{j}^{k}$ is solution to

$$
\begin{align*}
\Delta v_{j}^{k} & =0 \text { in } \Omega \backslash \partial O_{j}, v_{j}^{k}=0 \text { on } \partial \Omega \\
v_{j}^{k} & =g_{j}-\sum_{\ell=j+1}^{J} v_{\ell}^{k-1}-\sum_{\ell=1}^{j-1} v_{\ell}^{k} \text { on } \partial O_{j} \tag{14}
\end{align*}
$$

The equivalence between (14) and (10)-(11) is proved in the following theorem.
Theorem 1. Consider the sequences $\left\{u^{k}\right\}_{k \in \mathbb{N}} \subset H,\left\{d_{j}^{k}\right\}_{k \in \mathbb{N}^{+}}$with $d_{j}^{k} \in H_{j}$, and $\left\{v_{j}^{k}\right\}_{k \in \mathbb{N}}$ with $v_{j}^{k} \in H_{j}$ for $j=1, \ldots, J$, and assume that $v_{j}^{k}=v_{j}^{k-1}+d_{j}^{k}$ for $k \geq 1$ and $j=1, \ldots, J$. Then the following statements are equivalent:
(a) the $d_{j}^{k}$ solve (10)-(11) and $u^{k}=u^{k-1}+\sum_{j=1}^{J} d_{j}^{n}=\sum_{n=1}^{k} \sum_{j=1}^{J} d_{j}^{n}$.
(b) the $v_{j}^{k}$ solve (14) and $u^{k}=\sum_{j=1}^{J} v_{j}^{k}$.

Proof. The implication (a) $\Rightarrow(\mathrm{b})$ is proved in the discussion before the theorem. Hence, we need to show that (b) $\Rightarrow$ (a). First, using $v_{j}^{k}=v_{j}^{k-1}+d_{j}^{k}$ we get

$$
\begin{equation*}
u^{k}=\sum_{j=1}^{J} v_{j}^{k}=u^{0}+\sum_{n=1}^{k} \sum_{j=1}^{J}\left(v_{j}^{n}-v_{j}^{n-1}\right)=u^{0}+\sum_{n=1}^{k} \sum_{j=1}^{J} d_{j}^{n}=u^{k-1}+\sum_{j=1}^{J} d_{j}^{k} \tag{15}
\end{equation*}
$$

Notice that $d_{j}^{k}=v_{j}^{k}-v_{j}^{k-1}$. Clearly, $d_{j}^{k} \in H_{j}$ because $v_{j}^{k}, v_{j}^{k-1} \in H_{j}$. Hence we have to prove that $d_{j}^{k}$ satisfies the transmission condition of (10). For $k=1$, we use that $d_{j}^{1}=v_{j}^{1}-v_{j}^{0}$ and the boundary condition in (14) for $k=1$ to write

$$
d_{j}^{1}=v_{j}^{1}-v_{j}^{0}=g_{j}-\sum_{\ell=1}^{j-1} v_{j}^{1}-\sum_{\ell=j}^{J} v_{j}^{0}=g_{j}-\sum_{\ell=1}^{j-1} d_{j}^{1}-u^{0}
$$

In the case $k>1$, we write the transmission condition on $\partial O_{j}$ of (14) for $k$ and $k-1$ :

$$
v_{j}^{k}=g_{j}-\sum_{\ell=1}^{j-1} v_{\ell}^{k}-\sum_{\ell=j+1}^{J} v_{\ell}^{k-1}, \quad v_{j}^{k-1}=g_{j}-\sum_{\ell=1}^{j-1} v_{\ell}^{k-1}-\sum_{\ell=j+1}^{J} v_{\ell}^{k-2}
$$

Subtracting term by term these equations, we get the transmission condition of (10).

Using (14), we now rewrite Algorithm 1 in terms of $v_{j}^{k}$ to obtain Algorithm 2.

```
Algorithm 2 Alternating Method of Reflections (AltMR)
Input: \(K\) (maximum number of iterations), tol (tolerance).
    Set \(v_{j}^{0} \in H_{j}\) for \(j=1, \ldots, J\), and \(k=1\).
    for \(j=1: J\) do
        Compute \(v_{j}^{1}\) solving problem (14).
    end for
    Compute the approximation \(u^{1}=\sum_{j=1}^{J} v_{j}^{1}\).
    while \(k<K\) and \(\left\|u^{k}-u^{k-1}\right\|>\) tol do
        Update \(k=k+1\).
        for \(j=1: J\) do
            Compute \(v_{j}^{k}\) solving problem (14).
        end for
        Compute the approximation \(u^{k}=\sum_{j=1}^{J} v_{j}^{k}\).
    end while
```

We next show that this new formulation of the AltMR in terms of the decomposition functions $v_{j}^{k}$ allows us to easily obtain a substructured formulation of the AltMR procedure.
3.3. Substructured AltMR as block Gauss-Seidel method. In this section, we first write the AltMR in substructured form and then show that it can be interpreted as a block Gauss-Seidel method. To do so, we recall Lemma 1, which plays a very important role in our interpretation of the method of reflections, since it ensures that the unique solution $u \in H$ to (3) can be uniquely decomposed as $v_{1}+v_{2}+\cdots+v_{J}$ with $v_{j} \in H_{j}$. For each of these $v_{j}$, let us denote by $\widetilde{g}_{j}$ the corresponding Dirichlet trace on $\partial O_{j}$ so that, because of Definition (5), we have

$$
\begin{equation*}
\Delta v_{j}=0 \quad \text { in } \Omega \backslash \partial O_{j}, v_{j}=0 \quad \text { on } \partial \Omega, v_{j}=\widetilde{g}_{j} \text { on } \partial O_{j} . \tag{16}
\end{equation*}
$$

Now, we introduce the operator $G_{j}$ that maps Dirichlet boundary data into functions that are harmonic in $\Omega \backslash \partial O_{j}$, and such that the solution to (16) can be written as $v_{j}=G_{j}\left(\widetilde{g}_{j}\right)$. The operator $G_{j}$ can be explicitly expressed in terms of Green's representation formulas, and, in particular, as the composition of a single-layer potential integral operator and Dirichlet-to-Neumann operators: for $\bar{g} \in H^{1 / 2}\left(\partial O_{j}\right)$, we have

$$
G_{j}(\bar{g})(\mathbf{x}):=\int_{\partial O_{j}} G(\mathbf{x}, \mathbf{y})\left(\operatorname{DtN}_{\mathrm{j}, \mathrm{e}}(\bar{g})+\operatorname{DtN}_{\mathrm{j}, \mathrm{i}}(\bar{g})\right)(\mathbf{y}) d s(\mathbf{y})
$$

where $G(\mathbf{x}, \mathbf{y})$ is the Green's function associated to the problem, $\mathrm{DtN}_{\mathrm{j}, \mathrm{i}}: H^{1 / 2}\left(\partial O_{j}\right) \rightarrow$ $H^{-1 / 2}\left(\partial O_{j}\right)$ and $\operatorname{DtN}_{\mathrm{j}, \mathrm{e}}: H^{1 / 2}\left(\partial O_{j}\right) \rightarrow H^{-1 / 2}\left(\partial O_{j}\right)$. These Dirichlet-to-Neumann operators are obtained by first solving the interior and exterior problems for $O_{j}$,

$$
\begin{array}{rlrl}
\Delta v_{e} & =0 \text { in } \Omega \backslash \overline{O_{j}}, & \Delta v_{i} & =0 \text { in } O_{j} \\
v_{e} & =\bar{g} \text { on } \partial O_{j}, & v_{i} & =\bar{g} \text { on } \partial O_{j} \\
v_{e} & =0 \text { on } \partial \Omega, &
\end{array}
$$

and then extracting the Neumann trace on $\partial O_{j}$ of their solutions, that is $\mathrm{DtN}_{\mathrm{j}, \mathrm{e}}(\bar{g}):=$ $\left.\partial_{n} v_{e}(\bar{g})\right|_{\partial O_{j}}$ and $\mathrm{DtN}_{\mathrm{j}, \mathrm{i}}(\bar{g}):=-\left.\partial_{n} v_{i}(\bar{g})\right|_{\partial O_{j}}$, with $\partial_{n}=n \cdot \nabla$ the unit outward normal derivative with respect to $\Omega \backslash \overline{O_{j}}$. In addition, we also need the trace operator $\tau_{j}: H^{1}(\Omega) \rightarrow H^{1 / 2}\left(\partial O_{j}\right)$ such that $\tau_{j} v$ is the Dirichlet trace of $v$ on $\partial O_{j}$, $\tau_{j} v:=\left.v\right|_{\partial O_{j}}$.

With this framework, we can rewrite problem (2) in integral form: we look for boundary data $\widetilde{g}_{j}$, for $j=1, \ldots, J$, such that the function $u:=\sum_{j=1}^{J} G_{j}\left(\widetilde{g}_{j}\right)$ satisfies

$$
\begin{equation*}
\tau_{j} u=\tau_{j}\left(\sum_{j=1}^{J} G_{j}\left(\widetilde{g}_{j}\right)\right)=g_{j}, \text { on } \partial O_{j} \text { for } j=1, \ldots, J \tag{17}
\end{equation*}
$$

Equation (17) is equivalent to the linear system

$$
\begin{equation*}
A \widetilde{\mathbf{g}}=\mathbf{g} \tag{18}
\end{equation*}
$$

where

$$
A:=\left[\begin{array}{ccccc}
I_{1} & \tau_{1} G_{2} & \tau_{1} G_{3} & \cdots & \tau_{1} G_{J}  \tag{19}\\
\tau_{2} G_{1} & I_{2} & \tau_{2} G_{3} & \cdots & \tau_{2} G_{J} \\
\tau_{3} G_{1} & \tau_{3} G_{2} & I_{3} & \cdots & \tau_{3} G_{J} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
\tau_{J} G_{1} & \tau_{J} G_{2} & \tau_{J} G_{3} & \cdots & I_{J}
\end{array}\right], \widetilde{\mathbf{g}}:=\left[\begin{array}{c}
\widetilde{g}_{1} \\
\widetilde{g}_{2} \\
\widetilde{g}_{3} \\
\vdots \\
\widetilde{g}_{J}
\end{array}\right], \quad \text { and } \mathbf{g}:=\left[\begin{array}{c}
g_{1} \\
g_{2} \\
g_{3} \\
\vdots \\
g_{J}
\end{array}\right],
$$

and $I_{j}$ are identity operators, and we used the fact that $\tau_{j} G_{j}=I_{j}$.
The following proposition shows that the linear system (18) is an equivalent formulation of the problem (3).

Proposition 1. The system (18) is equivalent to (3). Moreover, the solution $u$ to (3) is uniquely decomposed as $u=\sum_{j=1}^{J} G_{j}\left(\widetilde{g}_{j}\right)$, where $\widetilde{\mathbf{g}}:=\left[\widetilde{g}_{1} \cdots \widetilde{g}_{J}\right]^{\top}$ is the solution to (18).
Proof. The equivalence between (3) and (18) follows by Lemma 1 and the fact that each component $v_{j}$ of $u$ is uniquely determined by the Dirichlet boundary data $\widetilde{g}_{j}$ on $\partial O_{j}$. The second statement is obtained by Lemma 1 and the following
argument: from (17) we have that $\tau_{j} u=g_{j}$ on $\partial O_{j}$ for $j=1, \ldots, J$, and we have that $\Delta u=\Delta \sum_{j=1}^{J} G_{j}\left(\widetilde{g}_{j}\right)=0$ on $\Omega \backslash \cup_{j} \partial O_{j}$.

Remark 1. The variables $\widetilde{g}_{j}$ are not necessarily of Dirichlet type: one could formulate the problem with, e.g., Neumann data. In this case the definition of $G_{j}$ would not require Dirichlet-to-Neumann operators, and $\tau_{j} G_{j}$ would be a Neumann-to-Dirichlet operator.

The linear system (18) can be solved by classical iterative methods as discussed in the next subsections. For this purpose, we consider the splitting $A=D+L+U$, where $D:=\operatorname{diag}(A)=I$,

$$
L:=\left[\begin{array}{ccccc}
0 & 0 & 0 & \cdots & 0 \\
\tau_{2} G_{1} & 0 & 0 & \cdots & 0 \\
\tau_{3} G_{1} & \tau_{3} G_{2} & 0 & \cdots & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
\tau_{J} G_{1} & \tau_{J} G_{2} & \tau_{J} G_{3} & \cdots & 0
\end{array}\right], U:=\left[\begin{array}{ccccc}
0 & \tau_{1} G_{2} & \tau_{1} G_{3} & \cdots & \tau_{1} G_{J} \\
0 & 0 & \tau_{2} G_{3} & \cdots & \tau_{2} G_{J} \\
0 & 0 & 0 & \cdots & \tau_{3} G_{J} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \cdots & 0
\end{array}\right] .
$$

This decomposition leads to the classical iterative methods based on the splitting $A=M-N$, and we can write the iterative methods in the standard form $M \widetilde{\mathbf{g}}^{k+1}=$ $N \widetilde{\mathbf{g}}^{k}+\mathbf{g}$ and in the difference form $M \boldsymbol{\delta}^{k+1}=N \boldsymbol{\delta}^{k}$, with $\boldsymbol{\delta}^{k}:=\widetilde{\mathbf{g}}^{k}-\widetilde{\mathbf{g}}^{k-1}$. The convergence of this class of methods is related to the iteration operator $G:=M^{-1} N$.

Let us now consider the block Gauss-Seidel method for the solution of system (18), which is obtained via the splitting $A=M-N$ with $M:=D+L$ and $N:=-U$,

$$
\begin{equation*}
(D+L) \widetilde{\mathbf{g}}^{k+1}=-U \widetilde{\mathbf{g}}^{k}+\mathbf{g} \tag{20}
\end{equation*}
$$

and the difference form is given by

$$
\begin{equation*}
(D+L) \boldsymbol{\delta}^{k+1}=-U \boldsymbol{\delta}^{k} \tag{21}
\end{equation*}
$$

More explicitly, the standard form is

$$
\left[\begin{array}{ccccc}
I_{1} & 0 & 0 & \cdots & 0  \tag{22}\\
\tau_{2} G_{1} & I_{2} & 0 & \cdots & 0 \\
\tau_{3} G_{1} & \tau_{3} G_{2} & I_{3} & \cdots & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
\tau_{J} G_{1} & \tau_{J} G_{2} & \tau_{J} G_{3} & \cdots & I_{J}
\end{array}\right]\left[\begin{array}{c}
\widetilde{g}_{1}^{k+1} \\
\widetilde{g}_{2}^{k+1} \\
\widetilde{g}_{3}^{k+1} \\
\vdots \\
\widetilde{g}_{J}^{k+1}
\end{array}\right]=\left[\begin{array}{c}
g_{1} \\
g_{2} \\
g_{3} \\
\vdots \\
g_{J}
\end{array}\right]-\left[\begin{array}{ccccc}
0 & \tau_{1} G_{2} & \tau_{1} G_{3} & \cdots & \tau_{1} G_{J} \\
0 & 0 & \tau_{2} G_{3} & \cdots & \tau_{2} G_{J} \\
0 & 0 & 0 & \cdots & \tau_{3} G_{J} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \cdots & 0
\end{array}\right]\left[\begin{array}{c}
\widetilde{g}_{1}^{k} \\
\widetilde{g}_{2}^{k} \\
\widetilde{g}_{3}^{k} \\
\vdots \\
\widetilde{g}_{J}^{k}
\end{array}\right] .
$$

Denoting the components of the difference $\boldsymbol{\delta}^{k}$ by $\delta_{j}^{k}:=\widetilde{g}_{j}^{k}-\widetilde{g}_{j}^{k-1}$, we can write the difference relation (21) in the explicit form

$$
\left[\begin{array}{ccccc}
I_{1} & 0 & 0 & \cdots & 0  \tag{23}\\
\tau_{2} G_{1} & I_{2} & 0 & \cdots & 0 \\
\tau_{3} G_{1} & \tau_{3} G_{2} & I_{3} & \cdots & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
\tau_{J} G_{1} & \tau_{J} G_{2} & \tau_{J} G_{3} & \cdots & I_{J}
\end{array}\right]\left[\begin{array}{c}
\delta_{1}^{k+1} \\
\delta_{2}^{k+1} \\
\delta_{3}^{k+1} \\
\vdots \\
\delta_{J}^{k+1}
\end{array}\right]=-\left[\begin{array}{ccccc}
0 & \tau_{1} G_{2} & \tau_{1} G_{3} & \cdots & \tau_{1} G_{J} \\
0 & 0 & \tau_{2} G_{3} & \cdots & \tau_{2} G_{J} \\
0 & 0 & 0 & \cdots & \tau_{3} G_{J} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \cdots & 0
\end{array}\right]\left[\begin{array}{c}
\delta_{1}^{k} \\
\delta_{2}^{k} \\
\delta_{3}^{k} \\
\vdots \\
\delta_{J}^{k}
\end{array}\right] .
$$

Now we show that (20) and (21) are equivalent to (14) and (10).
Theorem 2. Assume that $v^{0} \in H$ and $\widetilde{g}_{j}^{0}=\tau_{j} v^{0}$. Then the AltMR methods (14) and (10) are equivalent to (20) and (21).

Proof. Consider the transmission condition of (14), that is

$$
\begin{equation*}
\tau_{j} v_{j}^{k}=g_{j}-\sum_{\ell=j+1}^{J} \tau_{j} v_{\ell}^{k-1}-\sum_{\ell=1}^{j-1} \tau_{j} v_{\ell}^{k} \tag{24}
\end{equation*}
$$

for $j=1, \ldots, J$. We denote by $\widehat{g}_{j}^{k}=\tau_{j} v_{j}^{k}$ and, recalling that $v_{j}^{k-1}$ solves (14), we can use the operator $G_{j}$ defined above to write $v_{j}^{k-1}=G_{j}\left(\tau_{j} v_{j}^{k-1}\right)=G_{j}\left(\widehat{g}_{j}^{k-1}\right)$. Therefore, (24) can be equivalently written as

$$
\begin{equation*}
\widehat{g}_{j}^{k}=g_{j}-\sum_{\ell=1}^{j-1} \tau_{j} G_{j}\left(\widehat{g}_{\ell}^{k}\right)-\sum_{\ell=j+1}^{J} \tau_{j} G_{j}\left(\widehat{g}_{\ell}^{k-1}\right) \tag{25}
\end{equation*}
$$

If $\widehat{g}_{j}^{0}=\widetilde{g}_{j}^{0}$, then equation (25) is the equivalent pointwise form of (20)-(22) with $\widehat{g}_{j}^{k}=\widetilde{g}_{j}^{k}$. In a similar way, using the last equation in (10), we obtain the formula

$$
\delta_{j}^{k}=-\sum_{\ell=1}^{j-1} \tau_{j} G_{j}\left(\delta_{\ell}^{k}\right)-\sum_{\ell=j+1}^{J} \tau_{j} G_{j}\left(\delta_{\ell}^{k-1}\right)
$$

which is equivalent to (21)-(23).
3.4. Analogy with the alternating Schwarz method. In this section, we study the relations between the AltMR and the alternating Schwarz method (AltSM). In particular, we prove that if $J=2$, the AltMR is equivalent to the AltSM in the sense of Theorem 3. This identification does not hold in general for $J>2$.

To do so, consider problem (2), and define for $j=1, \ldots, J, \Omega_{j}:=\Omega \backslash \overline{O_{j}}$, and a smooth enough function $g: \Omega \rightarrow \mathbb{R}$ such that $g_{\mid \partial O_{j}}=g_{j}, g_{\mid \partial \Omega}=0$, and we assume that there exists a smooth extension of $u$ inside the objects that we denote by $\mathcal{E}_{j} u$. Then, we set

$$
\widetilde{u}:= \begin{cases}u-g & \text { in } \Omega_{j} \\ \mathcal{E}_{j} u-g & \text { in } O_{j}\end{cases}
$$

Notice that the extension $\mathcal{E}_{j} u$ must be smooth enough, e.g. twice differentiable, such that its Laplacian is well defined. Next, we define $f: \Omega \rightarrow \mathbb{R}$ as

$$
f:= \begin{cases}\Delta g & \text { in } \Omega_{j} \\ \Delta \mathcal{E}_{j} u+\Delta g & \text { in } \overline{O_{j}}\end{cases}
$$

Therefore, $\widetilde{u}$ solves the problem

$$
\begin{equation*}
\Delta \widetilde{u}=f \text { in } \Omega \text { with } \widetilde{u}=0 \text { on } \partial \Omega \tag{26}
\end{equation*}
$$

and by construction it satisfies $\widetilde{u}=0$ on $\bigcup_{j=1}^{J} \partial O_{j}$. We can now apply the AltSM to solve (26). To do so, we consider the domain decomposition $\Omega=\cup_{j=1}^{J} \Omega_{j}$, and denote by $\widetilde{u}_{j}^{k}$ the approximate solution at the $k$ th iteration on the subdomain $\Omega_{j}$. Starting with some initial approximation $\widetilde{u}_{j}^{0}$ for all $j$, the AltSM is defined as

$$
\begin{aligned}
\Delta \widetilde{u}_{j}^{k} & =f \text { in } \Omega_{j}, \widetilde{u}_{j}^{k}=0 \text { on } \partial \Omega \\
\widetilde{u}_{j}^{k} & =\sum_{\ell=1}^{j-1} \alpha_{j, \ell} \widetilde{u}_{\ell}^{k}+\sum_{\ell=j+1}^{J} \alpha_{j, \ell} \widetilde{u}_{\ell}^{k-1} \text { on } \partial O_{j}
\end{aligned}
$$

where the coefficients $\alpha_{j, \ell}$ are non-negative and such that $\sum_{\ell=1}^{J} \alpha_{j, \ell}=1$; see, e.g., $[15,43,13]$. This definition implies that the corresponding differences $\widetilde{d_{j}^{k}}:=$ $\widetilde{u}_{j}^{k}-\widetilde{u}_{j}^{k-1}$ solve

$$
\begin{align*}
\Delta \widetilde{d}_{j}^{k} & =0 \text { in } \Omega_{j}, \widetilde{d}_{j}^{k}=0 \text { on } \partial \Omega \\
\widetilde{d}_{j}^{k} & =\sum_{\ell=1}^{j-1} \alpha_{j, \ell} \widetilde{d}_{\ell}^{k}+\sum_{\ell=j+1}^{J} \alpha_{j, \ell} \widetilde{d}_{\ell}^{k-1} \text { on } \partial O_{j} . \tag{27}
\end{align*}
$$

Next, we prove in Theorem 3 the analogy between the AltMR and the AltSM for $J=2$, and afterwards we provide a counterexample to show that this analogy does not hold in general for $J>2$.
Theorem 3. Consider $J=2$ and assume that $\widetilde{d_{2}^{0}}=d_{2}^{0}$. Then for any $k \geq 1$ we have that $\widetilde{d}_{2}^{k}=d_{2}^{k}$ on $\bar{\Omega}_{1}$ and $\widetilde{d}_{1}^{k}=-d_{1}^{k}$ on $\bar{\Omega}_{2}$, where $\widetilde{d}_{j}^{k}$ and $d_{j}^{k}$ solve (27) and (10).

Proof. Since $J=2$, we have that $\alpha_{1,2}=\alpha_{2,1}=1$. We first prove the relation $\widetilde{d_{2}^{k}}=d_{2}^{k}$ by induction. The result is true for $k=0$ by assumption. Now, we assume that the relation holds for $k$ and prove that it is true for $k+1$ as well. Recalling the transmission conditions of (27) and (10) (for $J=2$ and $\alpha_{1,2}=\alpha_{2,1}=1$ ), it holds that

$$
\begin{equation*}
\widetilde{d}_{1}^{k+1}=\widetilde{d}_{2}^{k}=d_{2}^{k}=-d_{1}^{k+1} \text { on } \partial O_{1}, \tag{28}
\end{equation*}
$$

where we used the induction hypothesis. Equation (28), together with the existence of unique solutions to (27) and (10), implies that $\widetilde{d}_{1}^{k+1}=-d_{1}^{k+1}$ on $\bar{\Omega}_{1}$. Using this equality and the transmission conditions of (27) and (10), we have that

$$
\widetilde{d}_{2}^{k+1}=\widetilde{d}_{1}^{k+1}=-d_{1}^{k+1}=d_{2}^{k+1} \text { on } \partial O_{2}
$$

Hence the well-posedness of (27) and (10) implies that $\widetilde{d}_{2}^{k}=d_{2}^{k}$ on $\bar{\Omega}_{2}$, which is our first claim. The second relation follows directly by the first one together with (28).

The equivalence proved in Theorem 3 does not hold in general for $J>2$. In fact, we now show that classical choices of Schwarz methods (like the classical AltSM) does not lead to the result proved in Theorem 3. This negative result is suggested by the transmission conditions of (27) and (10). In fact, if we set $\alpha_{\ell, j}=-1$ for all $j, \ell$, the transmission condition of (27) coincides with the ones of (10). However, this is not possible in a Schwarz method framework, where the hypothesis $\sum_{\ell=1}^{J} \alpha_{j, \ell}=1$, with $\alpha_{j, \ell}$ non-negative, is required. Nevertheless, we provide the following example to show that the iterates $\widetilde{d}_{j}^{k}$ and $d_{j}^{k}$ are unrelated for $J>2$. To do so, consider a domain $\Omega=(0,1)$ and three holes $O_{j}=\left(a_{j}, b_{j}\right)$ for $j=1,2,3$ with $a_{j}=\frac{2 j-1}{7}$ and $b_{j}=\frac{2 j}{7}$. The alternating Schwarz method (27) becomes

$$
\begin{aligned}
& \left\{\begin{array}{c}
\Delta \widetilde{d}_{1}^{k}=0 \text { in } \Omega_{1}, \quad \widetilde{d}_{1}^{k}(0)=\widetilde{d}_{1}^{k}(1)=0, \\
\widetilde{d}_{1}^{k}(\widetilde{x})=\alpha_{1,2} \widetilde{d}_{2}^{k-1}(\widetilde{x})+\alpha_{1,3} \widetilde{d}_{3}^{k-1}(\widetilde{x}) \text { for } \widetilde{x}=a_{1}, b_{1},
\end{array}\right. \\
& \left\{\begin{aligned}
& \Delta \widetilde{d}_{2}^{k}=0 \text { in } \Omega_{2}, \quad \widetilde{d}_{2}^{k}(0)=\widetilde{d}_{2}^{k}(1)=0, \\
& \widetilde{d}_{2}^{k}(\widetilde{x})=\alpha_{2,1} \widetilde{d}_{1}^{k}(\widetilde{x})+\alpha_{2,3} \widetilde{d}_{3}^{k-1}(\widetilde{x}) \text { for } \widetilde{x}=a_{2}, b_{2},
\end{aligned}\right.
\end{aligned}
$$

$$
\left\{\begin{aligned}
\Delta \widetilde{d}_{3}^{k} & =0 \text { in } \Omega_{3}, \quad \widetilde{d}_{3}^{k}(0)=\widetilde{d}_{3}^{k}(1)=0, \\
\widetilde{d}_{3}^{k}(\widetilde{x}) & =\alpha_{3,1} \widetilde{d}_{1}^{k}(\widetilde{x})+\alpha_{3,2} \widetilde{d}_{2}^{k-1}(\widetilde{x}) \text { for } \widetilde{x}=a_{3}, b_{3}
\end{aligned}\right.
$$

Notice that by setting $\alpha_{\ell, j}=-1$, the above problems coincide with the AltMR (10). The general solutions to these three problems are

$$
\begin{gathered}
\widetilde{d}_{1}^{k}(x)=\left\{\begin{array}{ll}
\frac{A_{1}^{k} x}{a_{1}} & x \in\left[0, a_{1}\right], \\
\frac{B_{1}^{k}(1-x)}{1-b_{1}} & x \in\left[b_{1}, 1\right],
\end{array} \quad \widetilde{d}_{2}^{k}(x)= \begin{cases}\frac{A_{2}^{k} x}{a_{2}} & x \in\left[0, a_{2}\right], \\
\frac{B_{2}^{k}(1-x)}{1-b_{2}} & x \in\left[b_{2}, 1\right],\end{cases} \right. \\
\widetilde{d}_{3}^{k}(x)= \begin{cases}\frac{A_{3}^{k} x}{a_{3}} & x \in\left[0, a_{3}\right], \\
\frac{B_{3}^{k}(1-x)}{1-b_{3}} & x \in\left[b_{3}, 1\right],\end{cases}
\end{gathered}
$$

where $A_{j}^{k}$ and $B_{j}^{k}$ are constants depending on the transmission conditions. Defining $\mathbf{v}^{k}:=\left[A_{1}^{k}, B_{1}^{k}, A_{2}^{k}, B_{2}^{k}, A_{3}^{k}, B_{3}^{k}\right]^{\top}$ and using the transmission conditions, we obtain the iteration relation

$$
(I+\widetilde{L}) \mathbf{v}^{k}=-\widetilde{U} \mathbf{v}^{k-1}
$$

where

$$
\widetilde{L}=-\left[\begin{array}{cccccc}
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & \frac{\alpha_{2,1}\left(1-a_{2}\right)}{\left(1-b_{1}\right)} & 0 & 0 & 0 & 0 \\
0 & \frac{\alpha_{2,1}\left(1-b_{2}\right)}{\left(1-b_{1}\right)} & 0 & 0 & 0 & 0 \\
0 & \frac{\alpha_{3,1}\left(1-a_{3}\right)}{\left(1-b_{1}\right)} & 0 & \frac{\alpha_{3,2}\left(1-a_{3}\right)}{\left(1-b_{2}\right)} & 0 & 0 \\
0 & \frac{\alpha_{3,1}\left(1-b_{3}\right)}{\left(1-b_{1}\right)} & 0 & \frac{\alpha_{3,2}\left(1-b_{3}\right)}{\left(1-b_{2}\right)} & 0 & 0
\end{array}\right] \text { and } \widetilde{U}=-\left[\begin{array}{cccccc}
0 & 0 & \frac{\alpha_{1,2} a_{1}}{a_{2}} & 0 & \frac{\alpha_{1,3} a_{1}}{a_{3}} & 0 \\
0 & 0 & \frac{\alpha_{1,2} b_{1}}{a_{2}} & 0 & \frac{\alpha_{1,3} b_{1}}{a_{3}} & 0 \\
0 & 0 & 0 & 0 & \frac{\alpha_{2} a_{2}}{a_{3}} & 0 \\
0 & 0 & 0 & 0 & \frac{\alpha_{2,3}}{a_{3}} & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0
\end{array}\right] .
$$

Let us define the matrix $G=-\left(I_{6}+\widetilde{L}\right)^{-1} \widetilde{U}$, with $I_{6}$ the $6 \times 6$ identity. The iteration relation above then reads

$$
\mathbf{v}^{k}=G \mathbf{v}^{k-1}
$$

where the matrix $G$ depends on the weights $\alpha_{j, k}$. Notice that the same relation (with appropriately chosen weights in $G$ ) can be obtained for the AltMR. We consider three different cases. The first case is $\alpha_{1,2}=\alpha_{2,3}=\alpha_{3,1}=0$ and $\alpha_{1,3}=\alpha_{2,1}=$ $\alpha_{3,2}=1$, which corresponds to the classical AltSM, the second case is $\alpha_{j, k}=\frac{1}{2}$ for any $j, k$, which leads to a weighted AltSM, and the third case $\alpha_{j, k}=-1$ for any $j, k$, which corresponds to the AltMR. Computing explicitly the iteration matrix $G$ is these cases, we get
$G_{1}=\left[\begin{array}{cccccc}0 & 0 & 0 & 0 & \frac{1}{5} & 0 \\ 0 & 0 & 0 & 0 & \frac{2}{5} & 0 \\ 0 & 0 & 0 & 0 & \frac{8}{25} & 0 \\ 0 & 0 & 0 & 0 & \frac{6}{25} & 0 \\ 0 & 0 & 0 & 0 & \frac{4}{25} & 0 \\ 0 & 0 & 0 & 0 & \frac{2}{25} & 0\end{array}\right], G_{2}=\left[\begin{array}{cccccc}0 & 0 & \frac{1}{6} & 0 & \frac{1}{10} & 0 \\ 0 & 0 & \frac{1}{3} & 0 & \frac{1}{5} & 0 \\ 0 & 0 & \frac{2}{15} & 0 & \frac{19}{50} & 0 \\ 0 & 0 & \frac{1}{10} & 0 & \frac{23}{50} & 0 \\ 0 & 0 & \frac{1}{10} & 0 & \frac{29}{150} & 0 \\ 0 & 0 & \frac{1}{20} & 0 & \frac{29}{300} & 0\end{array}\right], G_{3}=\left[\begin{array}{cccccc}0 & 0 & -\frac{1}{3} & 0 & -\frac{1}{5} & 0 \\ 0 & 0 & -\frac{2}{3} & 0 & -\frac{2}{5} & 0 \\ 0 & 0 & \frac{8}{15} & 0 & -\frac{7}{25} & 0 \\ 0 & 0 & \frac{2}{5} & 0 & -\frac{14}{25} & 0 \\ 0 & 0 & 0 & 0 & \frac{8}{15} & 0 \\ 0 & 0 & 0 & 0 & \frac{4}{15} & 0\end{array}\right]$.
Their spectra are
$\sigma\left(G_{1}\right)=\left\{\frac{4}{25}, 0\right\}, \quad \sigma\left(G_{2}\right)=\left\{\frac{-(3 \sqrt{389}-49)}{300}, \frac{49+3 \sqrt{389}}{300}, 0\right\}, \quad \sigma\left(G_{3}\right)=\left\{\frac{8}{15}, 0\right\}$, and we thus obtain for the spectral radii

$$
\rho\left(G_{1}\right)=\frac{4}{25}=0.16, \quad \rho\left(G_{2}\right)=\frac{49+3 \sqrt{389}}{300} \approx 0.36, \quad \rho\left(G_{3}\right)=\frac{8}{15} \approx 0.53
$$

This shows that the iterations generated by $G_{1}$ and $G_{2}$ corresponding to alternating Schwarz methods must be different from the iterations generated by $G_{3}$ corresponding to the AltMR, and we have thus shown that there is in general no relation between the AltMR and the AltSM for $J>2$.

The reader may ask himself how the iteration matrices look like for $J=2$. To illustrate this, we consider a domain $\Omega=(0,1)$ with two holes $O_{j}=\left(a_{j}, b_{j}\right)$ for $j=1,2$ with $a_{j}=\frac{2 j-1}{5}$ and $b_{j}=\frac{2 j}{5}$, and recall that $\alpha_{1,2}=\alpha_{2,1}=1$. The same arguments as above allow us to obtain the iteration matrices

$$
G_{\mathrm{AltSM}}=\left[\begin{array}{cccccc}
0 & 0 & 0 & 0 & \frac{1}{3} & 0 \\
0 & 0 & 0 & 0 & \frac{2}{3} & 0 \\
0 & 0 & 0 & 0 & \frac{4}{9} & 0 \\
0 & 0 & 0 & 0 & \frac{2}{9} & 0
\end{array}\right], \quad G_{\mathrm{AltMR}}=\left[\begin{array}{cccccc}
0 & 0 & 0 & 0 & -\frac{1}{3} & 0 \\
0 & 0 & 0 & 0 & -\frac{2}{3} & 0 \\
0 & 0 & 0 & 0 & \frac{4}{9} & 0 \\
0 & 0 & 0 & 0 & \frac{2}{9} & 0
\end{array}\right]
$$

whose spectra coincide and are given by

$$
\sigma\left(G_{\mathrm{AltMR}}\right)=\sigma\left(G_{\mathrm{PSM}}\right)=\left\{\frac{4}{9}, 0\right\}
$$

It is clear that the two matrices generate similar iterates. Notice also the negative signs in $G_{\text {AltMR }}$ that produce $d_{1}^{k}$ with opposite sign to $\widetilde{d}_{1}^{k}$, in agreement with Theorem 3.
3.5. Convergence analysis for $J=2$ objects. Consider problem (2) for $J=2$ :

$$
\begin{align*}
\Delta u & =0 \text { in } \Omega \backslash\left(\partial O_{1} \cup \partial O_{2}\right) \\
u & =0 \text { on } \partial \Omega \\
u & =g_{1} \text { on } \partial O_{1}  \tag{29}\\
u & =g_{2} \text { on } \partial O_{2}
\end{align*}
$$

The AltMR (14) for the solution of (29) is

$$
\begin{align*}
\Delta v_{1}^{k} & =0 \text { in } \Omega \backslash \partial O_{1}, v_{1}^{k}=0 \text { on } \partial \Omega, & \Delta v_{2}^{k} & =0 \text { in } \Omega \backslash \partial O_{2}, v_{2}^{k}=0 \text { on } \partial \Omega, \\
v_{1}^{k} & =g_{1}-v_{2}^{k-1} \text { on } \partial O_{1}, & v_{2}^{k} & =g_{2}-v_{1}^{k} \text { on } \partial O_{2} . \tag{30}
\end{align*}
$$

Let us define the error at the $k$ th iteration by $e_{j}^{k}:=v_{j}-v_{j}^{k}$ for $j=1, \ldots, J$. In terms of the errors, the AltMR (30) reads

$$
\begin{align*}
\Delta e_{1}^{k} & =0 \text { in } \Omega \backslash \partial O_{1}, e_{1}^{k}=0 \text { on } \partial \Omega, & \Delta e_{2}^{k} & =0 \text { in } \Omega \backslash \partial O_{2}, e_{2}^{k}=0 \text { on } \partial \Omega,  \tag{31}\\
e_{1}^{k} & =-e_{2}^{k-1} \text { on } \partial O_{1}, & e_{2}^{k} & =-e_{1}^{k} \text { on } \partial O_{2} .
\end{align*}
$$

We can now prove the following theorem using similar techniques considered in [8].
Theorem 4. The AltMR for the solution to problem (29) converges geometrically, in the sense that

$$
\max _{\bar{\Omega}}\left|e_{j}^{k}\right| \leq \rho_{\mathrm{AltMR}}^{k} \max _{\bar{\Omega}}\left|e_{j}^{0}\right|
$$

for $j=1,2$, where $\rho_{\mathrm{AltMR}}=\left(\max _{\partial O_{1}} w_{2}\right)\left(\max _{\partial O_{2}} w_{1}\right)<1$, and $w_{j}$ solves

$$
\begin{equation*}
\Delta w_{j}=0 \text { in } \Omega \backslash \partial O_{j}, w_{j}=0 \text { on } \partial \Omega, w_{j}=1 \text { on } \partial O_{j} \tag{32}
\end{equation*}
$$

Proof. The functions $w_{j}$, for $j=1,2$ which are solutions to (32), satisfy because of the maximum principle that $w_{j} \geq 0$ and $\left|e_{j}^{k}\right| \leq w_{j} \max _{\partial O_{j}}\left|e_{j}^{k}\right|$. Now, since $e_{1}^{k}$


Figure 2. A unit disc $\Omega$ with two objects $O_{1}$ and $O_{2}$, which are two discs having the same radius denoted by $r$. The distance between the centers of $O_{1}$ and $O_{2}$ and the center of $\Omega$ is denoted by $\delta$.
and $e_{2}^{k}$ are harmonic functions and solve (31), the maximum principle allows us to write

$$
\begin{aligned}
\max _{\bar{\Omega}}\left|e_{1}^{k}\right| & =\max _{\partial O_{1}}\left|e_{1}^{k}\right|=\max _{\partial O_{1}}\left|e_{2}^{k-1}\right| \leq \max _{\partial O_{1}} w_{2} \max _{\partial O_{2}}\left|e_{2}^{k-1}\right|=\max _{\partial O_{1}} w_{2} \max _{\bar{\Omega}}\left|e_{2}^{k-1}\right| \\
\max _{\bar{\Omega}}\left|e_{2}^{k-1}\right| & =\max _{\partial O_{2}}\left|e_{2}^{k-1}\right|=\max _{\partial O_{2}}\left|e_{1}^{k-1}\right| \leq \max _{\partial O_{2}} w_{1} \max _{\partial O_{1}}\left|e_{1}^{k-1}\right|=\max _{\partial O_{2}} w_{1} \max _{\bar{\Omega}}\left|e_{1}^{k-1}\right|
\end{aligned}
$$

By combining these two estimates we get

$$
\max _{\bar{\Omega}}\left|e_{1}^{k}\right| \leq \max _{\partial O_{1}} w_{2} \max _{\bar{\Omega}}\left|e_{2}^{k-1}\right| \leq \max _{\partial O_{1}} w_{2} \max _{\partial O_{2}} w_{1} \max _{\bar{\Omega}}\left|e_{1}^{k-1}\right|=\rho_{\text {AltMR }} \max _{\bar{\Omega}}\left|e_{1}^{k-1}\right|
$$

Since $w_{1}$ and $w_{2}$ are harmonic functions in $\Omega \backslash \partial O_{1}$ and $\Omega \backslash \partial O_{2}$, Assumption (1) and the maximum principle imply that $\max _{\partial O_{1}} w_{2}<1$ and $\max _{\partial O_{2}} w_{1}<1$. Hence $\rho_{\text {AltMR }}<1$ and our proof is complete.

This theorem allows us to compute explicitly the contraction factor $\rho_{\text {AltMR }}$ as a function of the geometry of the domain. For example, consider a domain whose geometry is shown in Figure 2, where $\Omega \subset \mathbb{R}^{2}$ is the unit disc, the objects $O_{1}$ and $O_{2}$ are two discs whose centers are aligned on a straight line passing trough the center of $\Omega \subset \mathbb{R}^{2}$. The two discs have the same radius $r$ and the distance between the center of each of them and the center of $\Omega$ is denoted by $\delta$. We can prove the following result.

Corollary 1. Consider the problem (29) defined on a domain whose geometry is depicted in Figure 2. We have that

$$
\rho_{\mathrm{AltMR}}(r, \delta)=\left(\frac{\log \left(\frac{(-a-\delta+r)^{2}}{(a(r-\delta)-1)^{2}}\right)}{\log \left(\frac{(-a+\delta-r)^{2}}{(a(\delta-r)-1)^{2}}\right)}\right)^{2}
$$

where

$$
\begin{equation*}
a=\frac{1+\delta^{2}-r^{2}-\sqrt{r^{4}+\left(-2-2 \delta^{2}\right) r^{2}+\delta^{4}-2 \delta^{2}+1}}{2 \delta} . \tag{33}
\end{equation*}
$$



Figure 3. Möbius transformation of the unit disc $\Omega$ (dashed line) with a hole $O_{1}$ (solid line) into an annulus.

Proof. To prove the result, we first need to solve problems (32) and compute the $w_{j}$. To do so, we recall the Möbius transformation $h: \mathbb{C} \rightarrow \mathbb{C}$ :

$$
\begin{equation*}
h(z):=\frac{z-a}{\bar{a} z-1} \tag{34}
\end{equation*}
$$

which maps the unit circle into itself and circles into circles. In particular, by imposing the conditions $h(\delta-r)=R$ and $h(\delta+r)=-R$, one obtains a (real) coefficient $a$, as in (33), such that $h$ maps $\Omega \backslash \overline{O_{j}}$ into an annulus; see Figure 3. We consider now the problem

$$
\begin{equation*}
\Delta \widetilde{w}_{1}=0 \text { in } h(\Omega) \backslash h\left(\partial O_{1}\right), \widetilde{w}_{1}=0 \text { on } h(\partial \Omega), \widetilde{w}_{1}=1 \text { on } h\left(\partial O_{1}\right) \tag{35}
\end{equation*}
$$

Since (35) is radially symmetric, a separation of variables argument allows us to compute its solution,

$$
\widetilde{w}_{1}(\zeta)=\frac{\log \left(\sqrt{\operatorname{Re}(\zeta)^{2}+\operatorname{Im}(\zeta)^{2}}\right)}{\log (R)}
$$

By transforming back from $\zeta$ to $z$, we obtain

$$
w_{1}(\operatorname{Re}(z), \operatorname{Im}(z))=\widetilde{w}_{1}(h(z))=\frac{\log (|h(z)|)}{\log (R)}=\frac{\log \left(|h(z)|^{2}\right)}{\log \left(|R|^{2}\right)}=\frac{\log \left(\left|\frac{z-a}{a z-1}\right|^{2}\right)}{\log \left(\left|\frac{\delta-r-a}{a(\delta-r)-1}\right|^{2}\right)}
$$

which solves (32) for $j=1$ since $h$ is a conformal mapping. A further simplification leads to

$$
w_{1}(x, y)=\frac{\log \left(\frac{(x-a)^{2}+y^{2}}{(a x-1)^{2}+(a y)^{2}}\right)}{\log \left(\frac{(\delta-r-a)^{2}}{(a(\delta-r)-1)^{2}}\right)}
$$

Now, since $\widetilde{w}_{1}$ is symmetric with respect to the $\operatorname{Re}(\zeta)$-axis and decays monotonically in any radial direction, the same holds for $w_{1}$. Hence, the maximum of $w_{1}$ along $\partial O_{2}$ is attained at $(x, y)=(-\delta+r, 0)$ and has the value

$$
\max _{\partial O_{2}} w_{1}=w_{1}(-\delta+r, 0)=\frac{\log \left(\frac{(-\delta+r-a)^{2}}{(a(-\delta+r)-1)^{2}}\right)}{\log \left(\frac{(\delta-r-a)^{2}}{(a(\delta-r)-1)^{2}}\right)} .
$$

Since the solution $w_{2}$ to problem (32) for $j=2$ can be obtained by rotating $w_{1}$ by $\pi$ around the origin, it holds that

$$
\max _{\partial O_{1}} w_{2}=w_{2}(\delta-r, 0)=w_{1}(-\delta+r, 0)=\max _{\partial O_{2}} w_{1} .
$$

Recalling from Theorem 4 that $\rho_{\text {AltMR }}=\left(\max _{\partial O_{1}} w_{2}\right)\left(\max _{\partial O_{2}} w_{1}\right)$, the claim follows.

## 4. The parallel method of Reflections

The parallel method of reflections (PMR) was introduced by Golusin in 1934 [21] and formally recalled by Traytak [44] for the Laplace equation. Ichiki and Brady [27] present exactly the parallel method of Golusin, and they state: "It is easy to extend this procedure to the $N$ body problem by superposing distances by other particles". In practice, the parallel version is obtained by replacing in the righthand side of the boundary condition (9) the differences at the iteration $k$ with the corresponding ones at the iteration $k-1$. Hence, problem (10) becomes

$$
\begin{align*}
\Delta d_{j}^{k} & =0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{k}=0 \text { on } \partial \Omega, \\
d_{j}^{k} & =-\sum_{\ell=1, \ell \neq j}^{J} d_{\ell}^{k-1} \text { on } \partial O_{j} . \tag{36}
\end{align*}
$$

The sequences $\left\{d_{j}^{k}\right\}_{k \in \mathbb{N}^{+}}$are initialized by solving for each $j=1, \ldots, J$ the problem

$$
\begin{equation*}
\Delta d_{j}^{1}=0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{1}=0 \text { on } \partial \Omega, d_{j}^{1}=g_{j}-u^{0} \text { on } \partial O_{j}, \tag{37}
\end{equation*}
$$

for some $u^{0} \in H$, and the approximate solution at the $k$ th iteration is defined by

$$
\begin{equation*}
u^{k}=u^{k-1}+\sum_{j=1}^{J} d_{j}^{k} \tag{38}
\end{equation*}
$$

The PMR (36)-(37), as presented in [34], leads to Algorithm 3.

```
Algorithm 3 Parallel Method of Reflection (PMR)
Input: \(K\) (maximum number of iterations), tol (tolerance).
    Set \(u^{0} \in H\) and \(k=1\).
    for \(j=1: J\) (this loop is executed in parallel) do
        Compute \(d_{j}^{1}\) solving problem (37).
    end for
    Compute the approximation \(u^{1}=\sum_{j=1}^{J} d_{j}^{1}\).
    while \(k<K\) and \(\left\|u^{k}-u^{k-1}\right\|>\) tol do
        Update \(k=k+1\).
        for \(j=2: J\) (this loop is executed in parallel) do
            Compute \(d_{j}^{k}\) solving problem (36).
        end for
        Compute the approximation \(u^{k}=u^{k-1}+\sum_{j=1}^{J} d_{j}^{k}\).
    end while
```

Now, as in Section 3, we introduce the variable $v_{j}^{k}$ defined in (12), and we write the PMR in terms of $v_{j}^{k}$. To do so, we first use the transmission condition of (36)
and (38) to compute on $\partial O_{j}$ that

$$
d_{j}^{k}=-\sum_{\ell=1, \ell \neq j}^{J} d_{\ell}^{k-1}=d_{j}^{k-1}-\sum_{\ell=1}^{J} d_{\ell}^{k-1}=d_{j}^{k-1}-\left(u^{k-1}-u^{k-2}\right)
$$

which implies that

$$
d_{j}^{k}+u^{k-1}=d_{j}^{k-1}+u^{k-2}=\cdots=d_{j}^{1}+u^{0}
$$

Recalling that $d_{j}^{1}=g_{j}-u^{0}$ on $\partial O_{j}$, we obtain

$$
\begin{equation*}
d_{j}^{k}=g_{j}-u^{k-1} \text { on } \partial O_{j} \tag{39}
\end{equation*}
$$

Therefore, we have that (36) is equivalent to

$$
\begin{align*}
\Delta d_{j}^{k} & =0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{k}=0 \text { on } \partial \Omega \\
d_{j}^{k} & =g_{j}-u^{k-1} \text { on } \partial O_{j} \tag{40}
\end{align*}
$$

which is the parallel version of (8). Now, starting with the definition (12) and using (39) and (13), we compute on $\partial O_{j}$ that

$$
v_{j}^{k}=v_{j}^{k-1}+d_{j}^{k}=v_{j}^{k-1}+g_{j}-u^{k-1}=v_{j}^{k-1}+g_{j}-\sum_{\ell=1}^{J} v_{\ell}^{k-1}=g_{j}-\sum_{\ell=1, \ell \neq j}^{J} v_{\ell}^{k-1}
$$

Recalling that $v_{j}^{k} \in H_{j}$, we conclude that $v_{j}^{k}$ solves

$$
\begin{align*}
\Delta v_{j}^{k} & =0 \text { in } \Omega \backslash \partial O_{j}, v_{j}^{k}=0 \text { on } \partial \Omega \\
v_{j}^{k} & =g_{j}-\sum_{\ell=j+1}^{J} v_{\ell}^{k-1}-\sum_{\ell=1}^{j-1} v_{\ell}^{k-1} \text { on } \partial O_{j} . \tag{41}
\end{align*}
$$

The equivalence between (41) and (36)-(37) is proved in the following theorem.
Theorem 5. Consider the sequences $\left\{u^{k}\right\}_{k \in \mathbb{N}} \subset H,\left\{d_{j}^{k}\right\}_{k \in \mathbb{N}^{+}}$with $d_{j}^{k} \in H_{j}$, and $\left\{v_{j}^{k}\right\}_{k \in \mathbb{N}}$ with $v_{j}^{k} \in H_{j}$ for $j=1, \ldots, J$, and assume that $v_{j}^{k}=v_{j}^{k-1}+d_{j}^{k}$ for $k \geq 1$ and $j=1, \ldots, J$. Then the following statements are equivalent:
(a) the $d_{j}^{k}$ solve (36)-(37) with $u^{k}=u^{k-1}+\sum_{j=1}^{J} d_{j}^{n}=\sum_{n=1}^{k} \sum_{j=1}^{J} d_{j}^{n}$.
(b) the $v_{j}^{k}$ solve (41) with $u^{k}=\sum_{j=1}^{J} v_{j}^{k}$.

Proof. The implication $(\mathrm{a}) \Rightarrow(\mathrm{b})$ is proved in the discussion before the theorem, whereas $(\mathrm{b}) \Rightarrow(\mathrm{a})$ follows by similar arguments as in Theorem 1.

The PMR algorithm in terms of $v_{j}^{k}$ is given in Algorithm 4.
4.1. Substructured PMR as a block Jacobi method. In this section, we prove that the PMR is equivalent to the block Jacobi method. We recall that Ichiki and Brady [27, page 351] mention, without motivation or any rigorous argument, that "this iterative method (the PMR) is equivalent to the block Jacobi method". This sentence suggests a block Jacobi substructured interpretation of the PMR, as the block Gauss-Seidel one for the AltMR proved in Section 3.3. To see it, consider the operator $A$ defined in Section 3.3 and the decomposition $A=M-N$ with $M:=D$ and $N:=-(L+U)$. The block Jacobi method in the standard form is given by

$$
\begin{equation*}
D \widetilde{\mathbf{g}}^{k+1}=-(L+U) \widetilde{\mathbf{g}}^{k}+\mathbf{g} \tag{42}
\end{equation*}
$$

```
Algorithm 4 Parallel Method of Reflection (PMR)
Input: \(K\) (maximum number of iteration), tol (tolerance).
    Set \(v_{j}^{0} \in H_{j}\) for \(j=1, \ldots, J\), and \(k=1\).
    for \(j=1: J\) (this loop is executed in parallel) do
        Compute \(v_{j}^{1}\) solving problem (41).
    end for
    Compute the approximation \(u^{1}=\sum_{j=1}^{J} v_{j}^{1}\).
    while \(k<K\) and \(\left\|u^{k}-u^{k-1}\right\|>\) tol do
        Update \(k=k+1\).
        for \(j=1: J\) (this loop is executed in parallel) do
            Compute \(v_{j}^{k}\) solving problem (41).
        end for
        Compute the approximation \(u^{k}=\sum_{j=1}^{J} v_{j}^{k}\).
    end while
```

and in the difference form it is

$$
\begin{equation*}
D \boldsymbol{\delta}^{k+1}=-(L+U) \boldsymbol{\delta}^{k} \tag{43}
\end{equation*}
$$

where $\boldsymbol{\delta}^{k}:=\widetilde{\mathbf{g}}^{k}-\widetilde{\mathbf{g}}^{k-1}$, see also (23). More explicitly, the standard form is

$$
\left[\begin{array}{c}
\widetilde{g}_{1}^{k+1}  \tag{44}\\
\widetilde{g}_{2}^{k+1} \\
\widetilde{g}_{3}^{k+1} \\
\vdots \\
\widetilde{g}_{J}^{k+1}
\end{array}\right]=\left[\begin{array}{c}
g_{1} \\
g_{2} \\
g_{3} \\
\vdots \\
g_{J}
\end{array}\right]-\left[\begin{array}{ccccc}
0 & \tau_{1} G_{2} & \tau_{1} G_{3} & \cdots & \tau_{1} G_{J} \\
\tau_{2} G_{1} & 0 & \tau_{2} G_{3} & \cdots & \tau_{2} G_{J} \\
\tau_{3} G_{1} & \tau_{3} G_{2} & 0 & \cdots & \tau_{3} G_{J} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
\tau_{J} G_{1} & \tau_{J} G_{2} & \tau_{J} G_{3} & \cdots & 0
\end{array}\right]\left[\begin{array}{c}
\widetilde{g}_{1}^{k} \\
\widetilde{g}_{2}^{k} \\
\widetilde{g}_{3}^{k} \\
\vdots \\
\widetilde{g}_{J}^{k}
\end{array}\right]
$$

whereas the difference form (43) is

$$
\left[\begin{array}{c}
\delta_{1}^{k+1}  \tag{45}\\
\delta_{2}^{k+1} \\
\delta_{3}^{k+1} \\
\vdots \\
\delta_{J}^{k+1}
\end{array}\right]=-\left[\begin{array}{ccccc}
0 & \tau_{1} G_{2} & \tau_{1} G_{3} & \cdots & \tau_{1} G_{J} \\
\tau_{2} G_{1} & 0 & \tau_{2} G_{3} & \cdots & \tau_{2} G_{J} \\
\tau_{3} G_{1} & \tau_{3} G_{2} & 0 & \cdots & \tau_{3} G_{J} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
\tau_{J} G_{1} & \tau_{J} G_{2} & \tau_{J} G_{3} & \cdots & 0
\end{array}\right]\left[\begin{array}{c}
\delta_{1}^{k} \\
\delta_{2}^{k} \\
\delta_{3}^{k} \\
\vdots \\
\delta_{J}^{k}
\end{array}\right] .
$$

The next theorem shows that (42)-(44) and (43)-(45) are equivalent to the PMR.
Theorem 6. Assume that $v^{0} \in H$ and $\widetilde{g}_{j}^{0}=\tau_{j} v^{0}$. Then the PMR forms (41) and (40) are equivalent to (42) and (43).

Proof. We proceed as in Section 3.3 and consider the transmission condition of (41), that is

$$
\begin{equation*}
\tau_{j} v_{j}^{k}=g_{j}-\sum_{\ell=j+1}^{J} \tau_{j} v_{\ell}^{k-1}-\sum_{\ell=1}^{j-1} \tau_{j} v_{\ell}^{k-1} \tag{46}
\end{equation*}
$$

We define $\widehat{g}_{j}^{k}:=\tau_{j} v_{j}^{k}$ and recall that $v_{j}^{k-1}$ can be written as $v_{j}^{k-1}=G_{j}\left(\tau_{j} v_{j}^{k-1}\right)=$ $G_{j}\left(\widehat{g}_{j}^{k-1}\right)$. Hence, (46) becomes

$$
\widehat{g}_{j}^{k}=g_{j}-\sum_{\ell=j+1}^{J} \tau_{j} G_{j}\left(\widehat{g}_{\ell}^{k-1}\right)-\sum_{\ell=1}^{j-1} \tau_{j} G_{j}\left(\widehat{g}_{\ell}^{k-1}\right)
$$

which is equivalent to $(42)-(44)$ if $\widehat{g}_{j}^{0}=\widetilde{g}_{j}^{0}$. The same arguments can be used to show that (36) is equivalent to (43)-(45).
4.2. Analogies with Schwarz methods. Similar to Section 3.4, we show now that if $J=2$, then the PMR is the analog of the classical parallel Schwarz method (PSM), and if $J>2$, the PMR and the PSM define different iterates. To do so, we recall (26) and apply the PSM,

$$
\begin{align*}
\Delta \widetilde{u}_{j}^{k} & =f \text { in } \Omega_{j}, \widetilde{u}_{j}^{k}=0 \text { on } \partial \Omega \\
\widetilde{u}_{j}^{k} & =\sum_{\ell=1, \ell \neq j}^{J} \alpha_{j, \ell} \widetilde{u}_{\ell}^{k-1} \text { on } \partial O_{j} \tag{47}
\end{align*}
$$

where the coefficients $\alpha_{j, \ell}$ are non-negative and such that $\sum_{\ell} \alpha_{j, \ell}=1$. This definition implies that the corresponding differences $\widetilde{d}_{j}^{k}:=u_{j}^{k}-u_{j}^{k-1}$ solve

$$
\begin{align*}
\Delta \widetilde{d}_{j}^{k} & =0 \text { in } \Omega_{j}, \widetilde{d}_{j}^{k}=0 \text { on } \partial \Omega, \\
\widetilde{d}_{j}^{k} & =\sum_{\ell=1, \ell \neq j}^{J} \alpha_{j, \ell} \widetilde{d}_{\ell}^{k-1} \text { on } \partial O_{j} . \tag{48}
\end{align*}
$$

Theorem 7. Consider $J=2$ and assume that $\widetilde{d_{2}^{0}}=d_{2}^{0}$. Then, for any $k \geq 0$ we have that $\widetilde{d_{1}^{k}}=(-1)^{k} d_{1}^{k}$ on $\bar{\Omega}_{1}$ and $\widetilde{d_{2}^{k}}=(-1)^{k} d_{2}^{k}$ on $\bar{\Omega}_{2}$, where $\widetilde{d_{j}^{k}}$ and $d_{j}^{k}$ solve (48) and (41).

Proof. Since $J=2$, we have that $\alpha_{1,2}=\alpha_{2,1}=1$. We only prove the first relation $\widetilde{d_{1}^{k}}=(-1)^{k} d_{1}^{k}$, since the second follows by the same arguments. We proceed by induction. The statement is true for $k=0$ by assumption. Now, we assume that $\widetilde{d_{1}^{k}}=(-1)^{k} d_{1}^{k}$ holds and we show that it remains true for $k+1$. Recalling the transmission condition of (48) and (41), we have

$$
d_{1}^{k+1}=-d_{2}^{k}=-(-1)^{k} \widetilde{d}_{2}^{k}=(-1)^{k+1} \widetilde{d}_{2}^{k}=(-1)^{k} \widetilde{d}_{1}^{k+1} \text { on } \partial O_{1}
$$

where we used the induction hypothesis. Since (48) and (41) admits unique solutions, we have that $d_{1}^{k+1}=(-1)^{k} \widetilde{d}_{1}^{k+1}$ in $\bar{\Omega}_{1}$, which is our claim.

Now, to show that for $J>2$ the analogy proved in Theorem 7 is not in general true, we consider the same example provided in Section 3.4 and use the same notation. In this case, the PSM is

$$
\begin{aligned}
& \left\{\begin{aligned}
\Delta \widetilde{d}_{1}^{k} & =0 \text { in } \Omega_{1}, \quad \widetilde{d}_{1}^{k}(0)=\widetilde{d}_{1}^{k}(1)=0, \\
\widetilde{d}_{1}^{k}(\widetilde{x}) & =\alpha_{1,2} \widetilde{d}_{2}^{k-1}(\widetilde{x})+\alpha_{1,3} \widetilde{d}_{3}^{k-1}(\widetilde{x}) \text { for } \widetilde{x}=a_{1}, b_{1},
\end{aligned}\right. \\
& \left\{\begin{aligned}
\Delta \widetilde{d}_{2}^{k} & =0 \text { in } \Omega_{2}, \quad \widetilde{d}_{2}^{k}(0)=\widetilde{d}_{2}^{k}(1)=0, \\
\widetilde{d}_{2}^{k}(\widetilde{x}) & =\alpha_{2,1} \widetilde{d}_{1}^{k-1}(\widetilde{x})+\alpha_{2,3} \widetilde{d}_{2}^{k-1}(\widetilde{x}) \text { for } \widetilde{x}=a_{2}, b_{2},
\end{aligned}\right. \\
& \left\{\begin{aligned}
\Delta \widetilde{d}_{3}^{k} & =0 \text { in } \Omega_{3}, \quad \widetilde{d}_{3}^{k}(0)=\widetilde{d}_{3}^{k}(1)=0, \\
\widetilde{d}_{3}^{k}(\widetilde{x}) & =\alpha_{3,1} \widetilde{d}_{1}^{k-1}(\widetilde{x})+\alpha_{3,2} \widetilde{d}_{3}^{k-1}(\widetilde{x}) \text { for } \widetilde{x}=a_{3}, b_{3} .
\end{aligned}\right.
\end{aligned}
$$

The corresponding iteration in terms of the constants $A_{j}^{k}$ and $B_{j}^{k}$ is given by $\mathbf{v}^{k}=$ $G \mathbf{v}^{k-1}$, where $G:=-(\widetilde{L}+\widetilde{U})$, with $G$ depending on the weights $\alpha_{\ell, j}$. Notice that the same relation (with appropriately chosen weights in $G$ ) can be obtained for the PMR. We consider three different cases. The first case is $\alpha_{1,2}=\alpha_{2,3}=\alpha_{3,1}=0$
and $\alpha_{1,3}=\alpha_{2,1}=\alpha_{3,2}=1$, which corresponds to the classical PSM, the second case is $\alpha_{j, k}=\frac{1}{2}$ for any $j, k$, which leads to a weighted PSM, and the third case $\alpha_{j, k}=-1$ for any $j, k$, which corresponds to the PMR. Computing explicitly the iteration matrix $G$ is these cases, we get

$$
G_{1}=\left[\begin{array}{cccccc}
0 & 0 & 0 & 0 & \frac{1}{5} & 0 \\
0 & 0 & 0 & 0 & \frac{2}{5} & 0 \\
0 & \frac{4}{5} & 0 & 0 & 0 & 0 \\
0 & \frac{3}{5} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & \frac{2}{3} & 0 & 0 \\
0 & 0 & 0 & \frac{1}{3} & 0 & 0
\end{array}\right], G_{2}=\left[\begin{array}{cccccc}
0 & 0 & \frac{1}{6} & 0 & \frac{1}{10} & 0 \\
0 & 0 & \frac{1}{3} & 0 & \frac{1}{5} & 0 \\
0 & \frac{2}{5} & 0 & 0 & \frac{3}{10} & 0 \\
0 & \frac{3}{10} & 0 & 0 & \frac{2}{5} & 0 \\
0 & \frac{1}{5} & 0 & \frac{1}{3} & 0 & 0 \\
0 & \frac{1}{10} & 0 & \frac{1}{6} & 0 & 0
\end{array}\right], G_{3}=-\left[\begin{array}{cccccc}
0 & 0 & \frac{1}{3} & 0 & \frac{1}{5} & 0 \\
0 & 0 & \frac{2}{3} & 0 & \frac{2}{5} & 0 \\
0 & \frac{4}{5} & 0 & 0 & \frac{3}{5} & 0 \\
0 & \frac{3}{5} & 0 & 0 & \frac{4}{5} & 0 \\
0 & \frac{2}{5} & 0 & \frac{2}{3} & 0 & 0 \\
0 & \frac{1}{5} & 0 & \frac{1}{3} & 0 & 0
\end{array}\right],
$$

where $G_{1}, G_{2}$, and $G_{3}$ correspond to the three cases. Their spectra are

$$
\begin{aligned}
& \sigma\left(G_{1}\right)=\left\{\frac{i \sqrt{3} 4^{1 / 3}-4^{1 / 3}}{225^{1 / 3}},-\frac{4^{1 / 3}+i \sqrt{3} 4^{1 / 3}}{225^{1 / 3}}, \frac{4^{1 / 3}}{25^{1 / 3}}, 0\right\} \\
& \sigma\left(G_{2}\right)=\left\{-\frac{3+\sqrt{39}}{30}, \frac{\sqrt{39}-3}{30},-\frac{\sqrt{219}-3}{30}, \frac{\sqrt{219}+3}{30}, 0\right\} \\
& \sigma\left(G_{3}\right)=\left\{\frac{3+\sqrt{39}}{15},-\frac{\sqrt{39}-3}{15}, \frac{\sqrt{219}-3}{15},-\frac{\sqrt{219}+3}{15}, 0\right\}
\end{aligned}
$$

and the corresponding spectral radii are

$$
\rho\left(G_{1}\right)=\frac{4^{1 / 3}}{25^{1 / 3}} \approx 0.54, \quad \rho\left(G_{2}\right)=\frac{\sqrt{219}+3}{30} \approx 0.59, \quad \rho\left(G_{3}\right)=\frac{\sqrt{219}+3}{15} \approx 1.18
$$

We see again that the iterations generated by $G_{1}$ and $G_{2}$ corresponding to parallel Schwarz methods are different from the iterations generated by $G_{3}$ corresponding to the PMR, and we have therefore shown that there is no relation in general between PMR and PSM for $J>2$.
4.3. Convergence analysis for $J=2$ objects. Next, we prove for the PSM the same convergence result obtained for the AltMR in Theorem 4. The PMR for the errors $e_{1}^{k}$ and $e_{2}^{k}$ is

$$
\begin{align*}
\Delta e_{1}^{k} & =0 \text { in } \Omega \backslash \partial O_{1}, e_{1}^{k}=0 \text { on } \partial \Omega, & \Delta e_{2}^{k} & =0 \text { in } \Omega \backslash \partial O_{2}, e_{2}^{k}=0 \text { on } \partial \Omega,  \tag{49}\\
e_{1}^{k} & =-e_{2}^{k-1} \text { on } \partial O_{1}, & e_{2}^{k} & =-e_{1}^{k-1} \text { on } \partial O_{2}
\end{align*}
$$

Theorem 8. The PMR for the solution to problem (29) converges geometrically, in the sense that $\max _{\bar{\Omega}}\left|e_{j}^{2 k}\right| \leq \rho_{\mathrm{PMR}}^{2 k} \max _{\bar{\Omega}}\left|e_{j}^{0}\right|$, for $j=1,2$, where $\rho_{\mathrm{PMR}}=\sqrt{\rho_{\mathrm{AltMR}}}$ and $\rho_{\mathrm{AltMR}}$ is given in Corollary 1.

Proof. As in Corollary 1, the functions $w_{j}$ solving (32) satisfy by the maximum principle that $w_{j} \geq 0$ and $\left|e_{j}^{k}\right| \leq w_{j} \max _{\partial O_{j}}\left|e_{j}^{k}\right|$. Now, since $e_{1}^{k}$ and $e_{2}^{k}$ are harmonic functions and solve (31), the maximum principle allows us to write

$$
\begin{aligned}
\max _{\bar{\Omega}}\left|e_{1}^{k}\right| & =\max _{\partial O_{1}}\left|e_{1}^{k}\right|=\max _{\partial O_{1}}\left|e_{2}^{k-1}\right| \leq \max _{\partial O_{1}} w_{2} \max _{\partial O_{2}}\left|e_{2}^{k-1}\right|=\max _{\partial O_{1}} w_{2} \max _{\bar{\Omega}}\left|e_{2}^{k-1}\right|, \\
\max _{\bar{\Omega}}\left|e_{2}^{k-1}\right| & =\max _{\partial O_{2}}\left|e_{2}^{k-1}\right|=\max _{\partial O_{2}}\left|e_{1}^{k-2}\right| \leq \max _{\partial O_{2}} w_{1} \max _{\partial O_{1}}\left|e_{1}^{k-2}\right|=\max _{\partial O_{2}} w_{1} \max _{\bar{\Omega}}\left|e_{1}^{k-2}\right| .
\end{aligned}
$$

By combining these two estimates we get

$$
\max _{\bar{\Omega}}\left|e_{1}^{k}\right| \leq \max _{\partial O_{1}} w_{2} \max _{\bar{\Omega}}\left|e_{2}^{k-1}\right| \leq \max _{\partial O_{1}} w_{2} \max _{\partial O_{2}} w_{1} \max _{\bar{\Omega}}\left|e_{1}^{k-2}\right|=\rho_{\text {AltMR }} \max _{\bar{\Omega}}\left|e_{1}^{k-2}\right|
$$

Using this inequality recursively leads to $\max _{\bar{\Omega}}\left|e_{1}^{2 k}\right| \leq \rho_{\text {AltMR }}^{k} \max _{\bar{\Omega}}\left|e_{1}^{0}\right|$, which is the claim for $j=1$. The same arguments allow us to prove this result also for $j=2$.

Notice that Theorem 4 and Theorem 8 are proved using maximum principle arguments that are classical techniques for proving convergence of classical Schwarz methods; see, e.g., $[36,8]$. Notice also that the estimated contraction factors for AltMR and PMR are one the square of the other. This resembles classical results for Schwarz methods; see, e.g., $[14,15]$ and references therein.

## 5. Methods of Reflections with relaxation

In the following subsections, we present further forms of MR, and prove that they are equivalent to classical stationary iterative methods with relaxation. In particular, in Section 5.1, we present an averaged version of the parallel method introduced in [34] and show its equivalence with the damped block Jacobi iterations. In Section 5.2, we introduce a new MR that corresponds to a block variant of the well-known method of successive over relaxation (SOR).
5.1. The averaged parallel method of reflection. Golusin already said in [21] that the parallel method of reflections always converges for $J=2$ objects, but not for $J>2$, "but by modifying (17) (or (42)-(43)) appropriately one could increase the area of applicability of the preceding result". This fact has been mentioned in several publications; see, e.g., [21, 27]. In order to improve the convergence behavior of the PMR, Laurent et al. have recently proposed in [34] a modified version that is obtained (as mentioned by the authors) by averaging the different components $d_{j}^{k}$. To do so, the problem (36) is modified by adding a weight,

$$
\begin{align*}
\Delta d_{j}^{k} & =0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{k}=0 \text { on } \partial \Omega \\
d_{j}^{k} & =\left(1-\frac{1}{J}\right) d_{j}^{k-1}-\frac{1}{J} \sum_{\ell=1, \ell \neq j}^{J} d_{\ell}^{k-1} \text { on } \partial O_{j} \tag{50}
\end{align*}
$$

for $j=1, \ldots, J$, with the initialization problems

$$
\begin{equation*}
\Delta d_{j}^{1}=0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{1}=0 \text { on } \partial \Omega, d_{j}^{1}=g_{j}-u^{0} \text { on } \partial O_{j} \tag{51}
\end{equation*}
$$

for $j=1, \ldots, J$ and $u^{0} \in H$. The approximate solution $u^{k}$ can be obtained from $u^{k-1}$ by

$$
\begin{equation*}
u^{k}=u^{k-1}+\frac{1}{J} \sum_{j=1}^{J} d_{j}^{k} \tag{52}
\end{equation*}
$$

assuming that $u^{0}=0$. This new formulation of the method, that we call averaged parallel method of reflections (APMR), is proved to be always convergent in [34]. Now, we want to formulate this new version in terms of $v_{j}^{k}$. To do so, starting with the last boundary condition in (50) and using (52), we obtain on $\partial O_{j}$ that

$$
d_{j}^{k}=\left(1-\frac{1}{J}\right) d_{j}^{k-1}-\frac{1}{J} \sum_{\ell=1, \ell \neq j}^{J} d_{\ell}^{k-1}=d_{j}^{k-1}-\frac{1}{J} \sum_{\ell=1}^{J} d_{\ell}^{k-1}=d_{j}^{k-1}-\left(u^{k-1}-u^{k-2}\right)
$$

which implies that

$$
d_{j}^{k}+u^{k-1}=d_{j}^{k-1}+u^{k-2}=\cdots=d_{j}^{1}+u^{0}=g_{j}
$$

where we used that $d_{j}^{1}=g_{j}-u^{0}$ on $\partial O_{j}$. Hence, on $\partial O_{j}$ it holds that

$$
\begin{equation*}
d_{j}^{k}=g_{j}-u^{k-1} \tag{53}
\end{equation*}
$$

Now, we need to define the variable $v_{j}^{k}$ in a slightly different way than (12), namely

$$
\begin{equation*}
v_{j}^{k}:=v_{j}^{k-1}+\frac{1}{J} d_{j}^{k} \text { with } v_{j}^{0} \in H_{j} \tag{54}
\end{equation*}
$$

where $\sum_{j=1}^{J} v_{j}^{0}=u^{0}$. Then, we can use (52) to obtain, like for (13), the relation

$$
\begin{equation*}
u^{k}=u^{k-1}+\frac{1}{J} \sum_{j=1}^{J} d_{j}^{k}=u^{0}+\frac{1}{J} \sum_{n=1}^{k} \sum_{j=1}^{J} d_{j}^{n}=u^{0}+\sum_{n=1}^{k} \sum_{j=1}^{J}\left(v_{j}^{n}-v_{j}^{n-1}\right)=\sum_{j=1}^{J} v_{j}^{k} \tag{55}
\end{equation*}
$$

Now, using (53), (54), and (55), we get

$$
\begin{aligned}
v_{j}^{k} & =v_{j}^{k-1}+\frac{1}{J} d_{j}^{k}=v_{j}^{k-1}+\frac{1}{J} g_{j}-\frac{1}{J} u^{k-1}=v_{j}^{k-1}+\frac{1}{J} g_{j}-\frac{1}{J} \sum_{\ell=1}^{J} v_{\ell}^{k-1} \\
& =\left(1-\frac{1}{J}\right) v_{j}^{k-1}+\frac{1}{J} g_{j}-\frac{1}{J} \sum_{\ell=1, \ell \neq j}^{J} v_{\ell}^{k-1}
\end{aligned}
$$

which holds on $\partial O_{j}$. Recalling that $v_{j}^{k} \in H_{j}$, it solves the problem

$$
\begin{align*}
\Delta v_{j}^{k} & =0 \text { in } \Omega \backslash \partial O_{j}, v_{j}^{k}=0 \text { on } \partial \Omega \\
v_{j}^{k} & =\left(1-\frac{1}{J}\right) v_{j}^{k-1}+\frac{1}{J} g_{j}-\frac{1}{J} \sum_{\ell=1, \ell \neq j}^{J} v_{\ell}^{k-1} \text { on } \partial O_{j} \tag{56}
\end{align*}
$$

The equivalence between (56) and (50)-(51) is proved in the following theorem.
Theorem 9. Consider the sequences $\left\{u^{k}\right\}_{k \in \mathbb{N}} \subset H,\left\{d_{j}^{k}\right\}_{k \in \mathbb{N}^{+}}$with $d_{j}^{k} \in H_{j}$, and $\left\{v_{j}^{k}\right\}_{k \in \mathbb{N}}$ with $v_{j}^{k} \in H_{j}$ for $j=1, \ldots, J$, and assume that $v_{j}^{k}=v_{j}^{k-1}+d_{j}^{k}$ for $k \geq 1$ and $j=1, \ldots, J$. Then the following statements are equivalent:
(a) the $d_{j}^{k}$ solve (50)-(51) with $u^{k}=u^{k-1}+\sum_{j=1}^{J} d_{j}^{n}=\sum_{n=1}^{k} \sum_{j=1}^{J} d_{j}^{n}$.
(b) the $v_{j}^{k}$ solve (56) with $u^{k}=\sum_{j=1}^{J} v_{j}^{k}$.

Proof. The implication $(\mathrm{a}) \Rightarrow(\mathrm{b})$ is proved in the discussion before the theorem, whereas (b) $\Rightarrow$ (a) follows by similar arguments as in Theorem 1.

The block Jacobi method (42) we have presented earlier might not converge if the off-diagonal blocks are too heavy, i.e. the objects are too close (in agreement with the results of Golusin [21] and Traytak [44]); see also [10, Section 3] for detailed one-dimensional examples. However, one can consider a relaxation,

$$
\begin{equation*}
\widetilde{\mathbf{g}}^{k}=(1-\omega) \widetilde{\mathbf{g}}^{k-1}+\omega D^{-1}\left[-(L+U) \widetilde{\mathbf{g}}^{k-1}+\mathbf{g}\right] \tag{57}
\end{equation*}
$$

where the parameter $\omega$ has to be chosen in a proper way. This is the relaxed block Jacobi method, and we have the following result.

Theorem 10. We have the following equivalences:

- The relaxed block Jacobi method (57) is equivalent to the damped block Jacobi method, that is

$$
\begin{equation*}
\widetilde{\mathbf{g}}^{k}=\widetilde{\mathbf{g}}^{k-1}+\omega D^{-1}\left[\boldsymbol{g}-A \widetilde{\mathbf{g}}^{k-1}\right] \tag{58}
\end{equation*}
$$

- If one considers a sequence of relaxation parameters $\left\{\omega_{k}\right\}_{k}$, then the damped block Jacobi method is equivalent to Richardson's method [41], namely

$$
\widetilde{\mathbf{g}}^{k}=\widetilde{\mathbf{g}}^{k-1}+\omega_{k}\left[\boldsymbol{g}-A \widetilde{\mathbf{g}}^{k-1}\right] .
$$

- Assume that $v^{0} \in H$ and $\widetilde{g}_{j}^{0}=\tau_{j} v^{0}$ and $\omega=\frac{1}{J}$. Then the APMR (56) is equivalent to (58).
Proof. To show the first statement, we compute

$$
\begin{aligned}
\widetilde{\mathbf{g}}^{k} & =(1-\omega) \widetilde{\mathbf{g}}^{k-1}+\omega D^{-1}\left[-(L+U) \widetilde{\mathbf{g}}^{k-1}+\boldsymbol{g}\right] \\
& =\widetilde{\mathbf{g}}^{k-1}-\omega \widetilde{\mathbf{g}}^{k-1}+\omega D^{-1} \boldsymbol{g}-\omega D^{-1}(L+U) \widetilde{\mathbf{g}}^{k-1} \\
& =\widetilde{\mathbf{g}}^{k-1}-\omega \widetilde{\mathbf{g}}^{k-1}+\omega D^{-1} \boldsymbol{g}-\omega D^{-1}(A-D) \widetilde{\mathbf{g}}^{k-1} \\
& =\widetilde{\mathbf{g}}^{k-1}+\omega D^{-1}\left[\boldsymbol{g}-A \widetilde{\mathbf{g}}^{k-1}\right]
\end{aligned}
$$

which is the damped block Jacobi iteration (58). The second statement follows easily by recalling that $D=I$ and considering the parameter $\omega$ depending on the iteration $k$. Setting $\omega=\frac{1}{J}$ and recalling the operators $D, L$, and $U$ given in Section 3.3 , we can write (57) in the pointwise form

$$
\widetilde{g}_{j}^{k}=\left(1-\frac{1}{J}\right) \widetilde{g}_{j}^{k-1}+\frac{1}{J}\left[g_{j}-\sum_{\ell=1, \ell \neq j}^{J} \tau_{j} G_{\ell}\left(\widetilde{g}_{\ell}^{k-1}\right)\right]
$$

Then similar arguments as in Theorem 2 imply that this is equivalent to the APMR in (56).

We have seen that the APMR can be regarded as a relaxed (or damped) block Jacobi method or as Richardson's method with the specific choice $\omega_{k}=\omega=\frac{1}{J}$. This choice, motivated in [34] by a Hilbert projection analysis, is not guaranteed to be optimal. In what follows we derive under suitable conditions, an explicit formula for the optimal parameter $\omega^{\star}$ in discrete settings. Denoting by $G_{J}(\omega)$ the damped-block Jacobi iteration matrix obtained, for example, by a boundaryelement discretization of (18) and (57), the optimal parameter $\omega^{\star}$ is then given by

$$
\omega^{\star}=\underset{\omega \in[0,1]}{\arg \min } \rho\left(G_{\mathrm{J}}(\omega)\right)
$$

where $\rho\left(G_{\mathrm{J}}(\omega)\right)$ is the spectral radius of $G_{\mathrm{J}}(\omega)$. We have the following results.
Theorem 11. Let $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)$, for $k=1,2, \ldots$, be the eigenvalues of the block Jacobi iteration matrix. Assume that $\lambda_{k}\left(G_{\mathrm{J}}(1)\right) \in \mathbb{R}$ for any $k$ and that $\rho\left(G_{\mathrm{J}}(1)\right)<1$, that is the block Jacobi method converges. Then

$$
\omega^{\star}=\frac{2}{2-\left(\lambda_{\max }+\lambda_{\min }\right)}
$$

where $\lambda_{\max }$ and $\lambda_{\min }$ are the maximum and minimum eigenvalues of $G_{\mathrm{J}}(1)$.
Proof. Recalling (57), we observe that the damped block Jacobi iteration matrix has the form $G_{\mathrm{J}}(\omega)=(1-\omega) I+\omega G_{\mathrm{J}}(1)$. Hence its eigenvalues are

$$
\lambda_{k}(\omega)=1-\omega\left(1-\lambda_{k}\left(G_{\mathrm{J}}(1)\right)\right)
$$

for $k=1,2, \ldots$ Now, the optimal parameter $\omega^{\star}$ solves the problem

$$
\min _{\omega \in[0,1]} \max _{k}\left|1-\omega\left(1-\lambda_{k}\left(G_{\mathrm{J}}(1)\right)\right)\right| .
$$

Since $\rho\left(G_{\mathrm{J}}(1)\right)<1$, then $\left.\lambda_{k}\left(G_{\mathrm{J}}(1)\right)\right) \in(0,1)$. Therefore, the spectral radius $\rho\left(G_{\mathrm{J}}(\omega)\right)$ has the form

$$
\rho\left(G_{\mathrm{J}}(\omega)\right)=\max \left\{\left|1-\omega\left(1-\lambda_{\max }\right)\right|,\left|1-\omega\left(1-\lambda_{\min }\right)\right|\right\}
$$

and the optimal parameter $\omega^{\star}$ is the point where the two straight functions $\mid 1-$ $\omega\left(1-\lambda_{\max }\right) \mid$ and $\left|1-\omega\left(1-\lambda_{\min }\right)\right|$ intersect, that is $\omega=\frac{2}{2-\left(\lambda_{\max }+\lambda_{\min }\right)}$, which is our claim.

Theorem 11 provides the optimal parameter under the assumption that the block Jacobi method converges. If block Jacobi does however not converge, can $\omega$ be used to make it convergent? In other words, if an eigenvalue of $G_{\mathrm{J}}(1)$ is such that $\left|\lambda_{k}\left(G_{\mathrm{J}}(1)\right)\right| \geq 1$, is there a choice of $\omega$ capable to correct this behavior? To study this unfavorable case, we have to distinguish different cases:

- if $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)=1$, then $\lambda_{k}(\omega)=1-\omega\left(1-\lambda_{k}\left(G_{\mathrm{J}}(1)\right)\right)=1$ for any $\omega$. Hence, there exists no $\omega \in \mathbb{R}$ such that $\left|\lambda_{k}(\omega)\right|<1$.
- if $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)=-1$, then $\lambda_{k}(\omega)=1-2 \omega$, which implies that $\left|\lambda_{k}(\omega)\right|<1$ if and only if $0<\omega<1$.
- if $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)<-1$, then

$$
\left|\lambda_{k}(\omega)\right|=\left|1-\omega\left(1-\lambda_{k}\left(G_{\mathrm{J}}(1)\right)\right)\right|<1 \Leftrightarrow 0<\omega<\frac{2}{1-\lambda_{k}\left(G_{\mathrm{J}}(1)\right)}
$$

- if $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)>1$, then

$$
\left|\lambda_{k}(\omega)\right|=\left|1-\omega\left(1-\lambda_{k}\left(G_{\mathrm{J}}(1)\right)\right)\right|<1 \Leftrightarrow \frac{2}{1-\lambda_{k}\left(G_{\mathrm{J}}(1)\right)}<\omega<0
$$

We can summarize these facts in the following result.
Theorem 12. Let $G_{\mathrm{J}}(1)$ and $G_{\mathrm{J}}(\omega)$ be the block Jacobi and damped block Jacobi iteration matrices and assume that $\left|\lambda_{k}\left(G_{\mathrm{J}}(1)\right)\right| \geq 1$ for a given $k$. We have that

- if $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)=-1$, then $\left|\lambda_{k}(\omega)\right|<1$ for any $\omega \in(0,1)$.
- if $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)=1$, then $\left|\lambda_{k}(\omega)\right|=1$ for any $\omega \in \mathbb{R}$.
- if $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)<0$, then $\left|\lambda_{k}(\omega)\right|<1$ if and only if $0<\omega<\frac{2}{1-\lambda_{k}\left(G_{\mathrm{J}}(1)\right)}$.
- if $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)>1$, then $\left|\lambda_{k}(\omega)\right|<1$ if and only if $\frac{2}{1-\lambda_{k}\left(G_{\mathrm{J}}(1)\right)}<\omega<0$.

Theorem 12 shows that if block Jacobi does not converge, it is not always possible to find a parameter $\omega \in(0,1]$ that makes the method convergent. If the eigenvalues $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)$ are such that $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)<1$, i.e. there can be arbitrarily large negative eigenvalues, a case that we have observed numerically, then a similar proof as for Theorem 11 allows us to obtain the following result.

Theorem 13. Let $G_{\mathrm{J}}(1)$ and $G_{\mathrm{J}}(\omega)$ be the block Jacobi and damped-block Jacobi iteration matrices and assume that the eigenvalues $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)$ are real and such that $\lambda_{k}\left(G_{\mathrm{J}}(1)\right)<1$, then the optimal parameter is $\omega^{\star}=\frac{2}{2-\left(\lambda_{\max }+\lambda_{\min }\right)}$, where $\lambda_{\max }$ and $\lambda_{\min }$ are the maximum and minimum eigenvalues of $G_{\mathrm{J}}(1)$.
5.2. Successive-Over-Relaxation method (SOR). In the previous sections, we have seen that the different forms of the method of reflections are equivalent to block Gauss-Seidel, block Jacobi, damped block Jacobi, and Richardson's method. This suggests that a possible improvement of the method of reflections would be to
rewrite the AltMR as an SOR method, which converges in general much faster for a well chosen relaxation parameter [45]. SOR for the system (18) is

$$
\begin{equation*}
\widetilde{\mathbf{g}}^{k}=(D+\omega L)^{-1}\left[\omega \mathbf{g}+(-\omega U+(1-\omega) D) \widetilde{\mathbf{g}}^{k-1}\right] \tag{59}
\end{equation*}
$$

Recalling the form of $D, L$, and $U$ given in Section 3.3, (59) written component-wise is

$$
\begin{equation*}
\widetilde{g}_{j}^{k}=(1-\omega) \widetilde{g}_{j}^{k-1}+\omega\left[g_{j}-\sum_{\ell=1}^{j-1} \tau_{j} G_{\ell}\left(\widetilde{g}_{\ell}^{k}\right)-\sum_{\ell=j+1}^{J} \tau_{j} G_{\ell}\left(\widetilde{g}_{\ell}^{k-1}\right)\right] \tag{60}
\end{equation*}
$$

Defining $v_{j}^{k}:=G_{j}\left(\widetilde{g}_{j}^{k}\right)$, this iteration can be rewritten in the volume form,

$$
\begin{align*}
\Delta v_{j}^{k} & =0 \text { in } \Omega \backslash \partial O_{j}, v_{j}^{k}=0 \text { on } \partial \Omega \\
v_{j}^{k} & =(1-\omega) v_{j}^{k-1}+\omega\left[g_{j}-\sum_{\ell=1}^{j-1} v_{\ell}^{k}-\sum_{\ell=j+1}^{J} v_{\ell}^{k-1}\right] \text { on } \partial O_{j} \tag{61}
\end{align*}
$$

with $u^{k}=\sum_{j=1}^{J} v_{j}^{k}$. This is a new formulation of the method of reflections that can be also written in terms of the (more usual) $d_{j}^{k}$, that is

$$
\begin{align*}
\Delta d_{j}^{k} & =0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{k}=0 \text { on } \partial \Omega, \\
d_{j}^{k} & =(1-\omega) d_{j}^{k-1}-\omega\left[\sum_{\ell=1}^{j-1} d_{\ell}^{k}+\sum_{\ell=j+1}^{J} d_{\ell}^{k-1}\right] \text { on } \partial O_{j}, \tag{62}
\end{align*}
$$

with

$$
\begin{align*}
\Delta d_{j}^{1} & =0 \text { in } \Omega \backslash \partial O_{j}, d_{j}^{1}=0 \text { on } \partial \Omega \\
d_{j}^{1} & =g_{j}-u^{0}-\omega \sum_{\ell=1}^{j-1} d_{\ell}^{1} \text { on } \partial O_{j} \tag{63}
\end{align*}
$$

and $u^{k}=u^{k-1}+\omega \sum_{j=1}^{J} d_{j}^{k}=u^{0}+\omega \sum_{n=1}^{k} \sum_{j=1}^{J} d_{j}^{n}$. Notice that (62)-(63) and (60) are equivalent, because (60) is the equivalent substructured form of (61) and the equivalence between (61) and (62)-(63) is proved in the following theorem.
Theorem 14. Consider the sequences $\left\{u^{k}\right\}_{k} \subset H$ with $u^{0} \in H,\left\{d_{j}^{k}\right\}_{k}$ with $d_{j}^{k} \in$ $H_{j}$, and $\left\{v_{j}^{k}\right\}_{k}$ with $v_{j}^{k} \in H_{j}$ for $j=1, \ldots, J$, and assume that $v_{j}^{k}=v_{j}^{k-1}+\omega d_{j}^{k}$ for $k \geq 1$ and $\omega>0$ and $v_{j}^{0} \in H_{j}$ for $j=1, \ldots, J$. Then the following statements are equivalent:
(a) the $d_{j}^{k}, j=1, \ldots, J$, solve (62)-(63) and $u^{k}=u^{k-1}+\omega \sum_{j=1}^{J} d_{j}^{n}=\omega \sum_{n=1}^{k} \sum_{j=1}^{J} d_{j}^{n}$.
(b) the $v_{j}^{k}, j=1, \ldots, J$, solve (61) and $u^{k}=\sum_{j=1}^{J} v_{j}^{k}$.

Proof. (a) $\Rightarrow$ (b): first, using $v_{j}^{k}=v_{j}^{k-1}+\omega d_{j}^{k}$, we write that
(64) $u^{k}=u^{k-1}+\omega \sum_{j=1}^{J} d_{j}^{k}=u^{0}+\omega \sum_{n=1}^{k} \sum_{j=1}^{J} d_{j}^{n}=u^{0}+\sum_{n=1}^{k} \sum_{j=1}^{J}\left(v_{j}^{n}-v_{j}^{n-1}\right)=\sum_{j=1}^{J} v_{j}^{k}$.

Notice that $v_{j}^{k} \in H_{j}$, hence it is sufficient to prove that the boundary condition on $\partial O_{j}$ in (61) holds. Using the transmission condition of (62) and $u^{k-1}=u^{k-2}+$
$\omega \sum_{\ell=1}^{J} d_{\ell}^{k-1}$, we obtain

$$
\begin{aligned}
d_{j}^{k} & =d_{j}^{k-1}-\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k}-\omega \sum_{\ell=j}^{J} d_{\ell}^{k-1}=d_{j}^{k-1}-\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k}-\omega \sum_{\ell=1}^{J} d_{\ell}^{k-1}+\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k-1} \\
& =d_{j}^{k-1}-\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k}-\left(u^{k-1}-u^{k-2}\right)+\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k-1}
\end{aligned}
$$

which implies that

$$
d_{j}^{k}+u^{k-1}+\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k}=d_{j}^{k-1}+u^{k-2}+\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k-1}=\cdots=d_{j}^{1}+u^{0}+\omega \sum_{\ell=1}^{j-1} d_{\ell}^{1}=g_{j}
$$

where we used the transmission condition of (63). Hence, we have obtained that $d_{j}^{k}=g_{j}-u^{k-1}-\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k}$. Replacing this into the transmission condition of (62) and using $v_{j}^{k}=v_{j}^{k-1}+\omega d_{j}^{k}$ and (64), we obtain

$$
\begin{aligned}
\frac{1}{\omega}\left(v_{j}^{k}-v_{j}^{k-1}\right) & =d_{j}^{k}=d_{j}^{k-1}-\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k}-\omega \sum_{\ell=j}^{J} d_{\ell}^{k-1} \\
& =g_{j}-u^{k-2}-\omega \sum_{\ell=1}^{J} d_{\ell}^{k-1}-\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k}=g_{j}-u^{k-1}-\omega \sum_{\ell=1}^{j-1} d_{\ell}^{k} \\
& =g_{j}-u^{k-1}-\sum_{\ell=1}^{j-1}\left(v_{\ell}^{k}-v_{\ell}^{k-1}\right)=g_{j}-\sum_{\ell=1}^{j-1} v_{\ell}^{k}-\sum_{\ell=j}^{J} v_{\ell}^{k-1}
\end{aligned}
$$

which implies that

$$
v_{j}^{k}=v_{j}^{k-1}+\omega g_{j}-\omega \sum_{\ell=1}^{j-1} v_{\ell}^{k}-\omega \sum_{\ell=j}^{J} v_{\ell}^{k-1}=(1-\omega) v_{j}^{k-1}+\omega\left[g_{j}-\sum_{\ell=1}^{j-1} v_{\ell}^{k}-\sum_{\ell=j+1}^{J} v_{\ell}^{k-1}\right]
$$

which is our claim.
$(\mathrm{b}) \Rightarrow(\mathrm{a})$ : using $d_{j}^{k}=\frac{1}{\omega}\left(v_{j}^{k}-v_{j}^{k-1}\right)$, we obtain

$$
u^{k}=\sum_{j=1}^{J} v_{j}^{k}=u^{0}+\sum_{j=1}^{J}\left(v_{j}^{k}-v_{j}^{0}\right)=u^{0}+\sum_{n=1}^{k} \sum_{j=1}^{J}\left(v_{j}^{n}-v_{j}^{n-1}\right)=u^{0}+\omega \sum_{n=1}^{k} \sum_{j=1}^{J} d_{j}^{n}
$$

which implies that $u^{k}=u^{k-1}+\omega \sum_{j=1}^{J} d_{j}^{k}$. Now, since $d_{j}^{k}=\frac{1}{\omega}\left(v_{j}^{k}-v_{j}^{k-1}\right)$ we have that $d_{j}^{k} \in H_{j}$. It remains to show that the transmission conditions of (62)-(63) hold. For $k=1$ we have

$$
\begin{aligned}
d_{j}^{1} & =\frac{1}{\omega}\left(v_{j}^{1}-v_{j}^{0}\right)=g_{j}-\sum_{\ell=1}^{j-1} v_{\ell}^{1}-\sum_{\ell=j}^{J} v_{\ell}^{0} \\
& =g_{j}-\omega \sum_{\ell=1}^{j-1} d_{\ell}^{1}-\sum_{\ell=1}^{j-1} v_{\ell}^{0}-\sum_{\ell=j}^{J} v_{\ell}^{0}=g_{j}-\omega \sum_{\ell=1}^{j-1} d_{\ell}^{1}-u^{0}
\end{aligned}
$$

In the case $k>1$, we write the transmission condition of (61) on $\partial O_{j}$ for $k$ and $k-1$ :

$$
v_{j}^{k}=(1-\omega) v_{j}^{k-1}+\omega\left[g_{j}-\sum_{\ell=1}^{j-1} v_{\ell}^{k}-\sum_{\ell=j+1}^{J} v_{\ell}^{k-1}\right]
$$

$$
v_{j}^{k-1}=(1-\omega) v_{j}^{k-2}+\omega\left[g_{j}-\sum_{\ell=1}^{j-1} v_{\ell}^{k-1}-\sum_{\ell=j+1}^{J} v_{\ell}^{k-2}\right] .
$$

Subtracting term by term these equations, dividing by $\omega$, and recalling that $d_{j}^{k}=$ $\frac{1}{\omega}\left(v_{j}^{k}-v_{j}^{k-1}\right)$ we obtain the transmission condition of (62).

A good choice of the parameter $\omega$ is not trivial. In general one would consider the optimal parameter $\omega^{\star}$, that is the parameter that minimizes the contraction rate of the SOR method. From a discrete point of view, namely when one is interested in solving a linear system $A \widetilde{\mathbf{g}}=\mathbf{g}$ using a stationary iterative method of the form $\widetilde{\mathbf{g}}^{k+1}=M^{-1} N \widetilde{\mathbf{g}}^{k}+M^{-1} \mathbf{g}$, there are several results in the literature for point relaxation (in contrast to block relaxation here). A famous result proved by Kahan [29] says that for the SOR method one has $\rho\left(M^{-1} N\right) \geq|1-\omega|$, which implies that a necessary condition for convergence is $\omega \in(0,2)$. In the case the matrix $A$ has the so called Property A, David Young proved in [45] that the optimal parameter, that is the parameter that minimizes the spectral radius $\rho\left(G_{\mathrm{SOR}}(\omega)\right)$ with respect to $\omega$, where $G_{\mathrm{SOR}}(\omega)$ is the SOR -iteration matrix, is given by $\omega^{\star}=\frac{2}{1+\sqrt{1-\rho\left(G_{\mathrm{J}}\right)^{2}}}$, where $G_{\mathrm{J}}$ is the iteration matrix of the corresponding point-Jacobi method applied to the same linear system. A discrete form of our problem (18) for $J=2$ is characterized by a matrix $A$ of the form

$$
A=\left[\begin{array}{cc}
I_{1} & G_{2} \\
G_{1} & I_{2}
\end{array}\right]
$$

where $I_{1}$ and $I_{2}$ are identity matrices. This shows that $A$ is exactly in the form required to have the Property A. Therefore, we can apply the theory developed by David Young to get an explicit formula for the optimal parameter. This observation is not in general true for $J>2$. Moreover, to compute $\omega^{\star}$ one would need the spectral radius of the block Jacobi method (the PMR). This is not an easy task, even for $J=2$. However, we can estimate $\rho\left(G_{\mathrm{J}}\right)$ at a continuous level using maximum principle arguments, as we show in Section 3.5. Numerical experiments in Section 7 show that this leads to a very good estimate of $\omega^{\star}$ for $J=2$. In the case that $J>2$ the optimal parameter depends strongly on the geometry and we observed that when the distance between the objects is sufficiently large the optimal choice is $\omega^{\star} \approx 1$.

## 6. ScALABILITY ANALYSIS AND COARSE CORRECTION

An iterative method is said to be scalable, if its contraction factor does not deteriorate when the number of unknowns grows; see, e.g., [43] and references therein. In the field of domain-decomposition methods for the solution of partial differential equations, we distinguish different scalability problems depending on the way the number of unknowns increases:
(a) the number of discretization points increases, but number and size of subdomains is unchanged.
(b) the subdomains have fixed size and their number increases making also the size of the overall domain increase.
The first property is widely studied in the literature, see, e.g., [43] and references therein. The case (b) has been only recently observed in [5] and theoretically investigated in $[7,8,9]$ for classical Schwarz methods and for classical one-level


Figure 4. Geometry of the one-dimensional example used to study scalability of the AltMR. In this figure, $\delta$ is the distance between the centers of the $J$ equidistant objects that are represented by the thick segments.


Figure 5. Left: Spectral radius of the AltMR as a function of the number of objects. Right: Example of the spectrum of the (onedimensional) operator $G_{\mathrm{GS}}=-(D+L)^{-1} U$. Notice that many eigenvalues are clustered around 1.
domain-decomposition methods, like Neumann-Neumann, Dirichlet-Neumann and optimized-Schwarz methods in [6].

Studying methods of reflections, we observed numerically that scalability in the case (a) holds, that is the convergence of the method is not affected by the number of discretization points of the boundaries of the objects $\partial O_{j}$, for $j=1, \ldots, J$. On the other hand, scalability in the sense (b) does not hold, as we will show below. Moreover, if the domain $\Omega$ is fixed and one increases the number of objects, then obviously the distance between them reduces, and we have seen in Section 9 for $J=2$ that the contraction factor deteriorates in this case. The same behavior has been observed, e.g., in [34, 10]. This implies the need of a coarse-correction to address these unfavorable behaviors.

We consider a one-dimensional problem characterized by $J$ equidistant objects (subintervals); see Figure 4. This is a finite-dimensional problem because the unknowns are the Dirichlet data $\widetilde{\mathbf{g}} \in \mathbb{R}^{2 J}$ on the $2 J$ extrema of the objects. The problem can then be written in a form $A \widetilde{\mathbf{g}}=\mathbf{g}$, where $\mathbf{g} \in \mathbb{R}^{2 J}$ and $A \in \mathbb{R}^{2 J \times 2 J}$. Using the equivalence result proved in Theorem 2 , one can easily construct the block Gauss-Seidel (AltMR) iteration matrix, that is $G_{\mathrm{GS}}=-(D+L)^{-1} U$ and numerically compute its spectral radius for increasing number of objects $J$. The result is given in Figure 5, which clearly shows that the AltMR does not scale because $\rho\left(G_{\mathrm{GS}}\right)$ deteriorates for growing $J$. The heuristic reason for this behavior is that, once the approximate solution is corrected on one object, say the $j$ th one, then the procedure continues to correct it sequentially on all the other objects. Every correction induces an additional error in the approximate solution on the $j$ th object. Therefore, when $J$ increases, more and more objects have to be corrected, which adds more and more error induced by the alternating correction procedure on the


Figure 6. Example of an eigenvector of the one-dimensional iteration matrix $G_{\mathrm{GS}}$ corresponding to an eigenvalue $\lambda \approx 1$. The right plot shows a zoom into the highly oscillating area of the eigenvector depicted in the left plot.
$j$ th object. Similar arguments apply also for the other methods of reflections. We thus need a coarse correction to obtain a scalable method.

A spectral analysis of the one-dimensional iteration matrix $G_{G S}$ reveals that many of its eigenvalues are clustered around 1, see Figure 5 (right), and that the corresponding eigenvectors have a special structure: they globally oscillate with (local) maxima and minima attained on the objects, and many of them have constant value on the objects; see Figure 6. This suggests that a coarse-correction has mainly to deal with errors on the boundaries of the objects. Therefore, to design our coarse-space we denote by $\psi_{j, n}$ for $n=0,1, \ldots$ the eigenfunction of the Laplace-Beltrami operator on the $j$ th object $\partial O_{j}$, for $j=1, \ldots, J$, and introduce functions $\varphi_{j, n}$ as solution to the Dirichlet problem

$$
\begin{align*}
\Delta \varphi_{j, n} & =0 \text { in } \Omega \backslash \cup_{\ell=1}^{J} \partial O_{\ell} \\
\varphi_{j, n} & =\psi_{j, n} \text { on } \partial O_{j}  \tag{65}\\
\varphi_{j, n} & =0 \text { on } \partial \Omega \cup\left(\cup_{\ell=1, \ell \neq j}^{J} \partial O_{\ell}\right),
\end{align*}
$$

for $j=1, \ldots, J$ and $n=0,1, \ldots$ Notice that $\varphi_{j, n} \in H$ for any $j$ and $n$, and the functions $\psi_{j, n}$ are classical Fourier-basis functions in two dimensions and spherical harmonics in three dimensions. Our coarse-space is then defined as

$$
V_{c}^{N}:=\operatorname{span}_{j=1, \ldots, J, n=0,1, \ldots, N}\left\{\varphi_{j, n}\right\},
$$

where $N$ is the number of eigenfunctions considered. Notice that the dimension of $V_{c}^{N}$ is proportional to $N$. It is clear that the coarse space $V_{c}^{N}$ has mainly information condensed on the boundaries of the objects, similar to the Spectral Harmonically Enriched Multiscale coarse space SHEM in domain decomposition [19, 18], which contains mainly information on the interfaces between subdomains, see also $[17,16]$. It is important to remark that the construction of each function $\varphi_{j, n}$ would require the solution of problem (65), which requires the same computational effort of the original problem (3). However, this is not needed because the substructured formulation introduced in Section 3.3 allows us to work directly on the boundary of the objects: our coarse space contains thus spectral approximations of the substructured problem. For this reason, we introduce the restriction
operator and its corresponding adjoint (prolongation or extension) operator as

$$
R:=\left[\begin{array}{llll}
\mathcal{F}_{1}^{N} & & & \\
& \mathcal{F}_{2}^{N} & & \\
& & \ddots & \\
& & & \mathcal{F}_{J}^{N}
\end{array}\right] \quad \text { and } \quad R^{*}=\left[\begin{array}{llll}
\left(\mathcal{F}_{1}^{N}\right)^{*} & & & \\
& \left(\mathcal{F}_{2}^{N}\right)^{*} & & \\
& & \ddots & \\
& & & \left(\mathcal{F}_{J}^{N}\right)^{*}
\end{array}\right]
$$

where $\mathcal{F}_{j}^{N}: L^{2}\left(\partial O_{j}\right) \rightarrow \mathbb{R}^{N}$ is given by

$$
\mathcal{F}_{j}^{K}(w)=\left[\begin{array}{lll}
c_{1, j} & \ldots & c_{N, j}
\end{array}\right]^{\top}=: \mathbf{c}_{j}
$$

with $c_{k, j}=\int_{\partial O_{j}} w f_{k, j}$ for any $w \in L^{2}\left(\partial O_{j}\right)$ and $j=1, \ldots, J$, and $\left(\mathcal{F}_{j}^{N}\right)^{*}: \mathbb{R}^{K} \rightarrow$ $L^{2}\left(\partial O_{j}\right)$ is

$$
\left(\mathcal{F}_{j}^{N}\right)^{*}\left(\mathbf{c}_{j}\right)=\sum_{k=1}^{K} c_{k, j} \psi_{k, j}
$$

for $\mathbf{c}_{j} \in \mathbb{R}^{K}$. These operators allow us to restrict the operator $A$ introduced in (18) and (19) on the coarse space $V_{c}^{N}$ as $A_{c}=R A R^{*} \in \mathbb{R}^{J N \times J N}$.

We are now ready to state our two-level method of reflection: given an approximation $\widetilde{\mathbf{g}}^{k}$ to $\widetilde{\mathbf{g}}$ at the iteration $k$, a coarse-corrected method of reflections step is defined as

$$
\begin{array}{rlr}
\widetilde{\mathbf{g}}^{k+1 / 2} & :=G_{\mathrm{MR}} \widetilde{\mathbf{g}}^{k}+M_{\mathrm{MR}}^{-1} \mathbf{g} & \text { (one-level MR step), } \\
\widetilde{\mathbf{r}}^{k+1 / 2} & :=\mathbf{g}-A \widetilde{\mathbf{g}}^{k+1 / 2} & \text { (compute the residual) }, \\
\widetilde{\mathbf{g}}_{c}^{k+1 / 2} & :=A_{c}^{-1} R \widetilde{\mathbf{r}}^{k+1 / 2} & \text { (compute the correction), } \\
\widetilde{\mathbf{g}}^{k+1} & :=\widetilde{\mathbf{g}}^{k+1 / 2}+R^{*} \widetilde{\mathbf{g}}_{c}^{k+1 / 2} & \left(\text { correct } \widetilde{\mathbf{g}}^{k+1 / 2}\right. \text { ), }
\end{array}
$$

where $G_{\mathrm{MR}}$ is the one-level method of reflection operator and $M_{\mathrm{MR}}$ the corresponding preconditioning matrix, which can be, e.g., $G_{\mathrm{GS}}$ (with $M_{\mathrm{MR}}=(D+L)$ ) for the AltMR or $G_{\mathrm{J}}$ (with $M_{\mathrm{MR}}=D$ ) for the PMR. A direct calculation reveals that the two-level method of reflections iteration operator is

$$
G_{\mathrm{MR}-\text { c.c. }}=\left[I-R^{*} A_{c}^{-1} R A\right] G_{\mathrm{MR}}
$$

A coarse corrected method of reflection is given in Algorithm 5.

```
Algorithm 5 Method of Reflections with Coarse Correction
Input: \(K\) (maximum number of iterations), tol (tolerance), \(\widetilde{\mathbf{g}}^{0}\) (initial guess).
Input: \(f_{\mathrm{MR}}(\mathbf{v}, \mathbf{w}):=G_{\mathrm{MR}} \mathbf{v}+M_{\mathrm{MR}}^{-1} \mathbf{w}\) (function that performs one step of a method of
    reflection MR),
Input: g (data of the problem).
    Compute \(\widetilde{\mathbf{r}}^{0}=\mathbf{g}-A \widetilde{\mathbf{g}}^{0}\) and set \(k=0\).
    while \(k<K\) and \(\left\|\widetilde{\mathbf{r}}^{k}\right\|>\) tol do
        Set \(k=k+1\).
        \(\widetilde{\mathbf{g}}^{k}=f_{\mathrm{MR}}\left(\widetilde{\mathbf{g}}^{k-1}, \mathbf{g}\right)\).
        \(\widetilde{\mathbf{r}}^{k}=\mathbf{g}-A \widetilde{\mathbf{g}}^{k}\).
        \(\widetilde{\mathbf{g}}_{c}=A_{c}^{-1} R \widetilde{\mathbf{r}}^{k}\).
        \(\widetilde{\mathbf{g}}^{k+1}=\widetilde{\mathbf{g}}^{k}+R^{\top} \widetilde{\mathbf{g}}_{c}\).
    end while
```



Figure 7. Geometry of the problem under consideration, domain with some normal unit vectors. Left: objects relatively far from each other, right: objects relatively close to each other.

Notice that like in domain decomposition, Algorithm 5 requires the solution of the problem projected onto the coarse space. If the dimension of the coarse space is proportional to $J$, then this solution requires in principle $O\left(J^{3}\right)$ flops. Therefore, for very large $J$ the cost of the coarse correction becomes comparable to the cost of solving the full problem by a direct method. In this case, our two-level framework would need to be extended to a multi-level setting, which is however beyond the scope of this paper.

## 7. Numerical experiments

We now present some numerical tests to illustrate the results we obtained, and to explore cases not covered by the theory we developed. We start with the case of two objects, and then investigate the case of three objects. We then study the scalability of the method with respect to the number of objects and finally consider the method as a preconditioner for the GMRES algorithm. We used a publicly available package ${ }^{1}$ of Matlab functions to solve the integral equations of the problem at collocation points by the Nyström method.
7.1. Two objects. We first consider the case of two objects of radius $r=0.2$, in the unit disc, for which we have theoretical results. The theoretical value of the optimal $\omega^{\star}$ for the SOR variant of the method of reflections is given by

$$
\begin{equation*}
\omega^{\star}(r, \delta)=\frac{2}{1+\sqrt{1-\rho_{\mathrm{PMR}}(r, \delta)^{2}}} \tag{66}
\end{equation*}
$$

With this choice of $\omega$, the theoretical convergence factor is

$$
\begin{equation*}
\rho_{\mathrm{SOR}}(r, \delta)_{\omega=\omega^{\star}}=\omega^{\star}(r, \delta)-1 \tag{67}
\end{equation*}
$$

In a first experiment, we consider two cases where the distance between the objects is either relatively large, that is $\delta=0.5$, or relatively small that is $\delta=0.25$, see Figure 7. We consider both SOR and damped PMR. In the case of SOR, $\omega^{\star}$ is given by Equation (66), where, in our 2 object case, $\rho_{\mathrm{PMR}}(r, \delta)$ can be computed numerically and estimated theoretically. Alternatively the latter quantity can be estimated using formulas provided in Theorem 8 and Corollary 1. Results are presented in Figure 8. We see that the agreement between the theoretical and

[^1]


Figure 8. Contraction factor of the SOR variant of the method of reflections with respect to $\omega$ (stars), theoretical values of $\omega^{\star}$ given by (66) where $\rho_{\mathrm{PMR}}(r, \delta)$ is obtained either using Theorem 8 and Corollary 1 (dashed line), or computed numerically (dotted line). In the latter case, the optimal contraction factor is computed by Equation (67) (solid line). Left: the objects are far away from each other. Right: the objects are close to each other.


Figure 9. Contraction factor of the damped PMR with respect to $\omega$ (stars), theoretical estimate of the contraction factor (solid line) given by Theorem 8 and Corollary 1. Left: the objects are far from each other. Right: the objects are close to each other.
numerical values of $\omega^{\star}$ is good when the objects are far from each other, and less accurate when the objects are close to each other. We also see that the relaxation parameter $\omega^{\star}$ is only close to 1 in the first case, which means that the SOR variant provides a real improvement with respect to the standard AltMR when the objects are relatively close to each other. In the case of PMR, we observe that the damping does not improve the contraction factor, see Figure 9.

The theoretical estimate of the contraction factor is rather accurate and does not seem to depend on the distance between the objects. Note also that the eigenvalues of the PMR iteration matrix form a symmetric set with respect to zero. Indeed, one has

$$
G_{\mathrm{J}}=-D^{-1}(L+U)=-\left[\begin{array}{cc}
0 & \widetilde{G} \\
\widetilde{G} & 0
\end{array}\right]
$$

so that
$\operatorname{det}\left(G_{\mathrm{J}}-\lambda I\right)=\operatorname{det}\left[\begin{array}{cc}-\lambda I & -\widetilde{G} \\ -\widetilde{G} & -\lambda I\end{array}\right]=\operatorname{det}(-\lambda I) \operatorname{det}\left(-\lambda I+\frac{1}{\lambda} \widetilde{G}^{2}\right)=\operatorname{det}\left(-\lambda^{2} I+\widetilde{G}^{2}\right)$.


Figure 10. Contraction factors of AltMR (left) and PMR (right) with respect to $\delta$ and theoretical contraction factors (solid lines) given by Theorems 4 and 8 and Corollary 1, compared with numerical contraction factors (stars).


Figure 11. Left: geometry of the problem, domain with some normal unit vectors. Right: Contraction factor of the SOR variant with respect to $\omega$ (stars), theoretical values of $\omega^{\star}$ given by (66) where $\rho_{\text {PMR }}(r, \delta)$ is obtained either from Theorem 8 and Corollary 1 (dashed line) or computed numerically (dot line). In the latter case, the optimal contraction factor is computed by Equation (67) (solid line).

As a consequence, if $\lambda$ is an eigenvalue, then also $-\lambda$ is an eigenvalue. Because of Theorem 11, we find that in this case $\omega^{\star}=1$, as observed in Figure 9.

Finally, we compare theoretical and numerical values of the contraction factors both for AltMR and PMR when the distance $\delta$ between the objects varies; the results are shown in Figure 10, and we see a good agreement between theory and numerics.
7.2. Three objects. We now consider the case of three objects of radius $r=0.2$ in the unit disc, for which we do not have theoretical results. For the sake of simplicity, we restrict ourselves to SOR, and repeat the experiments done to obtain Figure 8. The results are presented in Figure 11. We observe again that SOR provides a real improvement compared to the standard AltMR, in the sense that the observed optimal $\omega$ differs from 1 . We however also see that the estimates obtained with Equation (66) are not accurate now.
7.3. Scalability with respect to the number of objects. In this test, we study the scalability of the AltMR and PMR when the number of objects varies, and the effect of the coarse correction on both methods. We consider for a given $L \in \mathbb{N}^{+}$



Figure 12. Geometry of the problem under consideration for $L=$ 3 (left) and $L=4$ (right).


Figure 13. Contraction factors of the AltMR (left) and PMR (right) with respect to the number of objects. Contraction factor without (crosses) and with (circles) coarse correction.
a square of size $L \times L$ with rounded corners containing $J=L^{2}$ objects of radius $r=0.1$, see Figure 12. We then evaluate numerically the contraction factor for both methods, with and without coarse correction, see Figure 13. We consider here the coarse space $V_{c}^{0}=\operatorname{span}_{j=1, \ldots, J}\left\{\varphi_{j, 0}\right\}$ of dimension $J$, where we use only the first (constant) Fourier mode $\psi_{j, 0}$ for each object, see (65). We observe that the coarse correction significantly reduces the contraction factor, by approximately one order of magnitude for the AltMR and even more for the PMR, where the coarse correction leads to a convergence factor below 1, so that the PMR with coarse correction remains convergent even in cases where the standard PMR diverges.

We next repeat these two tests with the method of reflection variants with relaxation, namely, with the SOR variant and the APMR, see Figure 14. We see that the coarse correction significantly improves the performance of the SOR variant, but not the APMR. It even seems that the convergence rate remains unchanged in the case of APMR. This is due to the relaxation parameter $\omega=\frac{1}{J}$ that decays as $J$ grows.

Next, we solve the problem using GMRES and employing the methods of reflections as preconditioners, both with and without coarse correction. For reference, we compare the performance of these methods with the coarse-corrected methods of reflections used as stationary iterative methods; the required iteration numbers are summarized in Table 1. We observe that with GMRES now the PMR preconditioner leads to a convergent method, even though as a stationary iterative method it was not convergent. This is very similar to the additive Schwarz method which is


Figure 14. Contraction factors of the SOR variant with $\omega=1.3$ (left) and APMR (right) with respect to the number of objects. Contraction factor without (crosses) and with (circles) coarse correction.

| \# Objects $(r=0.1)$ | 4 | 9 | 16 | 25 | 36 | 49 | 64 | 81 | 100 | 121 | 144 |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| GMRES-PMR | 6 | 9 | 11 | 14 | 16 | 18 | 19 | 20 | 21 | 23 | 24 |
| GMRES-AltMR | 5 | 8 | 11 | 14 | 16 | 19 | 22 | 24 | 27 | 30 | 32 |
| GMRES-PMR c.c. | 3 | 3 | 3 | 3 | 3 | 3 | 3 | 3 | 3 | 3 | 3 |
| GMRES-AltMR c.c. | 3 | 3 | 3 | 3 | 3 | 4 | 4 | 4 | 4 | 4 | 4 |
| AltMR c.c. | 8 | 8 | 9 | 9 | 9 | 9 | 9 | 9 | 9 | 9 | 9 |
| PMR c.c. | 6 | 7 | 8 | 8 | 9 | 9 | 9 | 9 | 9 | 9 | 9 |

TABLE 1. Number of GMRES iterations necessary to get a residual smaller than $10^{-10}$. The different columns correspond to different numbers of objects and show the number of iterations performed by GMRES and various preconditioners: AltMR, PMR with coarse correction, and AltMR with coarse correction. In the two last lines, we show the number of iterations required by coarse corrected AltMR and PMR used as stationary iterations.

| \# Objects $(r=0.3)$ | 4 | 9 | 16 | 25 | 36 | 49 | 64 | 81 | 100 | 121 | 144 |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| GMRES-PMR | 10 | 15 | 19 | 22 | 25 | 28 | 30 | 33 | 34 | 37 | 39 |
| GMRES-AltMR | 7 | 11 | 15 | 19 | 23 | 27 | 32 | 36 | 40 | 44 | 49 |
| GMRES-PMR c.c. | 5 | 6 | 6 | 6 | 6 | 6 | 7 | 6 | 7 | 7 | 7 |
| GMRES-AltMR c.c. | 5 | 6 | 6 | 7 | 7 | 8 | 8 | 8 | 8 | 8 | 8 |
| AltMR c.c. | 18 | 17 | 21 | 19 | 22 | 20 | 22 | 22 | 22 | 23 | 23 |
| PMR c.c. | 12 | 15 | 18 | 20 | 22 | 23 | 23 | 23 | 24 | 25 | 26 |

Table 2. Same experiment setting as in Table 1, but for a larger object radius $r=0.3$.
also not converging as a stationary iterative method [13, 15]. We also see that with our coarse correction one obtains scalability with respect to the number of objects. This scalability also holds for the stationary iterative variants. We next repeat the experiment using the same geometry but increasing the radius of the objects to $r=0.3$, which implies that the distance between the objects is smaller, and thus a slower convergence rate of the methods of reflections. The number of iterations are summarized in Table 2. We see that indeed iteration numbers are now larger, but the methods of reflections with our coarse correction are still scalable.

| Step size | 0.1000 | 0.0500 | 0.0250 | 0.0125 | 0.0063 | 0.0031 | 0.0016 | 0.0008 | 0.0004 |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| GMRES-PMR | 9 | 9 | 9 | 9 | 8 | 9 | 9 | 8 | 8 |
| GMRES-AltMR | 8 | 8 | 8 | 8 | 8 | 8 | 8 | 8 | 8 |
| GMRES-PMR c.c. | 3 | 3 | 3 | 3 | 3 | 3 | 3 | 3 | 3 |
| GMRES-AltMR c.c. | 3 | 3 | 3 | 3 | 3 | 3 | 3 | 3 | 3 |
| AltMR c.c. | 8 | 8 | 8 | 8 | 8 | 8 | 8 | 8 | 8 |
| PMR c.c. | 7 | 7 | 7 | 7 | 7 | 7 | 7 | 7 | 8 |

Table 3. Same experiment setting as in Table 1, but for a fixed number of objects when the mesh is refined.

We finally repeat this test with a fixed number of $J=L^{3}=9$ objects of radius $r=0.1$, but using various step sizes for the boundary discretization. The results are shown in Table 3. We observe scalability with respect to the mesh size in all cases.

## 8. Conclusions

We presented an extensive analysis of the alternating and parallel methods of reflections using domain decomposition techniques. We proved that for two objects, the methods of reflections can be identified with Schwarz domain decomposition methods, which led to new, sharp convergence estimates in this case. We also showed that for more than two objects, the methods of reflections are different from Schwarz domain decomposition methods. We then used substructuring techniques from domain decomposition to rewrite the methods of reflections only iterating on traces, which allowed us to identify these methods with block Gauss-Seidel and block Jacobi methods. Using this insight, we derived new, relaxed variants of the methods of reflections, which converge faster that the classical variants, and also introduced for the first time a coarse correction for the methods of reflections to make them scalable when the number of objects becomes large. The substructured formulation allowed us also naturally to use these methods as preconditioners for Krylov methods. All our results were obtained for the case of the Laplacian, but other operators could be used as well in the relations we found.

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