# DISCRETE ANALYSIS OF DOMAIN DECOMPOSITION APPROACHES FOR MESH GENERATION VIA THE EQUIDISTRIBUTION PRINCIPLE 

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#### Abstract

Moving mesh methods based on the equidistribution principle are powerful techniques for the space-time adaptive solution of evolution problems. Solving the resulting coupled system of equations, namely the original PDE and the mesh PDE, however, is challenging in parallel. Recently several Schwarz domain decomposition algorithms were proposed for this task and analyzed at the continuous level. However, after discretization, the resulting problems may not even be well posed, so the discrete algorithms requires a different analysis, which is the subject of this paper. We prove that when the number of grid points is large enough, the classical parallel and alternating Schwarz methods converge to the unique monodomain solution. Thus, such methods can be used in place of Newton's method, which can suffer from convergence difficulties for challenging problems. The analysis for the nonlinear domain decomposition algorithms is based on $M$-function theory and is valid for an arbitrary number of subdomains. An asymptotic convergence rate is provided and numerical experiments illustrate the results.


Key words. domain decomposition, Schwarz methods, moving meshes, equidistribution, discretization, $M$-functions

AMS subject classifications. $65 \mathrm{M} 55,65 \mathrm{~N} 22,65 \mathrm{Y} 05,65 \mathrm{M} 50,65 \mathrm{~N} 50$

1. Introduction. The numerical simulation of partial differential equations (PDEs) can often benefit from the use of non-uniform meshes chosen to adapt to the local solution structures, especially in cases where the solutions vary over disparate space and/or time scales. In this article, we consider a class of $r$-refinement or moving mesh methods. This PDE-based mesh generation approach adapts an initial (typically uniform) grid by relocating a fixed number of mesh nodes while preserving the mesh topology. The mesh is determined by solving a so-called moving mesh PDE (MMPDE) which is coupled to the physical PDE of interest.

In its simplest form, the $r$-refinement technique is guided in large part by the equidistribution principle (EP) of de Boor [15] and Burchard [7]. In one dimension, the EP seeks a mesh so that the error contribution on each interval is the same. This naturally places mesh points closer together in regions of the physical space $\Omega_{p}$ where the error in the solution $u(t, x)$ is large. Mathematically, given some positive measure of the error, $M(t, x, u)$, for $x \in \Omega_{p}$ and $t \in[0, T]$, the EP requires that the mesh points $x_{i}, i=0, \ldots, N$, satisfy

$$
\int_{x_{i-1}}^{x_{i}} M(t, \tilde{x}, u) d \tilde{x} \equiv \frac{1}{N} \int_{0}^{1} M(t, \tilde{x}, u) d \tilde{x}
$$

The function $M$ is referred to as the monitor or mesh density function.
In the continuous setting, we seek a time dependent mesh transformation between an underlying computational coordinate $\xi \in \Omega_{c}$ and the physical coordinate, $x \in \Omega_{p}$, so that

$$
\begin{equation*}
\int_{0}^{x\left(\xi_{i}, t\right)} M(t, \tilde{x}, u) d \tilde{x}=\frac{i}{N} \theta(t) \equiv \xi_{i} \theta(t) \tag{EP}
\end{equation*}
$$

[^0]where $\theta(t) \equiv \int_{0}^{1} M(t, \tilde{x}, u) d \tilde{x}$ is the total error in the solution. It follows directly, by differentiating the continuous form of (EP), that the required mesh transformation satisfies, for all $t$, the nonlinear differential equation
\[

$$
\begin{equation*}
\frac{\partial}{\partial \xi}\left\{M(t, x(\xi, t), u) \frac{\partial}{\partial \xi} x(\xi, t)\right\}=0 \tag{1.1}
\end{equation*}
$$

\]

For time dependent problems the EP is often relaxed [30], and the mesh may be found by solving

$$
\begin{equation*}
\frac{\partial x}{\partial t}=\frac{1}{\tau} \frac{\partial}{\partial \xi}\left(M(t, x(\xi, t), u) \frac{\partial x}{\partial \xi}\right) . \tag{MMPDE5}
\end{equation*}
$$

The relaxation parameter $\tau$ is chosen in practice (see [29]) so that the mesh evolves at a rate commensurate with that of the physical solution $u(t, x)$. A parabolic mesh equation provides a degree of temporal smoothing for the mesh generation problem. Coupling (MMPDE5) with a time dependent physical PDE and discretizing in space gives a initial value problem to solve. In contrast, coupling (1.1) with a physical PDE and discretizing in space gives an often more difficult DAE problem. The wellposedness of the steady mesh generation problem (1.1) and the semi-discretized (in time) form of (MMPDE5) may be found in [22]. A thorough treatise of $r$-refinement methods may be found in the recent book [32].

The mesh generators (1.1) and (MMPDE5) depend on the physical solution $u$ which itself is the unknown solution of a physical PDE. In practice, this coupled system of DEs may be solved using the $M P$ iterative procedure [32]. An initial guess for the mesh is chosen and the physical PDE is solved to give an approximate solution. This mesh and physical solution can be used to evaluate $M$ and then (1.1) or (MMPDE5) can be solved for the new mesh, and the process may be repeated. This is in contrast to the approach, taken in MOVCOL [31] for example, which solves for the mesh and physical solution as one large simultaneous system.

Recently we have been interested in the parallel solution of PDEs on (moving) equidistributing grids. Our solution strategy introduces spatial parallelism in both the computation of the underlying meshes and the physical PDE by a domain decomposition (DD) approach. DD solves a PDE by divide and conquer: the spatial domain $\Omega$ is partitioned into overlapping or non-overlapping subdomains, thus reformulating the original PDE into an equivalent coupled system of PDEs which is solved by iteration. The coupling is provided by transmission conditions which are designed to appropriately match solutions on neighboring subdomains. The matching is imposed in a way which accelerates the convergence of the DD iteration and to maintain as much coarse grain granularity in the computation as possible. This allows implementation on a distributed memory architecture using MPI (for example). A stochastic DD approach for mesh generation in 2D which does not require iteration has recently been considered in [4].

Another benefit of the DD approach is its robustness relative to other nonlinear solution schemes. When the monitor function $M$ is "difficult" (e.g. highly oscillatory), Newton's method often fails to converge when applied to the discretized versions of (EP) and (MMPDE5). Then it is natural to divide the interval [0,1] into several subdomains to obtain smaller, "easier" problems, and use a DD iteration to obtain the global solution. We will illustrate this approach in Section 6.3.

In [22], various parallel Schwarz algorithms were presented to solve (1.1) and the semi-discretized (in time) form of (MMPDE5) using domain decomposition in the
computational co-ordinate $\xi$. These algorithms were analyzed at the continuous level, were shown to converge, and contraction rates obtained. Alternating algorithms, which give faster convergence at the cost of losing immediate parallelization, were presented in [28]. These analyses assume $M$ is a function of the required mesh only, as is the case in the $M P$ procedure mentioned above.

In practice these Schwarz methods have to be discretized for implementation. For instance, (1.1) can be discretized as

$$
\begin{gather*}
\tilde{M}\left(x_{i}, x_{i+1}\right)\left(x_{i+1}-x_{i}\right)-\tilde{M}\left(x_{i-1}, x_{i}\right)\left(x_{i}-x_{i-1}\right)=0, \quad i=1, \ldots, N-1  \tag{1.2}\\
x_{0}=x_{L}, \quad x_{N}=x_{R}
\end{gather*}
$$

where $x_{L}=x(0)$ and $x_{R}=x(1)$ are the boundary conditions, and the quantity $\tilde{\sim} \tilde{\sim}(u, v)$ is chosen so that (1.2) is a consistent discretization of (1.1). For example, $\tilde{M}(u, v)=M((u+v) / 2)$ and $\tilde{M}(u, v)=(M(u)+M(v)) / 2$ gives the midpoint and trapezoidal discretizations respectively.

Likewise, we discretize (MMPDE5) using finite differences in space and the $\theta$ method in time, $\frac{1}{2} \leq \theta \leq 1$, to obtain the update from $t^{[\nu]}$ to $t^{[\nu+1]}$ as the solution of

$$
\begin{gather*}
\frac{x_{i}^{[\nu+1]}-x_{i}^{[\nu]}}{\Delta t}-\frac{\theta}{\tau \Delta \xi^{2}}\left(\tilde{M}\left(x_{i}^{[\nu+1]}, x_{i+1}^{[\nu+1]}\right)\left(x_{i+1}^{[\nu+1]}-x_{i}^{[\nu+1]}\right)-\tilde{M}\left(x_{i-1}^{[\nu+1]}, x_{i}^{[\nu+1]}\right)\left(x_{i}^{[\nu+1]}-x_{i-1}^{[\nu+1]}\right)\right) \\
-\frac{(1-\theta)}{\tau \Delta \xi^{2}}\left(\tilde{M}\left(x_{i}^{[\nu]}, x_{i+1}^{[\nu]}\right)\left(x_{i+1}^{[\nu]}-x_{i}^{[\nu]}\right)-\tilde{M}\left(x_{i-1}^{[\nu]}, x_{i}^{[\nu]}\right)\left(x_{i}^{[\nu]}-x_{i-1}^{[\nu]}\right)\right)=0  \tag{1.3}\\
x_{0}^{[\nu+1]}=x_{L} \quad x_{N}^{[\nu+1]}=x_{R}
\end{gather*}
$$

for $i=1, \ldots, N-1$.
The goal of this paper is to prove that whenever the number of grid points $N$ is large enough, the parallel and alternating Schwarz algorithms applied to the discrete problems (1.2) and (1.3) do indeed converge to the discrete monodomain solution. We will also estimate the rate of convergence in terms of the number of grid points and parameters related to the monitor function $M$. It is important to note that the analysis of the discrete algorithms is not a simple extension of the continuous theory presented in [22, 28]. That analysis depended heavily on implicit integral representations of the subdomain solutions and on properties of those integral representations. These properties do not carry over once quadrature has been applied. In fact, even in the single-domain case, the discretized problem may not be well-posed when the number of grid points is not large enough, see Example 2.1, and also Xu et al. [48], where the authors considered quadrature using piecewise constant or linear polynomials in the context of de Boor's algorithm. Instead, the theory here relies heavily on the notions of nonlinear $M$-functions and the method of super- and subsolutions in $\mathbb{R}^{n}$ popularized by Ortega and Rheinboldt, see for example [44, 41].

The literature involving DD approaches for linear elliptic PDEs is vast, see the extensive books [47, 39]. We focus in this paper on Schwarz methods, since these methods have been studied for time dependent problems in two fundamentally different approaches: the more classical approach, in which one applies the alternating or parallel Schwarz method to the sequence of elliptic problems obtained from semidiscretization (in time) of the PDE [8, 9], or the more recent approach, where one decomposes the space-time domain and applying the Schwarz waveform relaxation algorithm [5, 24, 25]. For both approaches, optimized variants have been developed, see
[18] and the references therein for the first, and [21, 19, 20, 2] for the second approach. The linear theory, however, does not apply to the mesh PDEs considered here, since they are highly nonlinear. While DD methods have been applied successfully to a large variety of nonlinear problems (see e.g. [27,3] for porous media flow problems, [45, 1] for problems from physiology, [26] for nonlinear parabolic systems, and [14] for the compressible Navier-Stokes equations), comparably much less has been written about the analysis of DD methods for nonlinear PDEs; see [10, 38, 37, 36, 16, 6, 46] for the steady case, and [17, 23] for evolution problems.

There are two common alternatives to the approach described in this paper. In the first, known as the Newton-Krylov-Schwarz (NKS) method, one discretizes the global nonlinear problem and solves the nonlinear algebraic equations with an inexact Newton method. The linear system that produces a Newton step is solved using a preconditioned Krylov method, where the inner (linear) iterations can be parallelized using a DD preconditioner. In the second approach, known as the nonlinear Additive Schwarz Preconditioned Inexact Newton (ASPIN) method, the discretized global nonlinear problem is not solved directly, but is replaced by a 'preconditioned' nonlinear function that has the same set of solutions as the original. Typically, the preconditioned function is a sum of local solutions of subdomain problems, so that it can be evaluated in parallel; thus, it is the analogue of the additive Schwarz method for linear problems. The preconditioned nonlinear problem is then solved using an inexact Newton method without additional linear preconditioning. See [11, 12, 42, 34] for more details. More recently, a two-level ASPIN algorithm has been proposed and studied in [35]. The multi-level approach and coarse space allows one to deal with deterioration of the convergence rate of the DD algorithm as the number of subdomains increases.

The paper is organized as follows. In Section 2, we begin by illustrating an example of (1.2) which is not well-posed. Then we introduce notation and review the main theoretical tools used in our analysis, namely the theory of $M$-functions. We then proceed in Sections 2.2 and 2.3 to specify sufficient conditions which guarantee the discrete systems for steady and time dependent mesh generation are $M$-functions. Our convergence results rely on the existence of super and subsolutions for the discrete systems; this is demonstrated in Section 3. Section 4 provides the main convergence results for both the parallel and alternating Schwarz methods for the discrete mesh generators. The asymptotic convergence rate of the discrete parallel Schwarz is given in Section 5, where we also show that the convergence rate coincides with the continuous analysis in the fine-grid limit. In Section 6, we illustrate our theoretical results with numerical examples, show the increased robustness of the DD approach relative to Newton's method, and illustrate the use of DD within the $M P$ iteration approach. Finally, we draw some conclusions in Section 7.
2. Uniqueness of solutions to the discrete equidistribution problem. As we mentioned in the introduction, once the equidistribution equation (1.1) is discretized by a quadrature rule, it is no longer clear that a unique solution always exists. We start this section by an example showing non-uniqueness of the discrete solution when there are not enough points on the grid.

Example 2.1. Consider the positive function $M:[0,1] \rightarrow \mathbb{R}^{+}$shown in Figure 2.1, defined as

$$
M(x)= \begin{cases}1+30 x-\frac{171}{4} x^{2}, & 0 \leq x \leq \frac{2}{3}, \\ \frac{4-5 x}{2 x-1}, & \frac{2}{3} \leq x \leq \frac{3}{4}, \\ \frac{1}{7}\left(836-2496 x+2328 x^{2}-640 x^{3}\right), & \frac{3}{4} \leq x \leq 1\end{cases}
$$

This $M$ is chosen so that $M(0)=1, M(1)=4$, and the function is continuously


Fig. 2.1. An example mesh density function leading to a non-unique mesh transformation.
differentiable at $2 / 3$ and $3 / 4$. Suppose we wish to solve the steady mesh generation problem with exactly one interior node $\alpha$, i.e., we want to find $\alpha \in(0,1)$ such that

$$
\int_{0}^{\alpha} M(u) d u=\int_{\alpha}^{1} M(u) d u .
$$

The solution to the continuous problem is unique, with $\alpha \approx 0.38$. However, if we replace the integral by the trapezoidal rule, i.e., if we solve

$$
\begin{equation*}
\alpha \frac{M(0)+M(\alpha)}{2}=(1-\alpha) \frac{M(\alpha)+M(1)}{2}, \tag{2.1}
\end{equation*}
$$

Then we observe that for $\frac{2}{3} \leq \alpha \leq \frac{3}{4}$, the left hand side of (2.1) becomes

$$
\alpha \frac{M(0)+M(\alpha)}{2}=\frac{\alpha(2 \alpha-1+4-5 \alpha)}{2(2 \alpha-1)}=\frac{3}{2} \cdot \frac{1-\alpha}{2 \alpha-1},
$$

and the right hand side gives

$$
(1-\alpha) \frac{M(\alpha)+M(1)}{2}=\frac{(1-\alpha) 4-5 \alpha+4(2 \alpha-1)}{2(2 \alpha-1)}=\frac{3}{2} \cdot \frac{1-\alpha}{2 \alpha-1} .
$$

Thus, every $\alpha$ in the interval $\left[\frac{2}{3}, \frac{3}{4}\right]$ is a solution of (2.1), so the discrete problem, defined on only 3 nodes, is not well-posed.

The problem in the above example is that there are too few points in the grid to "nail down" the area under the curve with sufficient precision to obtain a unique
solution. One of the main theoretical results of this paper states that whenever the grid is fine enough (where 'fine enough' depends on the mesh density function $M$ ), the discrete problem also has a unique solution. The main idea is to show that the equidistribution equations satisfy the $M$-function property in the sense of Ortega and Rheinboldt [44, 41] whenever the grid is fine enough. Uniqueness then follows from the $M$-function structure; in addition, the convergence of the parallel and alternating Schwarz methods is also a consequence of this $M$-function structure.
2.1. $M$-functions: definitions and properties. Consider a finite dimensional system of nonlinear equations $F(\mathbf{x})=0$ where $F: \mathbb{D} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ and $\mathbf{x}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)^{T}$ is a vector of unknowns. We will refer to the $i$ th equation in the system as $f_{i}(\mathbf{x})$. The natural partial ordering (component-wise) on $R^{n}$ is defined by

$$
\mathbf{x} \leq \mathbf{y}, \mathbf{x}, \mathbf{y} \in R^{n} \quad \text { if and only if } \quad x_{i} \leq y_{i}, \text { for } i=1, \ldots, n
$$

We will let $\mathbf{e}^{i}, i=1, \ldots, n$ denote the $i$-th standard basis vector in $\mathbb{R}^{n}$.
We begin by defining isotone and antitone mappings and then diagonally isotone and off-diagonally antitone mappings, see [41].

Definition 2.2. A mapping $F: \mathbb{D} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is isotone (antitone) if for all $x, y \in \mathbb{D}, F(\mathbf{x}) \leq F(\mathbf{y})(F(\mathbf{x}) \geq F(\mathbf{y}))$ whenever $\mathbf{x} \leq \mathbf{y}$. It is strictly isotone (antitone) if $F(\mathbf{x})<F(\mathbf{y})(F(\mathbf{x})>F(\mathbf{y}))$ whenever $\mathbf{x}<\mathbf{y}$.

Definition 2.3. A mapping $F: \mathbb{D} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is diagonally isotone if, for any $\mathbf{x} \in \mathbb{D}$, the $n$ functions $f_{i}\left(\mathbf{x}+t \mathbf{e}^{i}\right), i=1, \ldots, n$, are isotone when $\mathbf{x}+t \mathbf{e}^{i} \in \mathbb{D}$. $F$ is offdiagonally antitone $i f$, for any $\mathbf{x} \in \mathbb{R}^{n}$, the functions $f_{i}\left(\mathbf{x}+t \mathbf{e}^{j}\right), i \neq j, i, j=1, \ldots, n$, are antitone when $\mathbf{x}+t \mathbf{e}^{i} \in \mathbb{D}$.

Collatz [13] introduced the following converse notion of isotonicity on partially ordered topological spaces.

Definition 2.4. A mapping $F: \mathbb{D} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is inverse isotone if for any $\mathbf{x}, \mathbf{y} \in \mathbb{D}, F(\mathbf{x}) \leq F(\mathbf{y})$ implies $\mathbf{x} \leq \mathbf{y}$.

The following notion of an $M$-function was originally introduced by Ortega and developed by Rheinboldt [44].

Definition 2.5. A mapping $F: \mathbb{D} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is an $M$-function if $F$ is inverse isotone and off-diagonally antitone.

THEOREM 2.6 (Uniqueness). Let $F: \mathbb{D} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be an $M$-function and consider the problem $F(\mathbf{x})=\mathbf{b}$, where $\mathbf{b} \in \mathbb{R}^{n}$ is a given right hand side. Then the solution $\mathbf{x} \in \mathbb{D}$, if it exists, must be unique.

Proof. Let $\mathbf{x}$ and $\mathbf{y}$ be two solutions satisfying $F(\mathbf{x})=F(\mathbf{y})=\mathbf{b}$. Then since $F(\mathbf{x}) \leq F(\mathbf{y})$, the inverse isotonicity of $F$ implies that $\mathbf{x} \leq \mathbf{y}$. On the other hand, we also have $F(\mathbf{y}) \leq F(\mathbf{x})$, which implies $\mathbf{y} \leq \mathbf{x}$. It follows that $\mathbf{x}=\mathbf{y}$, i.e., the solution is unique.

Of course, the difficulty in determining if a map is a $M$-function is checking the inverse isotonicity. A sufficient condition requires the following definition, which first appeared in [44] and specifies the connectivity of the unknowns through $F$.

Definition 2.7. Let $F: \mathbb{D} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be off-diagonally antitone. Then a link $(i, j)$ is strict if the function $t \mapsto f_{i}\left(\mathbf{x}+t \mathbf{e}^{j}\right)$ is strictly antitone. A path $i \rightsquigarrow j$ exists if there exists a sequence of strict links $\left(i, j_{1}\right),\left(j_{1}, j_{2}\right), \ldots,\left(j_{k}, j\right)$.

A sufficient condition for $F: \mathbb{D} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ to be an $M$-function is given in the following lemma, a modification of Theorem 5.1 in Rheinboldt [44]. See the Appendix for a proof of this result.

Lemma 2.8. Let $\mathbb{D}$ be a convex and open subset of $\mathbb{R}^{n}$. Assume $F: \mathbb{D} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is off-diagonally antitone, and that for any $\mathbf{x} \in \mathbb{D}$, the functions $q_{i}: S_{i} \subset \mathbb{R} \rightarrow \mathbb{R}^{n}$ defined as

$$
q_{i}(\tau)=\sum_{j=1}^{n} f_{j}\left(\mathbf{x}+\tau \mathbf{e}^{i}\right), i=1, \ldots, n, \quad \text { with } \quad S_{i}=\left\{\tau: \mathbf{x}+\tau \mathbf{e}^{i} \in \mathbb{D}\right\}
$$

are isotone. Finally, assume for every $i$ that either $q_{i}$ is strictly isotone, or there exists a node $j=j(i)$ such that a path $j \rightsquigarrow i$ exists and $q_{j}$ is strictly isotone. Then $F$ is an $M$-function.

Here we are interested in mappings $F$ defined by the simple finite difference discretizations (1.2) and (1.3). We set $h=1 / N$ and choose the uniform grid $\xi_{i}=$ $i h, i=0, \ldots, N$, in the computational domain. Since the values $x_{0}, x_{1}, \ldots, x_{N}$ in (1.2) and (1.3) are meant to be samples of a continuous function $x(\xi)$ at $\xi_{i}$, we expect successive samples $x_{i-1}$ and $x_{i}$ to be close when the grid is fine enough. Thus, we will consider the nonlinear mappings defined by (1.2) and (1.3) on the subset $D$ of $\mathbb{R}^{N+1}$ given by

$$
\begin{equation*}
D=\left\{\mathbf{x} \in R^{N+1}:\left|x_{i}-x_{i-1}\right|<\frac{\check{m}}{B}, \text { for } i=1, \ldots, N\right\} \tag{2.2}
\end{equation*}
$$

where $\check{m}>0$ and $B>0$ are constants that depend on the function $M(x)$, to be defined later. From (2.2) it is clear that $D$ is open and convex, as required by Lemma 2.8.
2.2. Steady mesh generation. Throughout this paper, we assume the twopoint mesh density function $\tilde{M}(\cdot, \cdot)$ satisfies the following conditions:
(C1) there exists constants $\check{m}$ and $\hat{m}$ such that

$$
0<\check{m} \leq \tilde{M}(u, v) \leq \hat{m}<\infty \quad \text { for all } u, v \in \mathbb{R}
$$

and
(C2) $\tilde{M}(\cdot, \cdot)$ is continuously differentiable with respect to each variable, and there exists a positive constant $B<\infty$ such that

$$
\left|\frac{\partial \tilde{M}}{\partial u}(u, v)\right| \leq B, \quad\left|\frac{\partial \tilde{M}}{\partial v}(u, v)\right| \leq B \quad \text { for all } u, v \in \mathbb{R}
$$

(C3) for every $u, v \in \mathbb{R}$ such that $u \leq v$, there exists $\zeta \in[u, v]$ such that

$$
\tilde{M}(u, v)=M(\zeta)
$$

The condition (C3) simply assures consistency between $M$ and $\tilde{M}$. The constants $\check{m}$, $\hat{m}$ and $B$ can typically derived from the bounds and Lipschitz constant of the continuous mesh density function $M(x)$. For instance, for the trapezoidal rule $\tilde{M}(u, v)=$ $\frac{1}{2}(M(u)+M(v))$, we have $B=\frac{1}{2} \max _{u \in R}\left|M^{\prime}(u)\right|$.

Remark. Typically $M(x)$ is only defined on the interval $[0,1]$, since the mesh is only defined on that interval. However, for the purpose of proving existence, we will need to evaluate $M$ at points outside the unit interval; this can be done by extending $M(x)$ smoothly outside [0, 1] in such a way that $M$ becomes constant outside a slightly larger interval $\left[-x_{\min }, x_{\max }\right] \supset[0,1]$. Thus, it is reasonable to assume that $\check{m}$ and $\hat{m}$ exist, and one can expect these bounds to be close to the maximum and minimum of
$M$ inside the unit interval. A similar argument goes for the constant $B$, as long as the extension does not introduce gradients much steeper than those inside the unit interval.

We now use Lemma 2.8 to show that $H$ defined by (1.2) is an $M$-function on $D$. Theorem 2.9. The mapping $H: D \subset R^{N+1} \rightarrow R^{N+1}$ defined by

$$
h_{i}= \begin{cases}(\hat{m}+\check{m})\left(x_{0}-x_{L}\right), & i=0  \tag{2.3}\\ \tilde{M}\left(x_{i-1}, x_{i}\right)\left(x_{i}-x_{i-1}\right)-\tilde{M}\left(x_{i}, x_{i+1}\right)\left(x_{i+1}-x_{i}\right), & i=1, \ldots, N-1 \\ (\hat{m}+\check{m})\left(x_{N}-x_{R}\right), & i=N\end{cases}
$$

is an $M$-function on the set $D$ defined in (2.2) if $M(x)$ satisfies the conditions (C1) and (C2).

Proof. For $1 \leq i \leq N-1$, the $i$ th component of $H$ reads

$$
h_{i}=\phi_{i-1 / 2}\left(x_{i-1}, x_{i}\right)-\phi_{i+1 / 2}\left(x_{i}, x_{i+1}\right)
$$

where we have introduced the notation $\phi_{i-1 / 2}:=\tilde{M}\left(x_{i-1}, x_{i}\right)\left(x_{i}-x_{i-1}\right)$ representing the discrete flux between $x_{i}$ and $x_{i-1}$.

The function $h_{i}$ depends on $x_{j}$ only for $j=i-1, i$ and $i+1$. Computing directly for $i=1, \ldots, N-1$ we have

$$
\begin{aligned}
\frac{\partial h_{i}}{\partial x_{i-1}} & =\frac{\partial \phi_{i-1 / 2}}{\partial x_{i-1}}=-\tilde{M}\left(x_{i-1}, x_{i}\right)+\tilde{M}_{u}\left(x_{i-1}, x_{i}\right)\left(x_{i}-x_{i-1}\right) \\
\frac{\partial h_{i}}{\partial x_{i+1}} & =-\frac{\partial \phi_{i+1 / 2}}{\partial x_{i+1}}=-\tilde{M}\left(x_{i}, x_{i+1}\right)-\tilde{M}_{v}\left(x_{i}, x_{i+1}\right)\left(x_{i+1}-x_{i}\right)
\end{aligned}
$$

Here $\tilde{M}_{u}$ and $\tilde{M}_{v}$ denote the derivatives of $\tilde{M}(u, v)$ with respect to its first and second arguments respectively. Following Lemma 2.8, we begin by showing that $F$ is offdiagonally antitone on the set $D$ defined in (2.2).

Since the mesh $\left\{x_{i}\right\}_{i=0}^{N}$ belongs to $D$ and $\tilde{M}$ is Lipschitz with constant $B$ (from condition (C2)), we have

$$
\tilde{M}_{u}\left(x_{i-1}, x_{i}\right)\left(x_{i}-x_{i-1}\right) \leq\left|\tilde{M}_{u}\left(x_{i-1}, x_{i}\right)\right| \cdot\left|x_{i}-x_{i-1}\right|<B \frac{\check{m}}{B}=\check{m} \leq \tilde{M}\left(x_{i-1}, x_{i}\right)
$$

and

$$
-\tilde{M}_{v}\left(x_{i}, x_{i+1}\right)\left(x_{i+1}-x_{i}\right) \leq\left|\tilde{M}_{v}\left(x_{i}, x_{i+1}\right)\right| \cdot\left|x_{i+1}-x_{i}\right|<B \frac{\check{m}}{B}=\check{m} \leq \tilde{M}\left(x_{i}, x_{i+1}\right)
$$

Hence $\frac{\partial h_{i}}{\partial x_{i-1}}$ and $\frac{\partial h_{i}}{\partial x_{i+1}}$ are strictly negative for $i=1, \ldots, N-1$. Furthermore, $\partial h_{0} / \partial x_{i}=0$ for $i=1, \ldots, N$ and $\partial h_{N} / \partial x_{i}=0$ for $i=0, \ldots, N-1$ which implies $H$ is off-diagonally antitone.

We now construct the functions $q_{i}(t)=\sum_{j=0}^{N} h_{j}\left(\mathbf{x}+t \mathbf{e}^{i}\right), i=1, \ldots, N-1$, and take derivatives with respect to $t$. Direct calculation gives $d q_{i} / d t=0$ for $i=$ $2, \ldots, N-2$, while

$$
\begin{aligned}
\frac{d q_{0}}{d t} & =\hat{m}+\check{m}-\tilde{M}\left(x_{0}+t, x_{1}\right)+\tilde{M}_{u}\left(x_{0}+t, x_{1}\right)\left(x_{1}-x_{0}-t\right) \\
\frac{d q_{1}}{d t} & =\frac{d}{d t} \phi_{1 / 2}\left(x_{0}, x_{1}+t\right)=\tilde{M}_{v}\left(x_{0}, x_{1}+t\right)\left(x_{1}+t-x_{0}\right)+\tilde{M}\left(x_{0}, x_{1}+t\right)
\end{aligned}
$$

and

$$
\begin{aligned}
\frac{d q_{N-1}}{d t} & =-\tilde{M}_{u}\left(x_{N-1}+t, x_{N}\right)\left(x_{N}-x_{N-1}-t\right)+\tilde{M}\left(x_{N-1}+t, x_{N}\right) \\
\frac{d q_{N}}{d t} & =\hat{m}+\check{m}-\tilde{M}\left(x_{N-1}, x_{N}+t\right)-\tilde{M}_{v}\left(x_{N-1}, x_{N}+t\right)\left(x_{N}+t-x_{N-1}\right)
\end{aligned}
$$

Using condition (C2) and the assumption that the perturbed mesh $x+t e^{i}$ belongs to $D$ it is easy to see that $\frac{d q_{0}}{d t}, \frac{d q_{1}}{d t}, \frac{d q_{N-1}}{d t}, \frac{d q_{N}}{d t}$ are strictly positive. Hence the functions $q_{i}$ are isotone. Finally, the nodes $x_{1}$ and $x_{N-1}$ are connected to all other nodes, with $q_{1}$ and $q_{N-1}$ strictly isotone. So by Lemma 2.8, the function $H$ defined by (2.3) is an $M$-function on $D$.
2.3. Time dependent mesh generation. Given a physical solution at time level $t^{[\nu]}$, the updated mesh at time $t^{[\nu+1]}$ can be found by solving the nonlinear system (1.3). We now verify that this system is described by a $M$-function.

THEOREM 2.10. The mapping $G: D \subset \mathbb{R}^{N+1} \rightarrow \mathbb{R}^{N+1}$, for the unknowns $\left\{x_{i}^{[\nu+1]}\right\}_{i=0}^{N}$, defining (1.3) is an $M$-function on the set $D$ defined in (2.2) if $\tilde{M}(u, v)$ satisfies the conditions (C1) and (C2).

Proof. For $\nu$ fixed, the unknowns $x^{[\nu+1]}$ satisfy $G\left(x^{[\nu+1]}\right)=0$ where from (1.3) the $i$ th component of $G$ reads
$g_{i}= \begin{cases}\theta \gamma(\hat{m}+\check{m})\left(x_{0}^{[\nu+1]}-x_{L}\right), & i=0 \\ x_{i}^{[\nu+1]}-\theta \gamma\left(\phi_{i+1 / 2}^{[\nu+1]}-\phi_{i-1 / 2}^{[\nu+1]}\right)-x_{i}^{[\nu]}-(1-\theta) \gamma\left(\phi_{i+1 / 2}^{[\nu]}-\phi_{i-1 / 2}^{[\nu]}\right), & i=1, \ldots, N-1, \\ \theta \gamma(\hat{m}+\check{m})\left(x_{N}^{[\nu+1]}-x_{R}\right), & i=N\end{cases}$
where we have introduced the notation $\phi_{i-1 / 2}^{[\nu]}:=\tilde{M}\left(x_{i-1}^{[\nu]}, x_{i}^{[\nu]}\right)\left(x_{i}^{[\nu]}-x_{i-1}^{[\nu]}\right)$ and $\gamma:=$ $\Delta t /\left(\tau \Delta \xi^{2}\right)$.

Computing directly for $i=1, \ldots, N-1$ we have
$\frac{\partial g_{i}}{\partial x_{i-1}^{[\nu+1]}}=\theta \gamma \frac{\partial \phi_{i-1 / 2}}{\partial x_{i-1}^{[\nu+1]}}=\theta \gamma\left(-\tilde{M}\left(x_{i-1}^{[\nu+1]}, x_{i}^{[\nu+1]}\right)+\tilde{M}_{u}\left(x_{i-1}^{[\nu+1]}, x_{i}^{[\nu+1]}\right)\left(x_{i}^{[\nu+1]}-x_{i-1}^{[\nu+1]}\right)\right)$,
$\frac{\partial g_{i}}{\partial x_{i+1}^{[\nu+1]}}=-\theta \gamma \frac{\partial \phi_{i+1 / 2}}{\partial x_{i+1}^{[\nu+1]}}=-\theta \gamma\left(\tilde{M}\left(x_{i}^{[\nu+1]}, x_{i+1}^{[\nu+1]}\right)+\tilde{M}_{v}\left(x_{i}^{[\nu+1]}, x_{i+1}^{[\nu+1]}\right)\left(x_{i+1}^{[\nu+1]}-x_{i}^{[\nu+1]}\right)\right.$.
Once again $\tilde{M}_{u}$ and $\tilde{M}_{v}$ denote the derivatives of $\tilde{M}(u, v)$ with respect to its first and second arguments respectively. Moreoever $\partial g_{0} / \partial x_{i}=0$ for $i=1, \ldots, N$ and $\partial g_{N} / \partial x_{i}=0$ for $i=0, \ldots, N-1$. Arguing exactly as in the proof of Theorem 2.9 we can see that $G$ is off-diagonally antitone on the set $D$ defined in (2.2).

We now construct the functions $q_{i}(t)$ specified in Lemma 2.8. and take derivatives with respect to $t$. Direct calculation gives $d q_{i} / d t=1$ (strictly positive) for $i=$ $2, \ldots, N-2$. Also, we have

$$
\begin{aligned}
& \frac{d q_{0}}{d t}=\theta \gamma(\hat{m}+\check{m})+\theta \gamma\left(\tilde{M}_{u}\left(x_{0}^{[\nu+1]}+t, x_{1}^{[\nu+1]}\right)\left(x_{1}^{[\nu+1]}-x_{0}^{[\nu+1]}-t\right)-\tilde{M}\left(x_{0}^{[\nu+1]}+t, x_{1}^{[\nu+1]}\right)\right) \\
& \frac{d q_{1}}{d t}=1+\theta \gamma\left(\tilde{M}_{v}\left(x_{0}^{[\nu+1]}, x_{1}^{[\nu+1]}+t\right)\left(x_{1}^{[\nu+1]}+t-x_{0}^{[\nu+1]}\right)+\tilde{M}\left(x_{0}^{[\nu+1]}, x_{1}^{[\nu+1]}+t\right)\right)
\end{aligned}
$$

and

$$
\begin{aligned}
\frac{d q_{N-1}}{d t} & =1-\theta \gamma\left(\tilde{M}_{u}\left(x_{N-1}^{[\nu+1]}+t, x_{N}^{[\nu+1]}\right)\left(x_{N}^{[\nu+1]}-x_{N-1}^{[\nu+1]}-t\right)-\tilde{M}\left(x_{N-1}^{[\nu+1]}+t, x_{N}^{[\nu+1]}\right)\right) \\
\frac{d q_{N}}{d t} & =\theta \gamma(\hat{m}+\check{m})-\theta \gamma\left(\tilde{M}_{v}\left(x_{N-1}^{[\nu+1]}, x_{N}^{[\nu+1]}\right)\left(x_{N}^{[\nu+1]}+t-x_{N-1}^{[\nu+1]}\right)+\tilde{M}\left(x_{N-1}^{[\nu+1]}, x_{N}^{[\nu+1]}\right)\right)
\end{aligned}
$$

Using condition (C2) and the assumption that the perturbed mesh $x^{[\nu+1]}+t e^{i}$ belongs to $D$ it is easy to see that $\frac{d q_{0}}{d t}, \frac{d q_{1}}{d t}, \frac{d q_{N-1}}{d t}, \frac{d q_{N}}{d t}$ are strictly positive. Hence the functions $q_{i}$ are strictly isotone for all $i$, so by Lemma 2.8 the function $G$ defined by 2.4 is a $M$-function.

It follows from Theorems 2.9 and 2.10 that the steady and time-dependent mesh generation problems have unique solutions whenever the grid is fine enough. The existence of solutions, though, must be proved separately. In classical $M$-function theory, one proves the existence of solutions to $F(\mathbf{x})=0$ as follows [41]: First, one needs to construct a super- and a subsolution, i.e., vectors $\hat{\mathbf{x}}$ and $\check{\mathbf{x}}$ such that $\check{\mathbf{x}} \leq \hat{\mathbf{x}}$ and $F(\check{\mathbf{x}}) \leq 0 \leq F(\hat{\mathbf{x}})$. Next, one shows that the nonlinear Gauss-Seidel and Jacobi processes, starting with $\hat{\mathbf{x}}$ or $\check{\mathbf{x}}$, converge monotonically to a fixed point $\mathbf{x}^{*}$ satisfying $F\left(\mathbf{x}^{*}\right)=0$. Thus, the fixed point is a solution of the nonlinear problem and, by uniqueness, its only solution.

In the next section, we will construct super- and subsolutions to the steady and time-dependent mesh generation problems for grids that are fine enough. The existence of solutions then follows from the proof of convergence of a Jacobi or GaussSeidel process. This proof of convergence will be omitted, since it is identical to the one presented in [41]; however, in Section 4, we will show convergence of the parallel and alternating Schwarz methods, which are essentially block variants of the Jacobi and Gauss-Seidel processes.
3. Existence of Solutions. The existence of solutions as well as the convergence of our domain decomposition methods depend on the existence of a supersolution and a subsolution of the initial guess $\mathbf{x}^{0}$, both of which must lie within the domain $D$ of $F$. In both the steady and time-dependent mesh generation problems, the functions

$$
h_{i}\left(u_{i-1}, u_{i}, u_{i+1}\right):=\tilde{M}\left(u_{i-1}, u_{i}\right)\left(u_{i}-u_{i-1}\right)-\tilde{M}\left(u_{i}, u_{i+1}\right)\left(u_{i+1}-u_{i}\right),
$$

for $i=1, \ldots, N$, already defined in (2.3), feature prominently in the nonlinear systems. To prove the existence of super- and subsolutions, we need to know which values can be attained by $h_{i}$ when the $u_{i}$ are not allowed to be too far apart. The following technical lemma provides an answer to this question.

Lemma 3.1. Let $N>B \hat{m} / \check{m}^{2}$, where $B, \hat{m}$ and $\check{m}$ are defined in conditions (C1) and (C2). If $c \in \mathbb{R}$ is such that

$$
\begin{equation*}
\left|u_{R}-u_{L}\right| \frac{\hat{m}}{N}+N|c|<\frac{\check{m}^{2}}{B} \tag{3.1}
\end{equation*}
$$

then there exists a sequence $\left\{u_{i}\right\}_{i=0}^{N}$ such that $u_{0}=u_{L}, u_{N}=u_{R}$, and

$$
h_{i}\left(u_{i-1}, u_{i}, u_{i+1}\right)=c
$$

for $i=1,2, \ldots, N-1$. Moreover, the $u_{i}$ satisfy

$$
\left|u_{i}-u_{i-1}\right| \leq \frac{1}{\check{m}}\left(\left|u_{R}-u_{L}\right| \frac{\hat{m}}{N}+N|c|\right)<\frac{\check{m}}{B}
$$

Proof. Let us consider the flux $\phi_{i-1 / 2}$ between grid points $i-1$ and $i$ given by

$$
\phi_{i-1 / 2}=\tilde{M}\left(u_{i-1}, u_{i}\right)\left(u_{i}-u_{i-1}\right)
$$

Using this notation, we see that $h_{i}(\mathbf{u})=\phi_{i+1 / 2}-\phi_{i-1 / 2}$, so if we fix a parameter $s \in \mathbb{R}$ and require the fluxes $\phi_{i-1 / 2}$ to satisfy

$$
\begin{equation*}
\phi_{i-1 / 2}=\phi_{i-1 / 2}(s)=s+\sigma \chi_{i-1 / 2} \tag{3.2}
\end{equation*}
$$

where $\chi_{i-1 / 2}:=1-\frac{2 i-1}{N}$ is a decreasing sequence between -1 and 1 , then

$$
h_{i}(\mathbf{u}(s))=\phi_{i-1 / 2}(s)-\phi_{i+1 / 2}(s)=\sigma\left(\chi_{i-1 / 2}-\chi_{i+1 / 2}\right)=\frac{2 \sigma}{N} .
$$

So if we choose $\sigma=N c / 2$, then $h_{i}=c$ for all $i$. The extra parameter $s$, despite not having an influence on the value of $h_{i}$, is important because it will eventually allow us to satisfy the boundary conditions $u_{0}=u_{L}, u_{N}=u_{R}$. (This is similar in spirit to shooting methods for solving two-point boundary value problems; here, $s$ plays the role of the initial velocity, which will be chosen to satisfy the boundary condition $u_{N}=u_{R}$.)

Our next task is to show that for any parameter value $s$ satisfying $|s|<\check{m} / B-|\sigma|$, we can indeed choose the fluxes as specified in (3.2), i.e., the recurrence $u_{0}(s)=u_{L}$,

$$
\begin{equation*}
\tilde{M}\left(u_{i-1}(s), u_{i}(s)\right)\left(u_{i}(s)-u_{i-1}(s)\right)=\phi_{i-1 / 2}(s), \quad i=1, \ldots, N \tag{3.3}
\end{equation*}
$$

is solvable for these values of $s$ and the resulting $\mathbf{u}$ is in the set $D$. To show that this is possible, consider two points $t<\tilde{t}$, with $\max \left\{\left|t-u_{i-1}\right|,\left|\tilde{t}-u_{i-1}\right|\right\}<\check{m} / B$. Then

$$
\begin{aligned}
\tilde{M}\left(u_{i-1}, \tilde{t}\right)\left(\tilde{t}-u_{i-1}\right) & -\tilde{M}\left(u_{i-1}, t\right)\left(t-u_{i-1}\right) \\
& =\tilde{M}\left(u_{i-1}, \tilde{t}\right)(\tilde{t}-t)+\left[\tilde{M}\left(u_{i-1}, \tilde{t}\right)-\tilde{M}\left(u_{i-1}, t\right)\right]\left(t-u_{i-1}\right) \\
& \geq \check{m}(\tilde{t}-t)-B(\tilde{t}-t)\left|t-u_{i-1}\right|>0
\end{aligned}
$$

Thus, the function $t \mapsto M\left(u_{i-1}, t\right)\left(t-u_{i-1}\right)$ is strictly increasing within the interval $\left\{t:\left|t-u_{i-1}\right|<\check{m} / B\right\}$, and since

$$
\left|M\left(u_{i-1}, t\right)\left(t-u_{i-1}\right)\right| \geq \check{m}\left|t-u_{i-1}\right|,
$$

its image contains the interval $\left(-\check{m}^{2} / B, \check{m}^{2} / B\right)$. Thus, for

$$
|s|<\check{m}^{2} / B-|\sigma|
$$

we have $\left|\phi_{i-1 / 2}(s)\right|<\check{m}^{2} / B$ for all $i$, so (3.3) is uniquely solvable for $u_{i}(s)$.
Finally, we show that the boundary condition $u_{N}=u_{R}$ can be satisfied by choosing the parameter $s$ appropriately. We have

$$
u_{N}(s)-u_{L}=\sum_{i=1}^{N}\left(u_{i}(s)-u_{i-1}(s)\right)=\sum_{i=1}^{N} \frac{\phi_{i-1 / 2}(s)}{\tilde{M}\left(u_{i-1}(s), u_{i}(s)\right)} .
$$

Substituting the relation (3.2) into $\phi_{i-1 / 2}(s)$ then gives

$$
u_{N}(s)-u_{L}=\sum_{i=1}^{N} \frac{s+\sigma \chi_{i-1 / 2}}{\tilde{M}\left(u_{i-1}(s), u_{i}(s)\right)}
$$

Let $\alpha=\left(u_{R}-u_{L}\right) / N$. If $\alpha \geq 0$, then letting $s_{+}=\alpha \hat{m}+|\sigma|$ leads to

$$
u_{N}\left(s_{+}\right) \geq u_{L}+\sum_{i=1}^{N} \frac{\left(u_{R}-u_{L}\right) \hat{m} / N}{\tilde{M}\left(u_{i-1}(s), u_{i}(s)\right)} \geq u_{L}+u_{R}-u_{L}=u_{R}
$$

and letting $s_{-}=\alpha \check{m}-|\sigma|$ leads to

$$
u_{N}\left(s_{-}\right) \leq u_{L}+\sum_{i=1}^{N} \frac{\left(u_{R}-u_{L}\right) \check{m} / N}{\tilde{M}\left(u_{i-1}(s), u_{i}(s)\right)} \leq u_{L}+u_{R}-u_{L}=u_{R}
$$

If $\alpha<0$ (i.e., if $u_{L}>u_{R}$ ), then a similar calculation holds by letting $s_{+}=\alpha \check{m}+|\sigma|$ and $s_{-}=\alpha \hat{m}-|\sigma|$. Since $u_{N}$ is a continuous function of $s$, there exists $s^{*}$ between $s_{-}$ and $s_{+}$such that $u_{N}\left(s^{*}\right)=u_{R}$, as long as $\left|s_{ \pm}\right|<\check{m}^{2} / B-|\sigma|$. A sufficient condition is

$$
\left|u_{R}-u_{L}\right| \frac{\hat{m}}{N}+N|c|<\frac{\check{m}^{2}}{B}
$$

since this implies

$$
\left|s_{ \pm}\right| \leq\left|u_{R}-u_{L}\right| \frac{\hat{m}}{N}+\frac{N|c|}{2}<\frac{\check{m}^{2}}{B}-\frac{N|c|}{2}=\frac{\check{m}^{2}}{B}-|\sigma| .
$$

Finally, we have

$$
\left|u_{i}-u_{i-1}\right|=\frac{\left|\phi_{i-1 / 2}\right|}{\tilde{M}\left(u_{i-1}, u_{i}\right)} \leq \frac{s_{+}+|\sigma|}{\check{m}}<\frac{\check{m}}{B}
$$

so $\mathbf{u} \in D$, i.e., the solution stays in the region where $h$ is an $M$-function.
Lemma 3.1 immediately implies the existence of solutions to the steady-state problem when the grid is fine enough.

Theorem 3.2. Let $N>B \hat{m} / \check{m}^{2}$. Then there exists a vector $\mathbf{x} \in D$ that satisfies the steady-state equidistribution problem $H(\mathbf{x})=0$ defined in (1.2).

Proof. Letting $c=0, u_{L}=0$ and $u_{R}=1$ in Lemma 3.1 gives the required solution, with $x_{i}=u_{i}$ for $i=0,1, \ldots, N$. $\square$

For the time-dependent case, we can reformulate (1.3) as

$$
\begin{align*}
g_{i}\left(u_{i-1}, u_{i}, u_{i+1}\right)=\left(u_{i}-x_{i}^{[\nu]}\right) & -\theta \gamma\left(\phi_{i+1 / 2}\left(u_{i}, u_{i+1}\right)-\phi_{i-1 / 2}\left(u_{i-1}, u_{i}\right)\right) \\
& -(1-\theta) \gamma\left(\phi_{i+1 / 2}\left(x_{i}^{[\nu]}, x_{i+1}^{[\nu]}\right)-\phi_{i-1 / 2}\left(x_{i-1}^{[\nu]}, x_{i}^{[\nu]}\right)\right)=0 \tag{3.4}
\end{align*}
$$

where $\gamma=\Delta t /\left(\tau \Delta \xi^{2}\right)$. The above equation is to be solved for $u_{i}=x_{i}^{[\nu+1]}, i=$ $1, \ldots, N-1$.

To guarantee the existence of a solution to this problem, we need to make an assumption on $\mathbf{x}^{[\nu]}$, the solution at the previous time step. We assume that there exists a constant $\kappa$ (the "curvature") such that the residuals $r_{i}:=g_{i}\left(x_{i-1}^{[\nu]}, x_{i}^{[\nu]}, x_{i+1}^{[\nu]}\right)=$ $h_{i}\left(x_{i-1}^{[\nu]}, x_{i}^{[\nu]}, x_{i+1}^{[\nu]}\right)$ satisfy

$$
\begin{equation*}
\left|r_{i}\right| \leq \kappa / N^{2} \tag{3.5}
\end{equation*}
$$

This is true if, for instance, $x_{i}^{[\nu]}$ are equidistant samples of a smooth function $x^{[\nu]}(\xi)$ : in this case, we can write $r_{i}=\gamma(\psi(1)-\psi(0))=\gamma \psi^{\prime}(\eta)$ with $\eta \in(0,1)$, where
$\psi(t)=\tilde{M}\left(t x_{i-1}^{[\nu]}+(1-t) x_{i}^{[\nu]}, t x_{i}^{[\nu]}+(1-t) x_{i+1}^{[\nu]}\right)\left[t\left(x_{i}^{[\nu]}-x_{i-1}^{[\nu]}\right)+(1-t)\left(x_{i+1}^{[\nu]}-x_{i}^{[\nu]}\right)\right]$.
Computing directly we have

$$
\begin{aligned}
\psi^{\prime}(\eta)= & {\left[\tilde{M}_{u}(\cdot)\left(x_{i-1}^{[\nu]}-x_{i}^{[\nu]}\right)+\tilde{M}_{v}(\cdot)\left(x_{i}^{[\nu]}-x_{i+1}^{[\nu]}\right)\right]\left[\eta\left(x_{i}^{[\nu]}-x_{i-1}^{[\nu]}\right)+(1-\eta)\left(x_{i+1}^{[\nu]}-x_{i}^{[\nu]}\right)\right] } \\
& +M(\cdot)\left[2 x_{i}^{[\nu]}-x_{i-1}^{[\nu]}-x_{i+1}^{[\nu]}\right] .
\end{aligned}
$$

If $\left|x_{\xi}\right|<L_{x}$ and $\left|x_{\xi \xi}\right|<K_{x}$, then

$$
\left|\psi^{\prime}(\eta)\right| \leq \frac{1}{N^{2}}\left(2 B L_{x}^{2}+\hat{m} K_{x}\right)
$$

That is, (3.5) holds with

$$
\kappa=\gamma\left(2 B L_{x}^{2}+\hat{m} K_{x}\right)
$$

THEOREM 3.3. Let $N \geq B(\hat{m}+\kappa) / \check{m}^{2}$ where $B$, $\hat{m}$, and $\check{m}$ are defined in conditions (C1) and (C2) and $\kappa$ is defined in (3.5). Then for $\frac{1}{2} \leq \theta \leq 1$, there exists a unique solution $\left\{u_{i}\right\}_{i=0}^{N}$ to (3.4).

Proof. Assuming (3.5) holds, we can construct a supersolution for the timedependent problem. Let $c=\kappa / N^{2}, u_{0}=0, u_{N}=1$, and $\hat{\mathbf{u}}$ be the vector constructed from this value using Lemma 3.1. For this to work, we need

$$
\frac{\check{m}^{2}}{B}>\frac{\hat{m}}{N}+N|c|=\frac{\hat{m}+\kappa}{N}
$$

or

$$
N>\frac{B(\hat{m}+\kappa)}{\check{m}^{2}}
$$

Then we have

$$
h_{i}(\hat{\mathbf{u}})=c=\frac{\kappa}{N^{2}} \geq r_{i}=g_{i}\left(\mathbf{x}^{[\nu]}\right)=h_{i}\left(\mathbf{x}^{[\nu]}\right) .
$$

Since $H$ is an $M$-function, this implies $\hat{u}_{i} \geq x_{i}^{[\nu]}$ for all $i$. Thus, we have

$$
g_{i}(\hat{\mathbf{u}}) \geq \theta h_{i}(\hat{\mathbf{u}})+(1-\theta) r_{i} \geq(1-\theta)\left(\kappa / N^{2}-\left|r_{i}\right|\right) \geq 0
$$

where we have used the fact that $\theta \geq 1-\theta$, since $\theta \in\left[\frac{1}{2}, 1\right]$. Thus, $\hat{\mathbf{u}}$ is a supersolution. Similarly, we can choose $c=-\kappa / N^{2}$ to construct a subsolution ǔ, with the same condition on $N$ as above. Now, using $\hat{\mathbf{u}}$ or $\check{\mathbf{u}}$ as the initial guess, one can show [41, §13.5] that the nonlinear Jacobi/Gauss-Seidel process converges to a vector u satisfying $G(\mathbf{u})=0$, so the solution exists (and is unique by Theorem 2.6).

REmARK. Since $g_{i}\left(\mathbf{x}^{[\nu]}\right)=h_{i}\left(\mathbf{x}^{[\nu]}\right)=r_{i}$, we in fact have

$$
g_{i}(\hat{\mathbf{u}}) \geq \theta h_{i}(\hat{\mathbf{u}})+(1-\theta) r_{i} \geq \theta h_{i}\left(\mathbf{x}^{[\nu]}\right)+(1-\theta) r_{i}=r_{i},
$$

which implies $g_{i}(\hat{\mathbf{u}}) \geq g_{i}\left(\mathbf{x}^{[\nu]}\right)$. Similarly, we have $g_{i}(\check{\mathbf{u}}) \leq g_{i}\left(\mathbf{x}^{[\nu]}\right)$. These inequalities will be used in the next section to show that the alternating and parallel Schwarz methods converge when $\mathbf{x}^{[\nu]}$ is used as an initial guess.


Fig. 4.1. Decomposition into many subdomains.
4. Convergence of Domain Decomposition Methods. In this section, we consider domain decomposition algorithms for solving the equidistribution problems (1.2) and (1.3). Each system may be written as $F(\mathbf{x})=0$, where $F: D \subset \mathbb{R}^{N+1} \rightarrow$ $\mathbb{R}^{N+1}$ is given by $F=\left(f_{0}, \ldots, f_{N}\right)^{T}$, where $f_{i}=h_{i}$ from (2.3) in the steady case and $f_{i}=g_{i}$ from (2.4) in the time dependent case. We have seen in Theorems 2.9 and 2.10 that in each case $F$ is an $M$-function whenever the distance between successive points is not too large.

We now decompose the set of nodes $\{0, \ldots, N\}$ into $J$ overlapping subdomains as shown in Figure 4.1, i.e., we have $\Omega_{j}=\left\{i: a_{j}<i<b_{j}\right\}$, with $\xi_{a_{1}}=0, \xi_{b_{J}}=1$ and

$$
a_{j}<b_{j-1}<a_{j+1}<b_{j} \quad \text { for } 2 \leq j \leq J-1
$$

The subdomain problem associated with $\Omega_{j}$ is now defined as follows: we solve $F_{j}\left(\mathbf{x}^{j} ; x_{a_{j}}, x_{b_{j}}\right)=0$ for $\mathbf{x}^{j}$, where

$$
\begin{equation*}
F_{j}=\left(f_{a_{j}+1}, \ldots, f_{b_{j}-1}\right)^{T}, \quad \mathbf{x}^{j}=\left(x_{a_{j}+1}, \ldots, x_{b_{j}-1}^{T}\right) \tag{4.1}
\end{equation*}
$$

and the boundary values at $a_{j}$ and $b_{j}$ are given by $x_{a_{j}}$ and $x_{b_{j}}$. Note that since $F_{j}$ is a restriction of $F$, the functions $F_{j}\left(\cdot ; x_{a_{j}}, x_{b_{j}}\right): \mathbb{R}^{b_{j}-a_{j}-1} \rightarrow \mathbb{R}^{b_{j}-a_{j}-1}$ are $M$-functions. Moreover, the subdomain functions $F_{j}$ inherit the antitonicity property with respect to the boundary values, i.e., we have

$$
\begin{equation*}
F_{j}\left(\mathbf{x}^{j} ; \alpha, \beta\right) \geq F_{j}\left(\mathbf{x}^{j} ; \alpha^{\prime}, \beta^{\prime}\right) \tag{4.2}
\end{equation*}
$$

whenever $\alpha \leq \alpha^{\prime}$ and $\beta \leq \beta^{\prime}$.
We can now define the classical alternating Schwarz (or Gauss-Seidel) algorithm: for $k=0,1,2, \ldots$ and for $j=1, \ldots, J$ solve

$$
\begin{equation*}
F_{j}\left(\mathbf{x}^{j, k+1} ; x_{a_{j}}^{j-1, k+1}, x_{b_{j}}^{j+1, k}\right)=0 \quad \text { for } \mathbf{x}^{j, k+1} \tag{4.3}
\end{equation*}
$$

Similarly, we can define the classical parallel Schwarz (or Jacobi) algorithm: for $k=0,1,2, \ldots$ and for $j=1, \ldots, J$, solve

$$
\begin{equation*}
F_{j}\left(\mathbf{x}^{j, k+1} ; x_{a_{j}}^{j-1, k}, x_{b_{j}}^{j+1, k}\right)=0 \quad \text { for } \mathbf{x}^{j, k+1} \tag{4.4}
\end{equation*}
$$

Under the assumption of the existence of an appropriate subsolution and supersolution, the goal of this section is to prove that both iterations converge to the same unique global solution of (4.1). Thus, these DD iterations can be used in place of Newton's method, especially when the latter has trouble converging.
4.1. Convergence of the Alternating Schwarz iteration. We first prove the convergence of the alternating Schwarz iteration. We will show this by adapting the technique of [44] to handle the subdomain solves.

THEOREM 4.1. Let $N$ be large enough so that the global problem possesses a unique solution $\mathbf{x}^{*} \in D \subset \mathbb{R}^{N+1}$. Let $\mathbf{x}^{0} \in D \subset \mathbb{R}^{N+1}$ be a supersolution, i.e.,
$F\left(\mathbf{x}^{0}\right) \geq 0=F\left(\mathbf{x}^{*}\right)$. Then the subdomain iterates $\mathbf{x}^{j, k}$ obtained from the classical alternating Schwarz method (4.3) with initial guess $\mathbf{x}^{j, 0}=\left.\mathbf{x}^{0}\right|_{\Omega_{j}}$ are well defined and satisfy

$$
\mathbf{x}^{j, k} \geq \mathbf{x}^{j, k+1} \geq\left.\mathbf{x}^{*}\right|_{\Omega_{j}}
$$

for all $k \geq 0$ and $\left.\mathbf{x}^{j, k} \rightarrow \mathbf{x}^{*}\right|_{\Omega_{j}}$ as $k \rightarrow \infty$; that is we have convergence from above. Similarly, if $\mathbf{x}^{0}$ is a subsolution, then $\left.\mathbf{x}^{j, k} \uparrow \mathbf{x}^{*}\right|_{\Omega_{j}}$ as $k \rightarrow \infty$; that is we have convergence from below.

Proof. We only prove the supersolution case here, the subsolution case being similar. For any $1 \leq j \leq J$ and $k \geq 0$, we define the fractional iterate $\mathbf{x}^{k+j / J} \in \mathbb{R}^{N+1}$ as

$$
\mathbf{x}_{i}^{k+j / J}= \begin{cases}x_{i}^{\ell, k+1}, & a_{\ell} \leq i<a_{\ell+1}, \quad \ell<j  \tag{4.5}\\ x_{i}^{j, k+1}, & a_{j} \leq i<b_{j} \\ x_{i}^{\ell, k}, & \max \left\{b_{j}, a_{\ell}\right\} \leq i<a_{\ell+1}, \quad \ell>j\end{cases}
$$

In other words, each successive fractional iterate is obtained by replacing, after each subdomain solve, the portion of the global iterate corresponding to $\Omega_{j}$ by $\mathbf{x}^{j, k+1}$. We note that $\mathbf{x}^{k+j / J}$ is well defined, i.e., the subdomain problem (4.3) has a unique solution on $\Omega_{j}$, whenever $\mathbf{x}^{k+(j-1) / J}$ is a supersolution in $D$, since $F_{j}$ is an $M$-function on $D$ and the restriction of $\mathbf{x}^{k+(j-1) / J}$ onto $\Omega_{j}$ is a supersolution with respect to $F_{j}$. Moreover, the solution of the subdomain problem satisfies $\left|x_{i}^{j, k+1}-x_{i-1}^{j, k}\right|<\check{m} / B$, so that $\mathbf{x}^{k+j / J} \in D$.

Thus, it suffices to show inductively for all $1 \leq j \leq J$ and $k \geq 0$ that $F\left(\mathbf{x}^{k+j / J}\right) \geq$ 0 and

$$
\mathbf{x}^{k+(j-1) / J} \geq \mathbf{x}^{k+j / J} \geq \mathbf{x}^{*}
$$

the conclusion then follows by restricting $\mathbf{x}^{k+j / J}$ to $\Omega_{j}$. The assumption that $\mathbf{x}^{0}$ is a supersolution means that the base case $k=0$ and $j=0$ is verified, so assume inductively that $F\left(\mathrm{x}^{k+(j-1) / J}\right) \geq 0$ for some $1 \leq j \leq J$ and $k \geq 0$. By the induction hypothesis, we have

$$
F_{j}\left(\left.\mathbf{x}^{k+(j-1) / J}\right|_{\Omega_{j}} ; x_{a_{j}}^{k+1}, x_{b_{j}}^{k}\right) \geq 0
$$

whereas by construction, we have

$$
F_{j}\left(\left.\mathbf{x}^{k+j / J}\right|_{\Omega_{j}} ; x_{a_{j}}^{k+1}, x_{b_{j}}^{k}\right)=0
$$

Since $F_{j}\left(\cdot x_{a_{j}}^{k+1}, x_{b_{j}}^{k}\right)$ is an $M$-function, it is inverse isotone, which implies $\left.\mathbf{x}^{k+(j-1) / J}\right|_{\Omega_{j}} \geq$ $\left.\mathbf{x}^{k+j / J}\right|_{\Omega_{j}}$. But $\mathbf{x}^{k+(j-1) / J}$ and $\mathbf{x}^{k+j / J}$ are identical everywhere outside $\Omega_{j}$; so $\mathbf{x}^{k+(j-1) / J} \geq$ $\mathbf{x}^{k+j / J}$. Moreover, we have

$$
f_{i}\left(\mathbf{x}^{k+j / J}\right)= \begin{cases}f_{i}\left(\mathbf{x}^{k+(j-1) / J}\right) \geq 0 & \text { for } i<a_{j} \text { and for } i>b_{j} \\ 0 & \text { for } a_{j}<i<b_{j}\end{cases}
$$

so it remains to show that $f_{i}\left(\mathbf{x}^{k+j / J}\right) \geq 0$ for $i \in\left\{a_{j}, b_{j}\right\}$. But since $x_{i}^{k+(j-1) / J}=$ $x_{i}^{k+j / J}$ for $i \in\left\{a_{j}, b_{j}\right\}$ and $\mathbf{x}^{k+(j-1) / J} \geq \mathbf{x}^{k+j / J}$, the off-diagonal antitonicity of $F$ implies

$$
0 \leq f_{i}\left(\mathbf{x}^{k+(j-1) / J}\right) \leq f_{i}\left(\mathbf{x}^{k+j / J}\right), \quad i \in\left\{a_{j}, b_{j}\right\}
$$

So $F\left(\mathbf{x}^{k+j / J}\right) \geq 0$, as required. Thus, we have shown by induction $\left\{\mathbf{x}^{k+j / J}\right\}_{j, k}$ is a decreasing sequence bounded below by $\mathbf{x}^{*}$, so it converges to some global solution $\tilde{\mathbf{x}}^{*}$. But for any node $0 \leq i \leq N+1$, there is always a subsequence $\left\{r_{m}^{(i)}\right\}_{m \geq 0}$ for which $f_{i}\left(\mathbf{x}_{m}^{r_{m}^{(i)}}\right)=0$ for all $m \geq 0$, so the continuity of $F$ implies $F\left(\tilde{\mathbf{x}}^{*}\right)=0$. Since $F$ is an $M$-function, this implies $\tilde{\mathbf{x}}^{*}=\mathbf{x}^{*}$, so the $\mathbf{x}^{j, k}$ all converge to $\left.\mathbf{x}^{*}\right|_{\Omega_{j}}$. $\square$

As we will show in Section 6, the supersolution constructed from Lemma 3.1 is generally very far from the solution to the global problem and should not be used as the initial guess in an actual computation. A much better choice would be to use the uniform initial guess $\mathbf{x}^{0}=(0,1 / N, \ldots,(N-1) / N, 1)^{T}$. We show that the alternating Schwarz method also converges to the exact solution with this initial guess.

For the steady-state problem, we would like to find a supersolution as well as a subsolution for the uniform initial guess. There, the residual function is simply $h_{i}$, so by putting $c=2 B / N^{2}$, Lemma 3.1 gives a $\mathbf{u}$ that satisfies $h_{i}(\mathbf{u}) \geq h_{i}\left(\mathbf{x}^{0}\right)$, provided (3.1) holds, i.e.

$$
\begin{equation*}
\frac{\hat{m}}{N}+\frac{2 B}{N}<\frac{\check{m}^{2}}{B} \Longrightarrow N>\frac{B(\hat{m}+2 B)}{\check{m}^{2}} \tag{4.6}
\end{equation*}
$$

In other words, $\mathbf{u}$ is a supersolution. Similarly, putting $c=-2 B / N^{2}$ gives a subsolution, provided the same condition on $N$ holds. We now arrive at our first main result.

Theorem 4.2. Let $N>\frac{B(\hat{m}+2 B)}{\check{m}^{2}}$ with $B, \hat{m}$ and $\check{m}$ defined in (C1) and (C2). Suppose that for the uniform guess $\mathbf{x}^{0}=(1 / N, \ldots,(N-1) / N)^{T}$, there are enough points in each subdomain $\Omega_{j}, j=1, \ldots, N$, so that the alternating Schwarz method (4.3) applied to the nonlinear system $H$ in (2.3) produces well-defined subdomain iterates $\mathbf{x}^{j, k}$ for all $j$ and $k$. Then the subdomain iterates converge to $\left.\mathbf{x}^{*}\right|_{\Omega_{j}}$, the restriction of the exact solution onto $\Omega_{j}$.

Proof. First, $N$ is necessarily larger than the $B \hat{m} / \check{m}^{2}$ required for the existence of a solution by Theorem 3.2. In addition, the initial guess $\mathbf{x}^{0}$ belongs to the set $D$, since

$$
x_{i}-x_{i-1}=\frac{1}{N}<\frac{\check{m}^{2}}{B(\hat{m}+2 B)} \leq \frac{\check{m}}{B}
$$

Moreover, the initial residual satisfies

$$
\left|h_{i}\left(\mathbf{x}^{0}\right)\right|=\left|\tilde{M}\left(\frac{i-1}{N}, \frac{i}{N}\right) \cdot \frac{1}{N}-\tilde{M}\left(\frac{i}{N}, \frac{i+1}{N}\right) \cdot \frac{1}{N}\right| \leq \frac{2 B}{N^{2}} .
$$

As a result, by letting $c= \pm 2 B / N^{2}$, we get

$$
\frac{\hat{m}}{N}+N|c|=\frac{\hat{m}+2 B}{N}<\frac{\check{m}^{2}}{B}
$$

so the hypotheses of Lemma 3.1 are satisfied. We thus conclude that there exist subsolutions and supersolutions $\check{\mathbf{x}}^{0}, \hat{\mathbf{x}}^{0} \in D$ such that

$$
H\left(\check{\mathbf{x}}^{0}\right) \leq H\left(\mathbf{x}^{0}\right) \leq H\left(\hat{\mathbf{x}}^{0}\right)
$$

which implies

$$
\begin{equation*}
\check{\mathbf{x}}^{0} \leq \mathbf{x}^{0} \leq \hat{\mathbf{x}}^{0} \tag{4.7}
\end{equation*}
$$

We now prove inductively that for the alternating Schwarz (Gauss-Seidel) process, we have $\check{\mathbf{x}}^{j, k} \leq \mathbf{x}^{j, k} \leq \hat{\mathbf{x}}^{j, k}$ for $j=1, \ldots, J$ and for $k \geq 1$; this implies $\mathbf{x}^{j, k} \rightarrow \mathbf{x}^{*} \mid \Omega_{j}$, since both $\check{\mathbf{x}}^{j, k}$ and $\hat{\mathbf{x}}^{j, k}$ converge to $\left.\mathbf{x}^{*}\right|_{\Omega_{j}}$ by Theorem 4.1.

As an induction hypothesis, assume that there exists some $k \geq 0$ and $\ell \geq 1$ such that

$$
\begin{align*}
\check{\mathbf{x}}^{j, k} & \leq \mathbf{x}^{j, k} \leq \hat{\mathbf{x}}^{j, k} & & 1 \leq j \leq J,  \tag{4.8}\\
\check{\mathbf{x}}^{j, k+1} & \leq \mathbf{x}^{j, k+1} \leq \hat{\mathbf{x}}^{j, k+1} & & 1 \leq j \leq \ell-1, \tag{4.9}
\end{align*}
$$

where for $\ell=1$ the set in (4.9) is empty. Clearly, (4.8) and (4.9) hold for $k=0$ and $\ell=1$. We prove (4.9) for $\ell+1$ as follows: by (4.3), we know that

$$
\begin{aligned}
& H_{\ell}\left(\check{\mathbf{x}}^{\ell, k+1} ; \check{x}_{a_{\ell}-1, k+1}, \check{x}_{b_{\ell}^{\ell+1, k}}^{\ell+1, k}=0,\right. \\
& H_{\ell}\left(\mathbf{x}^{\ell, k+1} ; x_{a_{\ell}, k+1}^{\ell-1, x_{\ell}^{\ell, k}}\right)=0, \\
& H_{\ell}\left(\hat{\mathbf{x}}^{\ell, k+1} ; \hat{x}_{a_{\ell}}^{\ell-1, k+1}, \hat{x}_{b_{\ell}}^{\ell+1, k}\right)=0,
\end{aligned}
$$

where, using the notation established in Section $4, H_{\ell}$ refers to the $\ell$-th subdomain problem for $H$. For convenience we write the boundary conditions as

$$
\begin{aligned}
& \check{x}_{a_{1}}^{0, k+1}=x_{a_{1}}^{0, k+1}=\hat{x}_{a_{1}}^{0, k+1}=0, \\
& \check{x}_{b_{J}}^{J+1, k}=x_{b_{J}}^{J+1, k}=\hat{x}_{b_{J}}^{J+1, k}=1 .
\end{aligned}
$$

Now by (4.2), we have
$0=H_{\ell}\left(\mathbf{x}^{\ell, k+1} ; x_{a_{\ell}}^{\ell-1, k+1}, x_{b_{\ell}}^{\ell+1, k}\right)=H_{\ell}\left(\hat{\mathbf{x}}^{\ell, k+1} ; \hat{x}_{a_{\ell}}^{\ell-1, k+1}, \hat{x}_{b_{\ell}}^{\ell+1, k}\right) \leq H_{\ell}\left(\hat{\mathbf{x}}^{\ell, k+1} ; x_{a_{\ell}}^{\ell-1, k+1}, x_{b_{\ell}}^{\ell+1, k}\right)$.
Since $H_{\ell}$ is inverse isotone when the boundary conditions are fixed, this implies $\mathbf{x}^{\ell, k+1} \leq \hat{\mathbf{x}}^{\ell, k+1}$, and the same argument holds for the subsolution. Thus, (4.9) is verified for $\ell+1$. This completes the proof of (4.9), and hence of (4.8), which in turn shows convergence of $\mathbf{x}^{j, k}$ to the solution $\mathbf{x}^{*} \mid \Omega_{j}$.
Remarks.

1. Suppose each subdomain has at least $B(\hat{m}+2 B) / \check{m}^{2}$ nodes, so that the first iteration of the alternating Schwarz method is defined. Then all subsequent iterates are guaranteed to be defined: since each subdomain solution is necessarily monotonic by the local equidistribution principle, see Lemma B.1, we have $0 \leq u_{i} \leq 1$ for all $i$. Thus, all interface values are between 0 and 1 at the next iteration, so the same number of mesh points suffices. In actual computations, however, one generally needs fewer grid points, see Section 6.
2. If a non-uniform initial guess is used, a result analogous to Theorem 4.2 also holds, except we need to replace $2 B / N^{2}$ by $\max _{i}\left|h_{i}\left(\mathrm{x}^{0}\right)\right|$, and modify the number of grid points in (4.6) accordingly.
3. Remark 2 is important within the $M P$ iteration procedure which alternately solves the physical and mesh DEs. The previous approximate mesh can be used as a good initial guess for the DD solution of next mesh DE solve. This generally reduces the number of DD iterations as the $M P$ iteration proceeds. We provide an example of this in Section 6.
4.2. Convergence of the Parallel Schwarz iteration. Under the same hypothesis as in Theorem 4.1, we can show that the parallel Schwarz method also converges to the exact solution.

ThEOREM 4.3. Let $N$ be large enough so that a unique solution $\mathbf{x}^{*} \in D \subset \mathbb{R}^{N+1}$ of the global problem exists. Let $\mathbf{x}^{0} \in D \subset \mathbb{R}^{N+1}$ be a supersolution, i.e., $F\left(\mathbf{x}^{0}\right) \geq 0=$ $F\left(\mathbf{x}^{*}\right)$. Then the subdomain iterates $\mathbf{x}^{j, k}$ obtained from the classical parallel Schwarz method (4.4) with initial guess $\mathbf{x}^{j, 0}=\left.\mathbf{x}^{0}\right|_{\Omega_{j}}$ are well defined and satisfy

$$
\mathbf{x}^{j, k} \geq \mathbf{x}^{j, k+1} \geq\left.\mathbf{x}^{*}\right|_{\Omega_{j}}
$$

for all $k \geq 0$ and $\left.\mathbf{x}^{j, k} \rightarrow \mathbf{x}^{*}\right|_{\Omega_{j}}$ as $k \rightarrow \infty$; that is, we have convergence from above. Similarly, if $\mathbf{x}^{0}$ is a subsolution, then $\left.\mathbf{x}^{j, k} \uparrow \mathbf{x}^{*}\right|_{\Omega_{j}}$ as $k \rightarrow \infty$; that is, we have convergence from below.

Proof. The arguments for the supersolution case are similar to the proof of Theorem 4.1, but instead of defining the fractional iterates, we define $\mathbf{x}^{k+1}$ by

$$
x_{i}^{k+1}= \begin{cases}x_{i}^{j, k+1}, & \text { if } b_{j-1} \leq i \leq a_{j+1} \\ \min \left\{x_{i}^{j-1, k+1}, x_{i}^{j, k+1}\right\}, & \text { if } a_{j}<i<b_{j-1}\end{cases}
$$

In other words, for points inside the overlap, we choose the smaller of the two solutions. We now show by induction that $F\left(\mathbf{x}^{k}\right) \geq 0$ and $\mathbf{x}^{*} \leq \mathbf{x}^{k+1} \leq \mathbf{x}^{k}$ for all $k \geq 0$. The base case for $F\left(\mathbf{x}^{k}\right) \geq 0$ is true for $k=0$. Then we have

$$
F_{j}\left(\left.\mathbf{x}^{k}\right|_{\Omega_{j}} ; x_{a_{j}}^{k}, x_{b_{j}}^{k}\right) \geq 0
$$

for all $j$. For each index $i$, let $\Omega_{j}$ be the subdomain for which $x_{i}^{k+1}=x_{i}^{j, k+1}$. (The choice of $j$ is unambiguous when $i$ is outside the overlap, but when $i$ is inside the overlap, we pick the $j$ that yields the smaller value for $x_{i}^{j, k+1}$.) Then by construction, we have

$$
f_{i}\left(\mathbf{x}^{j, k+1}\right)=0 \leq f_{i}\left(\mathbf{x}^{k}\right)
$$

By the inverse isotonicity of $F_{j}$, this implies $\mathbf{x}^{j, k+1} \leq\left.\mathbf{x}^{k}\right|_{\Omega_{j}}$. Now $\left.\mathbf{x}^{k+1}\right|_{\Omega_{j}} \leq \mathbf{x}^{j, k+1}$, since $\mathbf{x}^{k+1}$ is defined as a pointwise minimum. By off-diagonal antitonicity of $f_{i}$, we see that

$$
0=f_{i}\left(\mathbf{x}^{j, k+1}\right) \leq f_{i}\left(\mathrm{x}^{k+1}\right)
$$

Since this is true for all $i$, we conclude that $F\left(\mathbf{x}^{k+1}\right) \geq 0$, as well as

$$
\mathbf{x}^{*} \leq \mathbf{x}^{k+1} \leq \mathbf{x}^{k}
$$

The argument for the subsolution is similar, except we should take the larger of the two subdomain values inside the overlap.

It is clear that for the steady-state problem, we can use the same super- and subsolution as in Theorem 4.2 to show that the parallel Schwarz method converges for the uniform guess $\mathbf{x}^{0}=(1 / N, \ldots,(N-1) / N)^{T}$. We now state a similar result for the time-dependent problem, which is the second main result of our paper.

THEOREM 4.4. Let $N>B(\hat{m}+\kappa) / \check{m}^{2}$, where $B, \hat{m}$, and $\check{m}$ are defined in conditions (C1) and (C2) and $\kappa$ is defined in (3.5). Suppose each subdomain $\Omega_{j}$
contains enough grid points so that the parallel and alternating Schwarz methods (4.4) and (4.3) applied to the nonlinear system $G$ in (2.4) produce well-defined subdomain iterates. Then for the initial guess $\mathbf{x}^{[\nu]}$, both algorithms converge to the solution $\mathbf{x}^{[\nu+1]}$ of the time-dependent mesh generation problem (3.4).

Proof. The condition on $N$ is the same as in Theorem 3.3, so a unique solution to the global problem exists. Consider once again the supersolution $\hat{\mathbf{u}}$ and subsolution $\mathbf{u}$ constructed in the proof of Theorem 3.3. As shown in the remark immediately after Theorem 3.3, $\hat{\mathbf{u}}$ and $\mathbf{u}$ satisfy

$$
g_{i}(\check{\mathbf{u}}) \leq g_{i}\left(\mathbf{x}^{[\nu]}\right) \leq g_{i}(\hat{\mathbf{u}})
$$

Thus, by the $M$-function property of $G$, we have

$$
\check{\mathbf{u}} \leq \mathbf{x}^{[\nu]} \leq \hat{\mathbf{u}} .
$$

We can now mimic the proof of Theorem 4.2 to show that the parallel and alternating Schwarz methods (4.3) and (4.4) both converge to the unique solution $\mathbf{x}^{*}$ for the initial guess $\mathbf{x}^{[\nu]}$.
REmARK. Suppose the $x_{i}^{[\nu]}$ are all between 0 and 1 , and assume that $2(1-\theta) \gamma \hat{m} \leq 1$. (This is always true for $\theta=1$, i.e., Backward Euler, for instance.) Then the discrete maximum principle (Lemma B.2) holds; as a result, if each subdomain has at least $B(\hat{m}+\kappa) / \check{m}^{2}$ nodes, then every iteration of the alternating and parallel Schwarz methods is well defined, for the same reason as in the steady case.
5. Asymptotic convergence rate of the Schwarz iterations. In this section, we analyze the asymptotic convergence rate of the discrete parallel Schwarz method of Section 4 as $k \rightarrow \infty$, i.e., when $\mathbf{x}^{k}$ is close to the solution $\mathbf{x}^{*}$. We only show the analysis for the stationary equidistribution problem

$$
\phi_{i-1 / 2}\left(x_{i-1}, x_{i}\right)-\phi_{i+1 / 2}\left(x_{i}, x_{i+1 / 2}\right)=0, \quad i=1, \ldots, N-1
$$

similar techniques can be used to analyze the time-dependent problem. Our first main result is the following:

Theorem 5.1. Suppose the conditions (C1)-(C3) hold. Then the parallel Schwarz method (4.4) for solving the stationary equidistribution problem (2.3) converges with the asymptotic convergence factor $\rho(Z)$, where $\rho(Z)$ denotes the spectral radius of the $(2 J-2) \times(2 J-2)$ matrix
$Z=\left[\begin{array}{cc|cc|cc|c|cc}0 & Q_{1} & & & & & & \\ P_{2} & 0 & 0 & R_{2} & & & & & \\ \hline S_{2} & 0 & 0 & Q_{2} & & & & & \\ & & P_{3} & 0 & 0 & R_{3} & & & \\ \hline & & S_{3} & 0 & 0 & Q_{3} & & & \\ & & & P_{4} & 0 & 0 & R_{4} & & \\ \hline & & \ddots & \ddots & & \ddots \\ \hline & & & & S_{J-1} & 0 & 0 & Q_{J-1} \\ & & & & & & & & P_{J} \\ & & & & & & & \end{array}\right]$.

Moreover, the coefficients $P_{j}, Q_{j}, R_{j}$ and $S_{j}$ satisfy the bounds

$$
\begin{aligned}
& \mu^{a_{j+1}-a_{j}} \cdot \frac{\mu^{-\left(b_{j}-a_{j+1}\right)}-1}{\mu^{-\left(b_{j}-a_{j}\right)}-1} \leq S_{j} \leq \mu^{-\left(a_{j+1}-a_{j}\right)} \cdot \frac{1-\mu^{b_{j}-a_{j+1}}}{1-\mu^{b_{j}-a_{j}}} \\
& \mu^{b_{j}-a_{j+1}} \cdot \frac{\mu^{-\left(a_{j+1}-a_{j}\right)}-1}{\mu^{-\left(b_{j}-a_{j}\right)}-1} \leq Q_{j} \leq \mu^{-\left(b_{j}-a_{j+1}\right)} \cdot \frac{1-\mu^{a_{j+1}-a_{j}}}{1-\mu^{b_{j}-a_{j}}} \\
& \mu^{b_{j-1}-a_{j}} \cdot \frac{\mu^{-\left(b_{j}-b_{j-1}\right)}-1}{\mu^{-\left(b_{j}-a_{j}\right)}-1} \leq P_{j} \leq \mu^{-\left(b_{j-1}-a_{j}\right)} \cdot \frac{1-\mu^{b_{j}-b_{j-1}}}{1-\mu^{b_{j}-a_{j}}} \\
& \mu^{b_{j}-b_{j-1}} \cdot \frac{\mu^{-\left(b_{j-1}-a_{j}\right)}-1}{\mu^{-\left(b_{j}-a_{j}\right)}-1} \leq R_{j} \leq \mu^{-\left(b_{j}-b_{j-1}\right)} \cdot \frac{1-\mu^{b_{j-1}-a_{j}}}{1-\mu^{b_{j}-a_{j}}}
\end{aligned}
$$

whenever $\mu \in(0,1]$ is a constant such that for all $i=1, \ldots, N-1$, we have

$$
\begin{equation*}
\mu\left|\frac{\partial \phi_{i-1 / 2}}{\partial x_{i}}\right| \leq\left|\frac{\partial \phi_{i+1 / 2}}{\partial x_{i}}\right| \leq \frac{1}{\mu}\left|\frac{\partial \phi_{i-1 / 2}}{\partial x_{i}}\right| . \tag{5.2}
\end{equation*}
$$

This result is analogous to the result in [22] for the continuous case: there, it is shown that for the continuous equidistribution problem, the error on the $j$ th subdomain with respect to the error measure

$$
E_{j}^{k}(\xi)=\int_{x_{j}(\xi)}^{x_{j}^{k}(\xi)} M(\tilde{x}) d \tilde{x}
$$

satisfies

$$
\begin{equation*}
E_{j}^{k+1}(\xi)=\frac{1}{\xi_{b_{j}}-\xi_{a_{j}}}\left[\left(\xi-\xi_{a_{j}}\right) E_{j+1}^{k}\left(\xi_{b_{j}}\right)+\left(\xi_{b_{j}}-\xi\right) E_{j-1}^{k}\left(\xi_{a_{j}}\right)\right] \tag{5.3}
\end{equation*}
$$

If we now define the vector $\mathbf{E}^{k}$ containing the errors at all interface points (ordered from left to right according to their location in the domain, see Fig. 4.1)

$$
\mathbf{E}^{k}=\left(E_{1}^{k}\left(\xi_{a_{2}}\right), E_{2}^{k}\left(\xi_{b_{1}}\right), E_{2}^{k}\left(\xi_{a_{3}}\right), E_{3}^{k}\left(\xi_{b_{2}}\right), \ldots, E_{J-1}^{k}\left(\xi_{a_{J}}\right), E_{J}^{k}\left(\xi_{b_{J-1}}\right)\right)^{T}
$$

then we can use (5.3) to obtain $\mathbf{E}^{k+1}=\tilde{Z} \mathbf{E}^{k}$, where
$\tilde{Z}=\left[\begin{array}{cc|cc|cc|c|c}0 & q_{1} & & & & & & \\ p_{2} & 0 & 0 & r_{2} & & & & \\ \hline s_{2} & 0 & 0 & q_{2} & & & & \\ & & p_{3} & 0 & 0 & r_{3} & & \\ \hline & & s_{3} & 0 & 0 & q_{3} & & \\ & & & p_{4} & 0 & 0 & r_{4} & \\ \hline & & & \ddots & \ddots & & \ddots \\ \hline & & & & s_{J-1} & 0 & 0 & q_{J-1} \\ & & & & & & p_{J} & 0\end{array}\right]$,
with

$$
p_{j}=\frac{\xi_{b_{j}}-\xi_{b_{j-1}}}{\xi_{b_{j}}-\xi_{a_{j}}}, \quad q_{j}=\frac{\xi_{a_{j+1}}-\xi_{a_{j}}}{\xi_{b_{j}}-\xi_{a_{j}}}, \quad r_{j}=\frac{\xi_{b_{j-1}}-\xi_{a_{j}}}{\xi_{b_{j}}-\xi_{a_{j}}}, \quad s_{j}=\frac{\xi_{b_{j}}-\xi_{a_{j+1}}}{\xi_{b_{j}}-\xi_{a_{j}}}
$$

Our second main result of this section shows that as the grid is refined, the discrete iteration matrix converges to the continuous one above, subject to some smoothness assumptions on $\tilde{M}$ :

ThEOREM 5.2. Let $N$ be the number of degrees of freedom in the equidistribution problem. Suppose the conditions (C1)-(C3) hold, and assume in addition that the partial derivatives $\frac{\partial \tilde{M}}{\partial u}$ and $\frac{\partial \tilde{M}}{\partial v}$ are both Lipschitz continuous. Then the coefficients $P_{j}, Q_{j}, R_{j}$ and $S_{j}$ in Theorem 5.1 satisfy

$$
P_{j} \rightarrow \frac{\xi_{b_{j}}-\xi_{b_{j-1}}}{\xi_{b_{j}}-\xi_{a_{j}}}, \quad Q_{j} \rightarrow \frac{\xi_{a_{j+1}}-\xi_{a_{j}}}{\xi_{b_{j}}-\xi_{a_{j}}}, \quad R_{j} \rightarrow \frac{\xi_{b_{j-1}}-\xi_{a_{j}}}{\xi_{b_{j}}-\xi_{a_{j}}}, \quad S_{j} \rightarrow \frac{\xi_{b_{j}}-\xi_{a_{j+1}}}{\xi_{b_{j}}-\xi_{a_{j}}}
$$

as $N \rightarrow \infty$.
The next two subsections contain the proofs of these two theorems.
5.1. Proof of Theorem 5.1. First, let us write the parallel Schwarz method (4.4) in fixed point form in terms of the interface values. Suppose we arrange the interface values in the same order as in the continuous setting:

$$
\mathbf{X}_{\Gamma}^{k}=\left(x_{a_{2}}^{1, k}, x_{b_{1}}^{2, k}, x_{a_{3}}^{2, k}, x_{b_{2}}^{3, k}, \ldots, x_{a_{J}}^{J-1, k}, x_{b_{J-1}}^{J, k}\right)^{T}
$$

Next, let $\Phi$ denote the fixed point mapping corresponding to the parallel Schwarz iteration, i.e., we have $\mathbf{X}_{\Gamma}^{k+1}=\Phi\left(\mathbf{X}_{\Gamma}^{k}\right)$. Then it is well known that the asymptotic convergence rate of the iteration is given by the spectral radius of the Jacobian $\frac{\partial \Phi}{\partial \mathbf{X}}$. Our first step is to determine the sparsity pattern of this Jacobian matrix.

Lemma 5.3. The Jacobian $\frac{\partial \Phi}{\partial X}$ has the same sparsity pattern as the matrix $Z$ in (5.1).

Proof. If $\mathbf{Y}=\Phi(\mathbf{X})$ with $\mathbf{Y}=\left(y_{a_{2}}, y_{b_{1}}, \ldots, y_{a_{J}}, y_{b_{J-1}}\right)^{T}$, then for all subdomains $j=\{1, \ldots, J\}$, the values $y_{a_{j}}$ and $y_{b_{j}}$ are defined implicitly by

$$
\begin{array}{lll}
y_{a_{j}}=\left(\mathbf{y}^{j-1}\right)_{a_{j}}, & \text { where } \quad H_{j-1}\left(\mathbf{y}^{j-1} ; x_{a_{j-1}}, x_{b_{j-1}}\right)=0 \\
y_{b_{j}}=\left(\mathbf{y}^{j+1}\right)_{b_{j}}, & \text { where } \quad H_{j+1}\left(\mathbf{y}^{j+1} ; x_{a_{j+1}}, x_{b_{j+1}}\right)=0,
\end{array}
$$

where the $H_{j}$ are the subdomain objective functions, as defined in (4.1). Thus, $y_{a_{j}}$ only depends on the values of $x_{a_{j-1}}$ and $x_{b_{j-1}}$, whereas $y_{b_{j}}$ only depends on $x_{a_{j+1}}$, and $x_{b_{j+1}}$, so their derivatives with respect to the remaining variables are identically zero. This yields a Jacobian with the sparsity pattern shown in (5.1).

Our next step is to compute the individual coefficients of the Jacobian. Recall that the global problem is defined by $H(\mathbf{y})=0$, where the components $h_{i}$ of $H$ read

$$
h_{i}\left(y_{i-1}, y_{i}, y_{i+1}\right)=\phi_{i-1 / 2}\left(y_{i-1}, y_{i}\right)-\phi_{i+1 / 2}\left(y_{i}, y_{i+1}\right), \quad a_{j-1}<i<b_{j-1} .
$$

In order to match the results in the continuous case, we will consider a diagonally scaled version of the Jacobian, i.e., we let $Z=D \frac{\partial \Phi}{\partial \mathbf{X}} D^{-1}$, where

$$
D=\operatorname{diag}\left(\frac{\partial \phi_{a_{2}+1 / 2}}{\partial y_{a_{2}}}, \frac{\partial \phi_{b_{1}+1 / 2}}{\partial y_{b_{1}}}, \ldots, \frac{\partial \phi_{a_{J}+1 / 2}}{\partial y_{a_{J}}}, \frac{\partial \phi_{b_{J-1}+1 / 2}}{\partial y_{b_{J-1}}}\right) .
$$

Comparing with the definition of $Z$ in (5.1), we see that

$$
\begin{array}{ll}
P_{j}=\left(\frac{\partial \phi_{a_{j}+1 / 2}}{\partial y_{a_{j}}}\right)^{-1} \frac{\partial y_{b_{j-1}}}{\partial x_{a_{j}}} \frac{\partial \phi_{b_{j-1}+1 / 2}}{\partial y_{b_{j-1}}}, \quad R_{j}=\left(\frac{\partial \phi_{b_{j}+1 / 2}}{\partial y_{b_{j}}}\right)^{-1} \frac{\partial y_{b_{j-1}}}{\partial x_{b_{j}}} \frac{\partial \phi_{b_{j-1}+1 / 2}}{\partial y_{b_{j-1}}}, \\
S_{j}=\left(\frac{\partial \phi_{a_{j}+1 / 2}}{\partial y_{a_{j}}}\right)^{-1} \frac{\partial y_{a_{j+1}}}{\partial x_{a_{j}}} \frac{\partial \phi_{a_{j+1}+1 / 2}}{\partial y_{a_{j+1}}}, \quad Q_{j}=\left(\frac{\partial \phi_{b_{j}+1 / 2}}{\partial y_{b_{j}}}\right)^{-1} \frac{\partial y_{a_{j+1}}}{\partial x_{b_{j}}} \frac{\partial \phi_{a_{j+1}+1 / 2}}{\partial y_{a_{j+1}}}, \tag{5.5}
\end{array}
$$

where all derivatives are to be evaluated at the solution $\mathbf{X}^{*}$.
Using implicit differentiation, we obtain

$$
\frac{\partial H_{j}}{\partial \mathbf{y}^{j}} \frac{\partial \mathbf{y}^{j}}{\partial x_{a_{j}}}+\frac{\partial H_{j}}{\partial x_{a_{j}}}=0 \Longrightarrow \frac{\partial \mathbf{y}^{j}}{\partial x_{a_{j}}}=-\left(\frac{\partial H_{j}}{\partial \mathbf{y}^{j}}\right)^{-1} \frac{\partial H_{j}}{\partial x_{a_{j}}}
$$

Thus, $\frac{\partial y_{a_{j+1}}}{\partial x_{a_{j}}}$ is simply the component of $\frac{\partial \mathbf{y}^{j}}{\partial x_{a_{j}}}$ corresponding to the location $a_{j+1}$, and the other derivatives can be calculated similarly. The matrix $K_{j}:=\frac{\partial H_{j}}{\partial \mathbf{y}^{j}}$ has a special structure: it is a tridiagonal matrix of the form

$$
K_{j}=\left[\begin{array}{ccccc}
c_{1}+d_{1} & -c_{2} & & &  \tag{5.6}\\
-d_{1} & c_{2}+d_{2} & -c_{3} & & \\
& \ddots & \ddots & \ddots & \\
& & -d_{N_{j}-2} & c_{N_{j}-1}+d_{N_{j}-1} & -c_{N_{j}} \\
& & & -d_{N_{j}-1} & c_{N_{j}}+d_{N_{j}}
\end{array}\right]
$$

where $N_{j}=b_{j}-a_{j}-1$ is the number of interior unknowns in subdomain $j$ and

$$
c_{i}=\frac{\partial \phi_{a_{j}+i-1 / 2}}{\partial y_{a_{j}+i}}, \quad d_{i}=-\frac{\partial \phi_{a_{j}+i+1 / 2}}{\partial y_{a_{j}+i}} .
$$

Next, we notice that the vector $\frac{\partial H_{j}}{\partial x_{a_{j}}}$ has exactly one non-zero component at the first position. This is because the only component that depends on $x_{a_{j}}$ corresponds to index $a_{j}+1$, i.e., the first unknown in the nonlinear system. Its value can be calculated explicitly:

$$
\frac{\partial h_{a_{j}+1}}{\partial x_{a_{j}}}=\frac{\partial \phi_{a_{j}+1 / 2}}{\partial x_{a_{j}}}=-d_{0}
$$

Thus, we see that the derivatives we want to compute are multiples of $\left(\mathbf{e}^{q}\right)^{T} K_{j}^{-1} \mathbf{e}^{p}$, where $\mathbf{e}^{q}$ is the $q$-th standard basis vector ( 1 at position $q$ and 0 elsewhere), $q$ is the position of index $a_{j+1}$ in subdomain $j$, and $p \in\left\{1, N_{j}\right\}$ depending on whether we are taking the derivative with respect to $x_{a_{j}}$ or $x_{b_{j}}$. In other words, we need to estimate individual entries of the inverse of a tridiagonal matrix. Estimates of this type exist in the literature, see e.g. [40] and [43]; unfortunately, those estimates are not precise enough for our purposes. In particular, to prove Theorem 5.2, we need the upper and lower bounds in 5.1 to converge to the same value as $\mu \rightarrow 1$. As a result, we will derive our own estimates with the help of the next lemma.

Lemma 5.4. Let $K_{j}$ be defined in (5.6). Then for $q \in\left\{1, \ldots, N_{j}\right\}$, we have

$$
\left(\mathbf{e}^{q}\right)^{T} K_{j}^{-1} \mathbf{e}^{1}=\frac{\mu_{1} \cdots \mu_{q-1}}{C_{q}\left(\frac{c_{q}}{C_{q}}+\frac{d_{q}}{D_{q}}\right)}, \quad\left(\mathbf{e}^{q}\right)^{T} K_{j}^{-1} \mathbf{e}^{N_{j}}=\frac{\mu_{q+1}^{-1} \cdots \mu_{N_{j}}^{-1}}{D_{q}\left(\frac{c_{q}}{C_{q}}+\frac{d_{q}}{D_{q}}\right)},
$$

where $\mu_{i}=d_{i} / c_{i}$ and

$$
\begin{aligned}
& C_{q}=1+\mu_{q-1}+\mu_{q-1} \mu_{q-2}+\ldots+\mu_{q-1} \cdots \mu_{1} \\
& D_{q}=1+\mu_{q+1}^{-1}+\mu_{q+1}^{-1} \mu_{q+2}^{-1}+\ldots+\mu_{q+1}^{-1} \cdots \mu_{N_{j}}^{-1}
\end{aligned}
$$

Proof. Let $\mathbf{w}=K_{j}^{-T} \mathbf{e}^{q}=\left(w_{1}, \ldots, w_{N_{j}}\right)^{T}$, so that for all $p$,

$$
\left(\mathbf{e}^{q}\right)^{T} K_{j}^{-1} \mathbf{e}^{p}=\mathbf{w}^{T} \mathbf{e}^{p}=w_{p}
$$

Then we have $K_{j} \mathbf{w}=\mathbf{e}^{q}$, which in component form reads

$$
c_{i}\left(w_{i}-w_{i-1}\right)-d_{i}\left(w_{i+1}-w_{i}\right)= \begin{cases}1, & i=q  \tag{5.7}\\ 0, & \text { else }\end{cases}
$$

where for convenience we have defined $w_{0}=w_{P+1}=0$. This implies

$$
w_{i}-w_{i-1}=\frac{d_{i}}{c_{i}}\left(w_{i+1}-w_{i}\right) \quad \text { for } i \neq q
$$

which allows us to unroll the recurrence to obtain

$$
\begin{aligned}
w_{q}=w_{q}-w_{0} & =\left(w_{q}-w_{q-1}\right)+\left(w_{q-1}-w_{q-2}\right)+\cdots+\left(w_{1}-w_{0}\right) \\
& =\left(w_{q}-w_{q-1}\right)\left[1+\frac{d_{q-1}}{c_{q-1}}+\frac{d_{q-1} d_{q-2}}{c_{q-1} c_{q-1}}+\cdots+\frac{d_{q-1} \cdots d_{1}}{c_{q-1} \cdots c_{1}}\right]=: C_{q}\left(w_{q}-w_{q-1}\right)
\end{aligned}
$$

Unrolling the recurrence in the other direction gives another expression for $w_{q}$ :

$$
\begin{aligned}
w_{q}=w_{q}-w_{P+1} & =\left(w_{q}-w_{q+1}\right)+\left(w_{q+1}-w_{q+2}\right)+\cdots+\left(w_{N_{j}}-w_{N_{j}+1}\right) \\
& =\left(w_{q}-w_{q+1}\right)\left[1+\frac{c_{q+1}}{d_{q+1}}+\frac{c_{q+1} c_{q+2}}{d_{q+1} d_{q+2}}+\cdots+\frac{c_{q+1} \cdots c_{N_{j}}}{d_{q+1} \cdots d_{N_{j}}}\right]=: D_{q}\left(w_{q}-w_{q+1}\right)
\end{aligned}
$$

Now substituting the two expressions above into the $q$-th equation of (5.7) gives

$$
\left(\frac{c_{q}}{C_{q}}+\frac{d_{q}}{D_{q}}\right) w_{q}=1
$$

Thus, we get

$$
\begin{aligned}
w_{N_{j}}=w_{N_{j}}-w_{N_{j}+1} & =\frac{c_{q+1} \cdots c_{N_{j}}}{d_{q+1} \cdots d_{N_{j}}}\left(w_{q}-w_{q+1}\right) \\
& =\frac{c_{q+1} \cdots c_{N_{j}}}{d_{q+1} \cdots d_{N_{j}}} \frac{w_{q}}{D_{q}}=\frac{c_{q+1} \cdots c_{N_{j}}}{d_{q+1} \cdots d_{N_{j}}} \frac{1}{D_{q}\left(\frac{c_{q}}{C_{q}}+\frac{d_{q}}{D_{q}}\right)} .
\end{aligned}
$$

Similarly, we have

$$
w_{1}=w_{1}-w_{0}=\frac{d_{q-1} \cdots d_{1}}{c_{q-1} \cdots c_{1}} \cdot \frac{1}{C_{q}\left(\frac{c_{q}}{C_{q}}+\frac{d_{q}}{D_{q}}\right)}
$$

The above lemma can be generalized to obtain expressions for the other components of $\mathbf{w}$, but we will omit this since it is not required for our analysis.

Example. For the case $c_{i}=d_{i}=1$ for all $i$ (i.e., the 1D Laplacian), we obtain $C_{q}=q, D_{q}=N_{j}-q+1$, so

$$
w_{1}=\frac{N_{j}-q+1}{N_{j}+1}, \quad w_{N_{j}}=\frac{q}{N_{j}+1}
$$

In particular, we have $w_{1}+w_{N_{j}}=1$.
We are now ready to prove Theorem 5.1.
Proof. (Theorem 5.1) For any invertible diagonal matrix $D$, the matrices $\frac{\partial \Phi}{\partial \mathbf{X}}$ and $D \frac{\partial \Phi}{\partial \mathbf{X}} D^{-1}$ have the same sparsity pattern, so by Lemma 5.3 , it suffices to show that the coefficients $P_{j}, Q_{j}, R_{j}$ and $S_{j}$ satisfy the stated bounds. We prove the bounds in detail for $S_{j}$ only, the others being similar.

Let $\mu \in(0,1]$ be a constant satisfying (5.2), so that $\mu \leq d_{i} / c_{i} \leq \mu^{-1}$ for all $i$. Then for $q$ fixed, the expressions in the statement of Lemma 5.4 imply that

$$
\frac{1-\mu^{q}}{1-\mu} \leq C_{q} \leq \frac{\mu^{-q}-1}{\mu^{-1}-1}, \quad \frac{1-\mu^{N_{j}-q+1}}{1-\mu} \leq D_{q} \leq \frac{\mu^{-\left(N_{j}-q+1\right)}-1}{\mu^{-1}-1}
$$

The denominator of $w_{1}=\left(\mathbf{e}^{q}\right)^{T} K_{j}^{-1} \mathbf{e}^{1}$ then satisfies

$$
\begin{aligned}
C_{q}\left(\frac{c_{q}}{C_{q}}+\frac{d_{q}}{D_{q}}\right) & =c_{q}\left(1+\frac{d_{q}}{c_{q}} \cdot \frac{C_{q}}{D_{q}}\right) \\
& \leq c_{q}\left(1+\mu^{-1} \cdot \frac{\mu^{-q}-1}{\mu^{-1}-1} \cdot \frac{1-\mu}{1-\mu^{N_{j}-q+1}}\right)=c_{q} \frac{\mu^{-\left(N_{j}+1\right)}-1}{\mu^{-\left(N_{j}-q+1\right)}-1} .
\end{aligned}
$$

On the other hand, we also have

$$
C_{q}\left(\frac{c_{q}}{C_{q}}+\frac{d_{q}}{D_{q}}\right) \geq c_{q}\left(1+\mu \cdot \frac{1-\mu^{q}}{1-\mu} \cdot \frac{\mu^{-1}-1}{\mu^{-\left(N_{j}-q+1\right)}-1}\right)=c_{q} \frac{1-\mu^{N_{j}+1}}{1-\mu^{N_{j}-q+1}} .
$$

This implies

$$
\frac{d_{1} \cdots d_{q-1}}{c_{1} \cdots c_{q}} \cdot \frac{\mu^{-\left(N_{j}-q+1\right)}-1}{\mu^{-\left(N_{j}+1\right)}-1} \leq w_{1} \leq \frac{d_{1} \cdots d_{q-1}}{c_{1} \cdots c_{q}} \cdot \frac{1-\mu^{N_{j}-q+1}}{1-\mu^{N_{j}+1}}
$$

Now let $q$ be the local index of $a_{j+1}$ within the $j$ th subdomain, i.e., $q=a_{j+1}-a_{j}$. Then by definition, we have

$$
\frac{\partial y_{a_{j}+1}}{\partial x_{a_{j}}}=d_{0} w_{1}
$$

But $\frac{\partial \phi_{a_{j}+1 / 2}}{\partial y_{a_{j}}}=-d_{0}$ and $\frac{\partial \phi_{a_{j+1}+1 / 2}}{\partial y_{a_{j+1}}}=-d_{q}$, so in fact

$$
S_{j}=\frac{d_{q}}{d_{0}} \cdot d_{0} w_{1}=d_{q} w_{1},
$$

from which it follows that

$$
\frac{d_{1} \cdots d_{q}}{c_{1} \cdots c_{q}} \cdot \frac{\mu^{-\left(N_{j}-q+1\right)}-1}{\mu^{-\left(N_{j}+1\right)}-1} \leq S_{j} \leq \frac{d_{1} \cdots d_{q}}{c_{1} \cdots c_{q}} \cdot \frac{1-\mu^{N_{j}-q+1}}{1-\mu^{N_{j}+1}}
$$

The bound on the ratios $d_{i} / c_{i}$ then yields

$$
\mu^{q} \cdot \frac{\mu^{-\left(N_{j}-q+1\right)}-1}{\mu^{-\left(N_{j}+1\right)}-1} \leq S_{j} \leq \mu^{-q} \cdot \frac{1-\mu^{N_{j}-q+1}}{1-\mu^{N_{j}+1}}
$$

Substituting $N_{j}=b_{j}-a_{j}-1$ and $q=a_{j+1}-a_{j}$ yields the required bounds on $S_{j}$. The bounds on $P_{j}, Q_{j}$ and $R_{j}$ are obtained similarly.
5.2. Proof of Theorem 5.2. The bounds in Theorem 5.1 depend on the value of $\mu$ that can be used. To understand the behavior of the algorithm as the number of grid points increases, we need the following lemma.

Lemma 5.5. Let $K_{j}$ be defined in (5.6) and $N$ be the number of degrees of freedom in the equidistribution problem. Suppose the conditions (C1)-(C3) hold, and assume in addition that the partial derivatives $\frac{\partial \tilde{M}}{\partial u}$ and $\frac{\partial \tilde{M}}{\partial v}$ are both Lipschitz continuous. Then for $N$ large enough, there exists a constant $\omega>0$ independent of $N$ such that if we let

$$
\mu=1-\frac{\omega}{N^{2}}
$$

then we have $\mu \leq d_{i} / c_{i} \leq \mu^{-1}$ for all $i$.
Proof. Let $\left(y_{i}\right)_{i=0}^{N}$ be the solution of the equidistribution problem. From the condition (C3), we see that $\sum_{i=1}^{N} \tilde{M}\left(y_{i-1}, y_{i}\right)\left(y_{i}-y_{i-1}\right)$ is simply a Riemann sum for the integral $\int_{0}^{1} M(y) d y$, so it converges to the integral as $N \rightarrow \infty$. Thus, the Riemann sum is bounded for all $N$ by a constant, say $\omega_{1}>0$. However, from the relation

$$
\tilde{M}\left(y_{0}, y_{1}\right)\left(y_{1}-y_{0}\right)=\tilde{M}\left(y_{1}, y_{2}\right)\left(y_{2}-y_{1}\right)=\cdots=\tilde{M}\left(y_{N-1}, y_{N}\right)\left(y_{N}-y_{N-1}\right)
$$

we see that for all $i$, we have

$$
\tilde{M}\left(y_{i-1}, y_{i}\right)\left(y_{i}-y_{i-1}\right)=\frac{1}{N} \sum_{i=1}^{N} \tilde{M}\left(y_{i-1}, y_{i}\right)\left(y_{i}-y_{i-1}\right) \leq \frac{\omega_{1}}{N}
$$

This implies

$$
0<y_{i}-y_{i-1}<\frac{\omega_{1}}{\check{m} N}, \quad i=1, \ldots, N
$$

Next, we have by definition

$$
\begin{aligned}
d_{i} & =-\frac{\partial \phi_{i+1 / 2}}{\partial y_{i}}=\frac{\partial}{\partial y_{i}}\left(\tilde{M}\left(y_{i}, y_{i+1}\right)\left(y_{i}-y_{i+1}\right)\right)=\tilde{M}\left(y_{i}, y_{i+1}\right)+\left.\frac{\partial \tilde{M}}{\partial u}\right|_{\left(y_{i}, y_{i+1}\right)}\left(y_{i}-y_{i+1}\right) \\
c_{i} & =\frac{\partial \phi_{i-1 / 2}}{\partial y_{i}}=\frac{\partial}{\partial y_{i}}\left(\tilde{M}\left(y_{i-1}, y_{i}\right)\left(y_{i}-y_{i-1}\right)\right)=\tilde{M}\left(y_{i-1}, y_{i}\right)+\left.\frac{\partial \tilde{M}}{\partial v}\right|_{\left(y_{i-1}, y_{i}\right)}\left(y_{i}-y_{i-1}\right)
\end{aligned}
$$

which means

$$
c_{i} \geq \check{m}-B \frac{\omega_{1}}{\check{m} N} \geq \frac{\check{m}}{2}>0
$$

for $N$ large enough, and the exact same bound holds for $d_{i}$. Moreover, we have

$$
d_{i}-c_{i}=\underbrace{\tilde{M}\left(y_{i}, y_{i+1}\right)-\tilde{M}\left(y_{i-1}, y_{i}\right)}_{(*)}+\left.\frac{\partial \tilde{M}}{\partial u}\right|_{\left(y_{i}, y_{i+1}\right)}\left(y_{i}-y_{i+1}\right)-\left.\frac{\partial \tilde{M}}{\partial v}\right|_{\left(y_{i-1}, y_{i}\right)}\left(y_{i}-y_{i-1}\right)
$$

The term in $(*)$ can be written as

$$
\begin{aligned}
\tilde{M}\left(y_{i}, y_{i+1}\right)-\tilde{M}\left(y_{i-1}, y_{i}\right) & =\tilde{M}\left(y_{i}, y_{i+1}\right)-\tilde{M}\left(y_{i}, y_{i}\right)+\tilde{M}\left(y_{i}, y_{i}\right)-\tilde{M}\left(y_{i-1}, y_{i}\right) \\
& =\left.\frac{\partial \tilde{M}}{\partial v}\right|_{\left(y_{i}, \zeta^{+}\right)}\left(y_{i+1}-y_{i}\right)+\left.\frac{\partial \tilde{M}}{\partial u}\right|_{\left(\zeta^{-}, y_{i}\right)}\left(y_{i}-y_{i-1}\right)
\end{aligned}
$$

where $y_{i-1}<\zeta^{-}<y_{i}$ and $y_{i}<\zeta^{+}<y_{i+1}$. This implies

$$
\begin{aligned}
\left|d_{i}-c_{i}\right| \leq & \left|\left(\left.\frac{\partial \tilde{M}}{\partial v}\right|_{\left(y_{i}, \zeta^{+}\right)}\left(y_{i+1}-y_{i}\right)-\left.\frac{\partial \tilde{M}}{\partial v}\right|_{\left(y_{i-1}, y_{i}\right)}\left(y_{i}-y_{i-1}\right)\right)\right| \\
& +\left|\left(\left.\frac{\partial \tilde{M}}{\partial u}\right|_{\left(\zeta^{-}, y_{i}\right)}\left(y_{i}-y_{i-1}\right)-\left.\frac{\partial \tilde{M}}{\partial u}\right|_{\left(y_{i}, y_{i+1}\right)}\left(y_{i+1}-y_{i}\right)\right)\right| \\
\leq & \frac{\omega_{1}}{N}\left|\left(\left.\frac{1}{\tilde{M}\left(y_{i}, y_{i+1}\right)} \frac{\partial \tilde{M}}{\partial v}\right|_{\left(y_{i}, \zeta^{+}\right)}-\left.\frac{1}{\tilde{M}\left(y_{i-1}, y_{i}\right)} \frac{\partial \tilde{M}}{\partial v}\right|_{\left(y_{i-1}, y_{i}\right)}\right)\right| \\
& +\frac{\omega_{1}}{N}\left|\left(\left.\frac{1}{\tilde{M}\left(y_{i-1}, y_{i}\right)} \frac{\partial \tilde{M}}{\partial u}\right|_{\left(\zeta^{-}, y_{i}\right)}-\left.\frac{1}{\tilde{M}\left(y_{i}, y_{i+1}\right)} \frac{\partial \tilde{M}}{\partial u}\right|_{\left(y_{i}, y_{i+1}\right)}\right)\right|
\end{aligned}
$$

Thus, within each absolute value term, we have a difference of bounded Lipschitz functions evaluated at points that are at most $\omega_{1} /(\check{m} N)$ apart, which means there exists a constant $\omega_{2}$ such that

$$
\left|d_{i}-c_{i}\right| \leq \omega_{2} / N^{2}
$$

This implies

$$
\frac{d_{i}}{c_{i}}=1+\frac{d_{i}-c_{i}}{c_{i}} \geq 1-\frac{\omega_{2}}{N^{2}} \cdot \frac{2}{\check{m}}
$$

for $N$ large enough. Hence, putting $\omega=2 \omega_{2} / \check{m}$ and $\mu=1-\omega / N^{2}$ leads to $d_{i} / c_{i} \geq \mu$. Similarly, we have

$$
\frac{c_{i}}{d_{i}}=1+\frac{c_{i}-d_{i}}{d_{i}} \geq 1-\frac{\omega_{2}}{N^{2}} \cdot \frac{2}{\check{m}}=\mu
$$

so we have $c_{i} / d_{i} \geq \mu$, or equivalently, $d_{i} / c_{i} \leq \mu^{-1}$, which completes the proof.
We are finally ready to prove the second main result of this section.
Proof. (Theorem 5.2) We need to show that as $N \rightarrow \infty$, we have

$$
P_{j} \rightarrow \frac{\xi_{b_{j}}-\xi_{b_{j-1}}}{\xi_{b_{j}}-\xi_{a_{j}}}, \quad Q_{j} \rightarrow \frac{\xi_{a_{j+1}}-\xi_{a_{j}}}{\xi_{b_{j}}-\xi_{a_{j}}}, \quad R_{j} \rightarrow \frac{\xi_{b_{j-1}}-\xi_{a_{j}}}{\xi_{b_{j}}-\xi_{a_{j}}}, \quad S_{j} \rightarrow \frac{\xi_{b_{j}}-\xi_{a_{j+1}}}{\xi_{b_{j}}-\xi_{a_{j}}} .
$$

We only prove the result for $S_{j}$, the other coefficients being similar. From Theorem 5.1, we have

$$
\begin{equation*}
\mu^{a_{j+1}-a_{j}} \cdot \frac{\mu^{-\left(b_{j}-a_{j+1}\right)}-1}{\mu^{-\left(b_{j}-a_{j}\right)}-1} \leq S_{j} \leq \mu^{-\left(a_{j+1}-a_{j}\right)} \cdot \frac{1-\mu^{b_{j}-a_{j+1}}}{1-\mu^{b_{j}-a_{j}}} \tag{5.8}
\end{equation*}
$$

so if we can show that both the upper and lower bounds tend to $\left(\xi_{b_{j}}-\xi_{a_{j+1}}\right) /\left(\xi_{b_{j}}-\xi_{a_{j}}\right)$ as $N \rightarrow \infty$, then we are done. Indeed, from Lemma 5.5, we can choose $\mu=1-\omega / N^{2}$, so that

$$
\begin{equation*}
\mu^{N}=\left(1-\frac{\omega}{N^{2}}\right)^{N} \rightarrow 1 \quad \text { as } N \rightarrow \infty \tag{5.9}
\end{equation*}
$$

Thus, by letting $\lambda=\mu^{N}$ and noting that $a_{j}=N \xi_{a_{j}}, b_{j}=N \xi_{b_{j}}$, we can rewrite (5.8) as

$$
\lambda^{\xi_{a_{j+1}}-\xi_{a_{j}}} \cdot \frac{\lambda^{-\left(\xi_{b_{j}}-\xi_{a_{j+1}}\right)}-1}{\lambda^{-\left(\xi_{b_{j}}-\xi_{a_{j}}\right)}-1} \leq S_{j} \leq \lambda^{-\left(\xi_{a_{j+1}}-\xi_{a_{j}}\right)} \cdot \frac{1-\lambda^{\xi_{b_{j}}-\xi_{a_{j+1}}}}{1-\lambda^{\xi_{b_{j}}-\xi_{a_{j}}}}
$$

where the exponents are now independent of $N$. Hence, by taking the limit as $\lambda \rightarrow 1$, we obtain

$$
\lambda^{\left(\xi_{a_{j+1}}-\xi_{a_{j}}\right)} \rightarrow 1, \quad \frac{\lambda^{-\left(\xi_{b_{j}}-\xi_{a_{j+1}}\right)}-1}{\lambda^{-\left(\xi_{b_{j}}-\xi_{a_{j}}\right)}-1} \rightarrow \frac{\xi_{b_{j}}-\xi_{a_{j+1}}}{\xi_{b_{j}}-\xi_{a_{j}}}, \quad \frac{1-\lambda^{\xi_{b_{j}}-\xi_{a_{j+1}}}}{1-\lambda^{\xi_{b_{j}}-\xi_{a_{j}}}} \rightarrow \frac{\xi_{b_{j}}-\xi_{a_{j+1}}}{\xi_{b_{j}}-\xi_{a_{j}}}
$$

Thus, we have

$$
S_{j} \rightarrow \frac{\xi_{b_{j}}-\xi_{a_{j+1}}}{\xi_{b_{j}}-\xi_{a_{j}}}
$$

as required.
REmARK. We now see the importance of the $N^{2}$ in the estimate of Lemma 5.5: if we only had the estimate $\mu=1-\omega / N$, then (5.9) would give $\mu^{N} \rightarrow e^{-\omega}$, and the upper and lower bounds of $S_{j}$ would not converge to the same limit.

## 6. Numerical Results.

6.1. A first example. We begin by considering the parallel Schwarz algorithm (4.4) applied to the steady mesh generation problem with

$$
M(x)=\sqrt{1+u_{x}^{2}} \quad \text { and } \quad u(x)=\frac{1-e^{R x}}{1-e^{R}}
$$

For $R=3$ the solution $u$ exhibits a boundary layer at $x=1$ and the mesh should have increased density in the location of the layer. For $x$ outside of $[0,1]$ we have extended $M$ by a constant.

The single domain solution, computed by solving the BVP with a Newton iteration is shown in Figure 6.1. We see a concentration of mesh points near $x=1$. The supersolution and subsolution are computed using $N$ points where $N$ is calculated by equation (3.1) using the bounds $\hat{m}=3.4, \check{m}=1$ and $B=9.1$. From equation (4.6) this gives a requirement of $N=198$ points. The shooting approach outlined in the proof of Lemma 3.1 is used along with a simple bisection approach to find $s^{*}$. We plot the supersolution, subsolution, the uniform initial guess and the single domain solution in Figure 6.2.


Fig. 6.1. Single domain solution $x(\xi)$ obtained with a Newton iteration for the boundary layer problem


Fig. 6.2. Supersolution and subsolution for the boundary layer example.

In Figures 6.3 and 6.4 we illustrate the iterates obtained using the parallel DD iteration on two subdomains starting from the supersolution and subsolution. The plot shows iterates converging to the single domain solution.


Fig. 6.3. $D D$ iterations from the supersolution.


Fig. 6.4. $D D$ iterations from the subsolution.

In Figures 6.5 and 6.6 we illustrate the convergence histories on subdomains one and two with varying amounts of overlap $K$. We see the convergence improves as the amount of overlap increases. This is consistent with the asymptotic convergence rate obtained in Section 5.


Fig. 6.5. The $\|\cdot\|_{\infty}$ difference between the $D D$ iterates on subdomain one.


Fig. 6.6. The $\|\cdot\|_{\infty}$ difference between the $D D$ iterates on subdomain two.

We see very close agreement between the theoretical and numerically estimated convergence rates in Table 6.1. The numerically estimated convergence rates are obtained by averaging the ratio or errors over five iterations as we near convergence.

|  | Overlap |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  | 10 | 20 | 30 | 40 | 50 |
| $\rho$ | 0.8302 | 0.6791 | 0.5526 | 0.4482 | 0.3625 |
| $\hat{\rho}$ | 0.8329 | 0.6791 | 0.5526 | 0.4482 | 0.3625 |
| TABLE 6.1 |  |  |  |  |  |

Expected ( $\hat{\rho}$ ) and obtained ( $\rho$ ) convergence rates as a function of overlap for the boundary layer problem.
6.2. A more challenging problem. We now consider a more difficult problem by choosing

$$
M(x)=1+40 e^{-(x-0.25)^{2} / 0.01}+10 e^{-(x-0.75)^{2} / 0.05}
$$

This choice of $M$ will concentrate mesh nodes where $M$ is large, i.e. near $x=0.25$ and $x=0.75$, as shown in Figure 6.7. A quick calculation gives the parameters $\check{m}=1.7$, $\hat{m}=6$ and $B=45$. Our theory indicates that using $N=1496$ points is sufficient to guarantee the existence of a subsolution and supersolution for this problem.

The single domain Newton solver requires 6 iterations and the resulting solution is shown in Figure 6.7.


Fig. 6.7. Single domain solution computed using a Newton iteration.
In Figure 6.8 we illustrate the supersolution and subsolution for this problem along with the single domain solution (in solid black). The uniform mesh is shown in a dashed black line. The DD iterates obtained starting from the supersolution are shown in Figure 6.9.


FIG. 6.8. Supersolution and subsolution for problem two along with the single domain solution.


Fig. 6.9. DD iterates from supersolution for problem two.

The theory presented in Section 4 suggests that the sequence of meshes, $x^{j, k}$ obtained from an initial uniform mesh, should lie between $\check{x}^{j, k}$ and $\hat{x}^{j, k}$, the sequence of
meshes obtained by starting the DD iteration with the subsolution $\check{x}$ and supersolution $\hat{x}$. This is confirmed in Table 6.2 where we report the differences between $x^{1, k}$ and $\hat{x}^{1, k}$ and $\check{x}^{1, k}$ at the interface $\xi_{a_{2}}$. Indeed these differences remain positive and are monotonically decreasing to zero as $k$ increases.

| $k$ | 1 | 2 | 3 | 4 | 5 | 6 | 7 | 8 | 9 | 10 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\left(\hat{x}^{1, k}-x^{1, k}\right)\left(\xi_{a_{2}}\right)$ | 6.18 | 5.31 | 4.58 | 3.86 | 3.30 | 2.74 | 2.32 | 1.89 | 1.56 | 1.23 |
| $\left(x^{1, k}-\check{x}^{1, k}\right)\left(\xi_{a_{2}}\right)$ | 10.26 | 8.76 | 6.68 | 5.59 | 3.99 | 3.16 | 1.93 | 1.39 | 0.36 | 0.18 |

The sequence of meshes starting with a uniform initial guess are trapped between the converging sequence of meshes starting from the subsolution and supersolution.

In Table 6.3 we see if we can recover the asymptotic rate of convergence as $N$ increases. As we increase $N$ we adjust the number of mesh points in the overlap region so as keep the overlap width approximately constant. For each value of $N$ we estimate the numerical rate of convergence by averaging the ratio of errors over five iterations as we near convergence. We also record the theoretical convergence rate. The theoretical convergence rate $\hat{\rho}$ changes as the overlap width varies slightly with increasing $N$.

|  | $N$ |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | 25 | 50 | 75 | 100 | 125 | 150 |
| $\rho$ | 0.5233 | 0.4789 | 0.4663 | 0.4605 | 0.4569 | 0.4547 |
| $\hat{\rho}$ | 0.5102 | 0.4756 | 0.4648 | 0.4596 | 0.4564 | 0.4545 |
| TABLE 6.3 |  |  |  |  |  |  |

Expected ( $\hat{\rho}$ ) and obtained ( $\rho$ ) convergence rates as a function of overlap for the boundary layer problem.
6.3. A note on robustness. As $M$ becomes more difficult the solution of the nonlinear BVP for the mesh transformation using a Newton iteration becomes quite difficult. Indeed a very good initial guess is often required. Here we show that the nonlinear Schwarz DD method, with a sufficient number of subdomains, gives additional robustness, and is able to recover the single domain solution from a uniform initial guess in a situation where the single domain Newton iteration fails.

Here we consider a more challenging $M$ given by

$$
\begin{equation*}
M(x)=1+30 e^{-(x-0.25)^{2} / 0.01}+40 e^{-(x-0.75)^{2} / 0.05} \tag{6.1}
\end{equation*}
$$

In Table 6.4 we show the difficulty that the single domain Newton solve has for this nonlinear BVP (starting from a uniform initial guess) for various values of $N$ near 100. During the DD solution we alter the global number of mesh points slightly to ensure that each interior subdomain has the same number of degrees of freedom. We see that for $N$ near 100 the Newton iteration has difficulty converging on a single domain. Indeed for $N=101$ and $N=105$ no convergence is obtained even after 100 Newton iterations. For the other values of $N$ the convergence only seems to occur by happenchance!

| $N$ | 100 | 101 | 103 | 104 | 105 |
| :---: | :---: | :---: | :---: | :---: | :---: |
| \# iters | 17 | NC | 67 | 26 | NC |

TABLE 6.4
The number of single domain Newton iterations versus the number of mesh points $N$ for the mesh generation $B V P$ with $M$ given in (6.1).

We now consider dividing the unit interval into subdomains and using a parallel Schwarz iteration to obtain the global solution. Within each subdomain, we use Newton's method to solve the nonlinear subdomain problems. In Table 6.5 we give the number of Newton iterations required for each subdomain solve during the first DD iteration as a function of the number of subdomains used. We can see that as the number of subdomains increases the nonlinear problems solved on the subdomains become easier. For example, using $N=101$ points the global Newton iteration fails. Using 4 subdomains convergence is achieved, albeit with great difficulty, requiring 43 Newton iterations on subdomain 1. With 8 subdomains, and again 101 global mesh points, only 4-7 Newton iterations are required on each subdomain. The DD iteration itself converges for both 4 and 8 subdomains and hence we are able to recover the global solution which is not possible with the global Newton iteration.

| $S$ | $N$ | $\Omega_{1}$ | $\Omega_{2}$ | $\Omega_{3}$ | $\Omega_{4}$ | $\Omega_{5}$ | $\Omega_{6}$ | $\Omega_{7}$ | $\Omega_{8}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 100 | 17 |  |  |  |  |  |  |  |
| 2 | 101 | 75 | 5 |  |  |  |  |  |  |
| 3 | 101 | 31 | 6 | 5 |  |  |  |  |  |
| 4 | 101 | 43 | 6 | 5 | 5 |  |  |  |  |
| 5 | 100 | 43 | 6 | 5 | 4 | 5 |  |  |  |
| 6 | 101 | 8 | 5 | 5 | 5 | 4 | 5 |  |  |
| 7 | 103 | 7 | 5 | 6 | 5 | 4 | 4 | 5 |  |
| 8 | 101 | 7 | 5 | 5 | 4 | 5 | 4 | 4 | 5 |
|  |  |  |  | TABLE 6.5 |  |  |  |  |  |

The number of Newton iterations required on each subdomain as a function of the number of subdomains, $S$, using approximately $N=100$ mesh points.

In Tables 6.6 and 6.7 we repeat the experiment above but now with approximately 150 global mesh points. A similar pattern emerges. The single domain Newton iteration fails or requires many iterations. With a sufficient number of subdomains the global solution is recovered.

| $N$ | 150 | 151 | 152 | 153 | 155 | 157 |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: |
| \# iters | 65 | 28 | NC | NC | NC | 76 |
|  | Table 6.6 |  |  |  |  |  |
| \# iters versus $N$ |  |  |  |  |  |  |


| $S$ | $N$ | $\Omega_{1}$ | $\Omega_{2}$ | $\Omega_{3}$ | $\Omega_{4}$ | $\Omega_{5}$ | $\Omega_{6}$ | $\Omega_{7}$ | $\Omega_{8}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 150 | 65 |  |  |  |  |  |  |  |
| 2 | 151 | NC | 5 |  |  |  |  |  |  |
| 3 | 152 | 21 | 6 | 5 |  |  |  |  |  |
| 4 | 153 | 32 | 6 | 5 | 5 |  |  |  |  |
| 5 | 150 | 9 | 6 | 5 | 4 | 5 |  |  |  |
| 6 | 155 | 8 | 4 | 5 | 5 | 4 | 5 |  |  |
| 7 | 152 | 7 | 5 | 5 | 5 | 4 | 4 | 5 |  |
| 8 | 157 | 7 | 5 | 5 | 4 | 5 | 4 | 4 | 4 |

The number of Newton iterations required on each subdomain as a function of the number of subdomains, $S$, using approximately $N=150$ mesh points.
6.4. DD within the $M P$ framework. We conclude our presentation of numerical results with a short illustration of how the domain decomposition method performs within the $M P$ alternating procedure described in the introduction.

Consider the solution of the two-point boundary value problem

$$
-\epsilon u^{\prime \prime}+u^{\prime}=1, u(0)=u(1)=0
$$

with $\epsilon=0.05$. The physical solution

$$
u(x)=x-\frac{1-e^{x / \epsilon}}{1-e^{1 / \epsilon}}
$$

has a sharp boundary layer near $x=1$.
The MP procedure proceeds as follows. Pick an (often uniform) initial mesh $X^{(0)}$. For $k=0,1, \ldots$

- Solve the physical PDE for $U^{(k)}$ on the mesh $X^{(k)}$.
- Solve for $X^{(k+1)}$ using the mesh DE, where the monitor function $M(x)$ is calculated based on $X^{(k)}$ and $U^{(k)}$, see below.
In the mesh solve step, we solve

$$
\begin{equation*}
M\left(\left(x_{i}+x_{i-1}\right) / 2\right)\left(x_{i}-x_{i-1}\right)-M\left(\left(x_{i}+x_{i+1}\right) / 2\right)\left(x_{i+1}-x_{i}\right)=0 \tag{6.2}
\end{equation*}
$$

for the new mesh $\left\{x_{i}\right\}$, where the continuous function $M$ is constructed by interpolating the ordered pairs $\left(\hat{x}_{i-1 / 2}, M_{i-1 / 2}\right)$ where the $M_{i-1 / 2}$ are computed using the previous mesh $\{\hat{x}\}$ and physical solution on that mesh. We choose

$$
M_{i-1 / 2}=\sqrt{1+\left(\frac{u_{i}-u_{i-1}}{\hat{x}_{i}-\hat{x}_{i-1}}\right)^{2}}
$$

as an approximation to $M$ at the half nodes $\hat{x}_{i-1 / 2}=\left(\hat{x}_{i}+\hat{x}_{i-1}\right) / 2$.
The convergence of the $M P$ iteration procedure, starting with an uniform initial mesh with $N=51$ mesh points, is shown in Figure 6.10.


Fig. 6.10. $M P$ iteration convergence for a boundary layer problem.

The initial uniform mesh and the solution to the boundary value problem is shown in Figure 6.11.


Fig. 6.11. Solution of the BVP and an initial uniform mesh.

The solution of the BVP and the first mesh $X^{(1)}$ obtained by solving (6.2) during the $M P$ process described above is shown in Figure 6.12.


Fig. 6.12. Solution of the $B V P$ and the first mesh obtained during our MP procedure.

During the $M P$ procedure one could linearize the mesh solves by simply evaluating the mesh density function at the previously computed mesh and solution. In other words, instead of constructing $\tilde{M}$ by interpolation and evaluating at $\left(x_{i}+x_{i-1}\right) / 2$ and $\left(x_{i}+x_{i+1}\right) / 2$, one could instead use the values of $M_{i-1 / 2}$ and $M_{i+1 / 2}$ directly. This results in a different mesh $X^{(1)}$, shown in Figure 6.13. Comparing Figures 6.12 and 6.13 , we see that the initial mesh quality using the linearized variant is clearly not as good as the $M P$ procedure we describe.


Fig. 6.13. Solution of the BVP and the first mesh obtained during the linearized MP procedure.

During each $M P$ iteration, the previous mesh $\left\{\hat{x}_{i}\right\}$ is used as initial guess for the DD iteration to obtain the new mesh. Using two subdomains, the number of DD iterations required for the difference in DD approximations to agree within $10^{-5}$ for the first $7 M P$ iterations are given in Table 6.8.

| $M P$ iteration | 1 | 2 | 3 | 4 | 5 | 6 | 7 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| \# DD iterations | 63 | 49 | 54 | 33 | 19 | 22 | 11 |

The number of $D D$ iterations for each MP iteration.

We can see that starting the DD iteration for (6.2) using the mesh found during the previous $M P$ iteration generally reduces the number of $D D$ iterations required to converge to the specified tolerance.
7. Conclusions. In this paper, we have shown that when the integral form of the moving mesh PDE is discretized using a quadrature rule, the resulting nonlinear algebraic system may not have a solution if the number of grid points is too small. Using $M$-function theory, we have proved the existence and uniqueness of the solution of the discrete problem when the number of grid points is large enough. Moreover, we have shown that the parallel and alternating Schwarz methods converge to the unique solution, so they can be used to solve the discrete solution in parallel. In fact, as the mesh is refined, these methods behave asymptotically the same way as their continuous counterparts, with the same iteration matrix and the same convergence rates. This is also confirmed by our numerical results. Our numerical experiments show that the domain decomposition algorithms are often more robust than the methods that try to find the single-domain solution all at once, such as Newton's method. And finally, we show that the DD approach can effectively be used within the $M P$ iteration to solve the coupled problem.

We believe that this $M$-function theory can be applied to analyze the behavior of domain decomposition methods when applied to other nonlinear PDEs; this is a subject of ongoing investigation. Ongoing work also focuses on understanding the behavior when optimized conditions are used, see [22], as well as for the 2D mesh equation in [33].

Appendix A. M-function charaterization. To prove Lemma 2.8, we will first need the following geometrical fact.

Proposition A.1. Let $\mathbb{D}$ be an open convex subset of $\mathbb{R}^{n}$. For two given points $x, y \in \mathbb{D}$, let $r>0$ be a constant such that the open balls $B_{r}(x)$ and $B_{r}(y)$, of radius $r$ and centered at $x$ and $y$ respectively, are contained in $\mathbb{D}$. Then for any intermediate point $p=(1-t) x+t y$ with $t \in[0,1]$, we have $B_{r}(p) \subset \mathbb{D}$.

Proof. Let $v \in \mathbb{R}^{n}$ be a vector with $\|v\|<r$, so that $x+v \in \mathbb{D}$ and $y+v \in \mathbb{D}$. Then by the convexity of $\mathbb{D}$, we have

$$
(1-s)(x+v)+s(y+v)=(1-s) x+s y+v \in \mathbb{D} \quad \forall s \in[0,1]
$$

Choosing $s=t$ shows that $p+v \in \mathbb{D}$. Since $v$ is arbitrary, we conclude that $B_{r}(p) \subset \mathbb{D}$. $\square$

Now we proceed to prove Lemma 2.8 which we restate here.
Lemma 2.8 Let $\mathbb{D}$ be a convex and open subset of $\mathbb{R}^{n}$. Assume $F: \mathbb{D} \subset R^{n} \rightarrow R^{n}$ is off-diagonally antitone, and that for any $\mathbf{x} \in \mathbb{D}$, the functions $q_{i}: S_{i} \subset \mathbb{R} \rightarrow \mathbb{R}^{n}$ defined as

$$
q_{i}(\tau)=\sum_{j=1}^{n} f_{j}\left(\mathbf{x}+\tau \mathbf{e}^{i}\right), i=1, \ldots, n, \quad \text { with } \quad S_{i}=\left\{\tau: \mathbf{x}+\tau \mathbf{e}^{i} \in \mathbb{D}\right\}
$$

are isotone. Finally, assume for every $i$ that either $q_{i}$ is strictly isotone, or there exists a node $j=j(i)$ such that a path $j \rightsquigarrow i$ exists and $q_{j}$ is strictly isotone. Then $F$ is an $M$-function.

Proof. To prove Lemma 2.8 we need to show that $F$ is inverse isotone, so let $\mathbf{x} \neq \mathbf{y} \in \mathbb{D}$ be such that $F(\mathbf{x}) \leq F(\mathbf{y})$. We must now show that $\mathbf{x} \leq \mathbf{y}$. Let $r>0$ be such that $B_{r}(x)$ and $B_{r}(y)$ are both contained in $\mathbb{D}$. Moreover, let $\mathbf{x}_{s}=(1-s) \mathbf{y}+s \mathbf{x}$, so that $s \in[0,1] \Longrightarrow \mathbf{x}_{s} \in \mathbb{D}$. Now define $t \in(0,1]$ so that

$$
\frac{1}{t}=\max \left\{1,\left\lceil\frac{N \max _{i}\left|x_{i}-y_{i}\right|}{r}\right\rceil\right\}=: M \in \mathbb{N}
$$

Then the points $\mathbf{y}=\mathbf{x}_{0}, \mathbf{x}_{t}, \mathbf{x}_{2 t}, \ldots, \mathbf{x}_{M t}=\mathbf{x}$ all satisfy $B_{r}\left(\mathbf{x}_{k t}\right) \subset \mathbb{D}$ for $k=0, \ldots, M$ by Proposition A.1. Now define the sets

$$
N^{-}=\left\{i \in \mathbb{N}: y_{i}<x_{i}\right\}, \quad N^{+}=\left\{i \in \mathbb{N}: y_{i} \geq x_{i}\right\}
$$

Let us assume $N^{-} \neq \emptyset$, and let $i_{1}, \ldots, i_{m}$ be nodes of $N^{-}$and set

$$
\mathbf{z}^{j}=\left(x_{i_{j}}-y_{i_{j}}\right) \mathbf{e}^{i_{j}}, \quad j=1, \ldots, m, \quad \mathbf{z}=\sum_{j=1}^{m} \mathbf{z}^{j}
$$

where $\mathbf{e}^{i_{j}}$ is the $i_{j}$-th coordinate vector. Noting that $\left\|t \mathbf{z}^{j}\right\| \leq r / N$, we see that the sequence

$$
\mathbf{y}+t \mathbf{z}^{1}, \mathbf{y}+t \mathbf{z}^{1}+t \mathbf{z}^{2}, \ldots, \mathbf{y}+t \sum_{j=1}^{m} \mathbf{z}^{j}=\mathbf{y}+t \mathbf{z}
$$

all belong to $\mathbb{D}$, since $\mathbb{D}$ is convex. Then by the isotonicity of the $q_{i}$, we have

$$
\begin{equation*}
\sum_{j=1}^{n} f_{j}(\mathbf{y}) \leq \sum_{j=1}^{n} f_{j}\left(\mathbf{y}+t \mathbf{z}^{1}\right) \leq \sum_{j=1}^{n} f_{j}\left(\mathbf{y}+t \mathbf{z}^{1}+t \mathbf{z}^{2}\right) \leq \cdots \leq \sum_{j=1}^{n} f_{j}(\mathbf{y}+t \mathbf{z}) \tag{A.1}
\end{equation*}
$$

Splitting the last sum into sums over $N^{-}$and $N^{+}$, we use the antitonicity of the $f_{j}$ to obtain

$$
\begin{equation*}
\sum_{j=1}^{n} f_{j}(\mathbf{y}+t \mathbf{z})=\sum_{j \in N^{-}} f_{j}(\mathbf{y}+t \mathbf{z})+\sum_{j \in N^{+}} f_{j}(\mathbf{y}) \leq \sum_{j \in N^{-}} f_{j}\left(\mathbf{x}_{t}\right)+\sum_{j \in N^{+}} f_{j}(\mathbf{y}) \tag{A.2}
\end{equation*}
$$

Thus, we have

$$
\begin{equation*}
\sum_{j \in N^{-}} f_{j}(\mathbf{y}) \leq \sum_{j \in N^{-}} f_{j}\left(\mathbf{x}_{t}\right) \tag{A.3}
\end{equation*}
$$

We can now continue this argument to obtain

$$
\begin{equation*}
\sum_{j \in N^{-}} f_{j}(\mathbf{y}) \leq \sum_{j \in N^{-}} f_{j}\left(\mathbf{x}_{t}\right) \leq \sum_{j \in N^{-}} f_{j}\left(\mathbf{x}_{2 t}\right) \leq \cdots \leq \sum_{j \in N^{-}} f_{j}\left(\mathbf{x}_{M t}\right)=\sum_{j \in N^{-}} f_{j}(\mathbf{x}) \tag{A.4}
\end{equation*}
$$

But since we had assumed $F(\mathbf{x}) \leq F(\mathbf{y})$, we must have $f_{j}(\mathbf{y})=f_{j}(\mathbf{x})$ for all $j \in N^{-}$, so that all the inequalities in (A.4) (and hence (A.1), (A.2) and (A.3)) are in fact
equalities. This, in turn, shows that $q_{i}$ cannot be strictly isotone for any $i \in N^{-}$. Thus, there must exist a strict $\operatorname{link}\left(i, i^{\prime}\right)$ such that $i \in N^{+}, i^{\prime} \in N^{-}$, i.e., we have

$$
f_{i}(\mathbf{y})>f_{i}\left(\mathbf{y}+t\left(x_{i^{\prime}}-y_{i^{\prime}}\right) \mathbf{e}^{i^{\prime}}\right)
$$

as well as

$$
f_{j}(\mathbf{y}) \geq f_{j}\left(\mathbf{y}+t\left(x_{i^{\prime}}-y_{i^{\prime}}\right) \mathbf{e}^{i^{\prime}}\right) \quad \text { whenever } j \in N^{+}
$$

Thus,

$$
\begin{aligned}
\sum_{j=1}^{n} f_{j}(\mathbf{y}) & =\sum_{j \in N^{-}} f_{j}(\mathbf{y}+t \mathbf{z})+\sum_{j \in N^{+}} f_{j}(\mathbf{y}) \\
& >\sum_{j \in N^{-}} f_{j}(\mathbf{y}+t \mathbf{z})+\sum_{j \in N^{+}} f_{j}\left(\mathbf{y}+t\left(x_{i^{\prime}}-y_{i^{\prime}}\right) \mathbf{e}^{i^{\prime}}\right) \\
& \geq \sum_{j \in N^{-}} f_{j}(\mathbf{y}+t \mathbf{z})+\sum_{j \in N^{+}} f_{j}(\mathbf{y}+t \mathbf{z}) \\
& =\sum_{j=1}^{n} f_{j}(\mathbf{y}+t \mathbf{z})=\sum_{j=1}^{n} f_{j}(\mathbf{y})
\end{aligned}
$$

which is a contradiction. Thus $N^{-}$must be empty, so that $F$ is inverse isotone, and hence an $M$-function.

Appendix B. Discrete maximum principles. In this section, we prove two discrete maximum principles, which are used to deduce the number of grid points necessary to guarantee the existence of subdomain iterates.

Lemma B. 1 (Steady mesh generation). Let $\left\{u_{i}\right\}_{i=0}^{N}$ be such that $u_{0}=u_{L}, u_{N}=$ $u_{R}$ and

$$
\tilde{M}\left(u_{i}, u_{i+1}\right)\left(u_{i+1}-u_{i}\right)-\tilde{M}\left(u_{i-1}, u_{i}\right)\left(u_{i}-u_{i-1}\right)
$$

for $1 \leq i \leq N-1$. Then

$$
\min \left(u_{L}, u_{R}\right) \leq u_{i} \leq \max \left(u_{L}, u_{R}\right)
$$

Proof. We have

$$
\tilde{M}\left(u_{0}, u_{1}\right)\left(u_{1}-u_{0}\right)=\tilde{M}\left(u_{1}, u_{2}\right)\left(u_{2}-u_{1}\right)=\cdots=\tilde{M}\left(u_{N-1}, u_{N}\right)\left(u_{N}-u_{N-1}\right)
$$

Since $M(u, v) \geq \check{m}>0$ for all $u$ and $v$, the differences $u_{i}-u_{i-1}$ all have the same sign, so there can be no internal minima or maxima. Thus the maximum and minimum values are $u_{0}=u_{L}$ and $u_{N}=u_{R}$ respectively.

Lemma B. 2 (Time-dependent mesh generation). Let $\left\{u_{i}\right\}_{i=0}^{N}$ be such that $u_{0}=$ $u_{L}, u_{N}=u_{R}$ and

$$
\begin{align*}
& u_{i}-\theta \gamma\left(\tilde{M}\left(u_{i}, u_{i+1}\right)\left(u_{i+1}-u_{i}\right)-\tilde{M}\left(u_{i-1}, u_{i}\right)\left(u_{i}-u_{i-1}\right)\right) \\
& \quad=x_{i}^{[\nu]}+(1-\theta) \gamma\left(\tilde{M}\left(x_{i}^{[\nu]}, x_{i+1}^{[\nu]}\right)\left(x_{i+1}^{[\nu]}-x_{i}^{[\nu]}\right)-\tilde{M}\left(x_{i-1}^{[\nu]}, x_{i}^{[\nu]}\right)\left(x_{i}^{[\nu]}-x_{i-1}^{[\nu]}\right)\right) \tag{B.1}
\end{align*}
$$

for $1 \leq i \leq N-1$. If $2(1-\theta) \gamma \hat{m} \leq 1$, then

$$
\min \left(u_{L}, u_{R}, \min _{i} x_{i}\right) \leq u_{i} \leq \max \left(u_{L}, u_{R}, \max _{i} x_{i}\right)
$$

Proof. We show only the maximum, the minimum being similar. Let

$$
\begin{array}{ll}
m_{1}=\tilde{M}\left(u_{i-1}, u_{i}\right), & m_{2}=\tilde{M}\left(u_{i}, u_{i+1}\right) \\
m_{3}=\tilde{M}\left(x_{i-1}^{[\nu]}, x_{i}^{[\nu]}\right), & m_{4}=\tilde{M}\left(x_{i}^{[\nu]}, x_{i+1}^{[\nu]}\right),
\end{array}
$$

so that $\check{m} \leq m_{j} \leq \hat{m}$ for $j=1,2,3,4$, and let $\bar{x}=\max _{i} x_{i}^{[\nu]}$. Then for $i=1, \ldots, N-1$, the right-hand side of (B.1) can be bounded by

$$
\begin{aligned}
x_{i}^{[\nu]} & +(1-\theta) \gamma\left(m_{4}\left(x_{i+1}^{[\nu]}-x_{i}^{[\nu]}\right)-m_{3}\left(x_{i}^{[\nu]}-x_{i-1}^{[\nu]}\right)\right) \\
& =(\underbrace{1-\gamma(1-\theta)\left(m_{3}+m_{4}\right)}_{\geq 0}) x_{i}^{[\nu]}+\gamma(1-\theta) m_{3} x_{i-1}^{[\nu]}+\gamma(1-\theta) m_{4} x_{i+1}^{[\nu]} \\
& \leq\left(1-\gamma(1-\theta)\left(m_{3}+m_{4}\right)\right) \bar{x}+\gamma(1-\theta) m_{3} \bar{x}+\gamma(1-\theta) m_{4} \bar{x}=\bar{x} .
\end{aligned}
$$

On the other hand, suppose $u_{i}$ is the maximum value, i.e., $u_{i} \geq u_{j}$ for all $j$. Then we either have $i=0$ or $i=N$, or the left-hand side of (B.1) satisfies

$$
u_{i}-\theta \gamma(\underbrace{m_{1}\left(u_{i-1}-u_{i}\right)+m_{2}\left(u_{i+1}-u_{i}\right)}_{\leq 0}) \geq u_{i} .
$$

Thus, if the maximum occurs at $1 \leq i \leq N-1$, then we must have $u_{i} \leq \bar{x}$.

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