# GLOBAL APPROXIMATE NEWTON METHODS * 

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#### Abstract

We derive a class of globally convergent and quadratically converging algorithms for a system of nonlinear equations $g(u)=0$, where $g$ is a sufficiently smooth homeomorphism. Particular attention is directed to key parameters which control the iteration. Several examples are given that have successful in solving the coupled nonlinear PDEs which arise in semiconductor device modelling.


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1. Introduction. In this paper we derive an algorithm for solving the nonlinear system

$$
\begin{equation*}
g(u)=0 \tag{1.1}
\end{equation*}
$$

where $g=\left(g_{1}, g_{2}, \ldots, g_{n}\right)^{T}$ is a sufficiently smooth homeomorphism from $\mathbb{R}^{n}$ to $\mathbb{R}^{n}$. Recall that a homeomorphism is a bijection (1-1, onto) with both $g$ and the inverse map, $g^{-1}$, continuous. Physically, a homeomorphism means that the process modelled by $g$ has a unique solution $x$ for any set of input conditions $y$, i.e., $g(x)=y$, and that the solution $x$ varies continuously with input $y$. Sometimes this notion is referred to as a "well-posed" process $g$. Actually, the requirement that $g$ be a homeomorphism is a special case of our assumptions, but we defer a more detailed and general discussion to Section 2.

In its generic form the algorithm we propose is well known. Starting at some initial guess $u_{0}$, we solve and $n \times n$ linear system

$$
\begin{equation*}
M_{k} x_{k}=-g\left(u_{k}\right) \equiv-g_{k} \tag{1.2}
\end{equation*}
$$

and then set

$$
\begin{equation*}
u_{k+1}=u_{k}+t_{k} x_{k} \tag{1.3}
\end{equation*}
$$

We call the method an approximate Newton method because $M_{k}$ will be chosen to be related to the Jacobian $g^{\prime}\left(u_{k}\right) \equiv g_{k}^{\prime}$, in such a manner that $x_{k}$ approximates the Newton step $w_{k}=-\left\{g\left(u_{k}\right)^{\prime}\right\}^{-1} g\left(u_{k}\right)$, and because usually $t_{k} \neq 1$. In many applications, (1.2) can be interpreted as an "inner" iterative method for solving the linear system

$$
\begin{equation*}
g_{k}^{\prime} w_{k}=-g_{k} \tag{1.4}
\end{equation*}
$$

When $g$ is a smooth homeomorphism, we will show how to choose the damping parameters $t_{k}$ and the approximate Newton steps $x_{k}$ such that the $u_{k}$ converge to $u^{*}$ with $g\left(u^{*}\right)=0$ quadratically for any initial $u_{0}$ (see Section 2 for the notions of quadratic and more general higher order convergence). The choice of $x_{k}=w_{k}$ in (1.4), the damped Newton method, in an important special case.

In Section 2 we show that, for any choice of norm $\|\cdot\|$, the choice $t_{k}=\left(1+\mathcal{K}\left\|g_{k}\right\|\right)^{-1}$ for some sequence $0 \leq \mathcal{K}_{k} \leq \mathcal{K}_{0}$ produces the convergence mentioned above. More

[^0]precisely, for $\mathcal{K}_{0}$ sufficiently large, the sequence $\left\|g_{k}\right\|$ decreases monotonically and quadratically to zero. While it is possible in theory to take $\mathcal{K}_{k}=\mathcal{K}_{0}$ for all $k$, such a strategy often leads to the quagmire of slow initial convergence and can prove disastrous in practice (see Sections 3-4). As we shall see, this rule for choosing $t_{k}$ is motivated by the requirement $t_{k} \rightarrow 1$ such that $1-t_{k}=O\left(\left\|g_{k}\right\|\right)$. By specifying a formula for picking $t_{k}$, we attempt to avoid most of the searching common to other damping strategies.

We also show in Section 2 how to choose the $x_{k}$ (or $M_{k}$ ) in (1.2) such that the $x_{k}$ approximates $w_{k}$ of (1.4). In this setting our analysis continues and extends investigations of approximate Newton methods initiated by Dennis and Moré (see $[5,6])$. Motivated by problems in optimization where it may be difficult or undesirable to deal with the Jacobian, $g_{k}^{\prime}$, they choose $M_{k}$, for example, such that

$$
\begin{equation*}
\left\|\left\{M_{k}-g\left(u^{*}\right)\right\}\left(u_{k+1}-u_{k}\right)\right\|\left\|u_{k+1}-u_{k}\right\|^{-1} \rightarrow 0 \tag{1.5}
\end{equation*}
$$

to obtain superlinear convergence. Equation (1.5) may not be immediately useful in contexts where $M_{k}$ represents an iterative process for solving (1.4); contexts, for example, arising from nonlinear PDEs where it is often possible to evaluate $g_{k}^{\prime}$ but it may not be easy to solve (1.4) exactly. Sherman [12] discusses such Newton-iterative methods, showing that to obtain quadratic convergence it suffices to take $m_{k}=O\left(2^{k}\right)$ inner iterations as $k \rightarrow \infty$.

Computationally, it is more convenient to measure the extent to which $x_{k}$ approximates $w_{k}$ by monitoring the quantity $\left\|g_{k}^{\prime} x_{k}+g_{k}\right\|$, that is, checking the residual of (1.4) when $x_{k}$ replaces $w_{k}$. To obtain quadratic convergence, for example, we choose $x_{k}$ such that

$$
\begin{equation*}
\alpha_{k} \equiv \frac{\left\|g_{k}^{\prime} x_{k}+g_{k}\right\|}{\left\|g_{k}\right\|} \leq c\left\|g_{k}\right\| \tag{1.6}
\end{equation*}
$$

$c>0$, with $x_{k} \neq 0$, a suggestion also made by Dembo, Eisenstat and Steihaug [4].
The discussions by the above named researchers all deal with local convergence; that is, they examine convergence in a local region containing a root $u^{*}$ such that the choice $t_{k}=1$ is appropriate. We derive (1.6) within a global framework consistent with our choice of $t_{k}$; that is we require $\alpha_{k}=O\left(\left\|g_{k}\right\|\right)$ to balance the quantity $1-t_{k}=$ $O\left(\left\|g_{k}\right\|\right)$.

We also show that conditions such as (1.6) are equivalent to the original conditions imposed by Dennis and Moré. The whole key to our analysis is the judicious use of the Taylor expansion:

$$
\begin{align*}
g_{k+1}= & g_{k}+g_{k}^{\prime}\left\{u_{k+1}-u_{k}\right\}  \tag{1.7}\\
& +\int_{0}^{1}\left\{g^{\prime}\left(u_{k}+s\left(u_{k+1}-u_{k}\right)\right)-g_{k}^{\prime}\right\}\left\{u_{k+1}-u_{k}\right\} d s
\end{align*}
$$

or, using (1.3) and the notation of Section 2

$$
\begin{align*}
g_{k+1}= & \left(1-t_{k}\right) g_{k}+t_{k}\left\|g_{k}\right\|\left(g_{k}^{\prime} x_{x}+g_{k}\right)\left\|g_{k}\right\|^{-1}  \tag{1.8}\\
& +\int_{0}^{1} G\left(s ; u_{k+1}, u_{k}\right) t_{k} x_{k} d s
\end{align*}
$$

Taking norms will lead to

$$
\begin{equation*}
\left\|g_{k+1}\right\| \leq\left\|g_{k}\right\|\left\{\left(1-t_{k}\right)+t_{k} \alpha_{k}+t_{k}^{2} \beta_{k}\left\|g_{k}\right\|\right\} \tag{1.9}
\end{equation*}
$$

under appropriate conditions on the smoothness of $g$ and the sequence $M_{k}^{-1}$. Since $\beta_{k}$ will be bounded, (1.9) shows it is possible to insure that the $\left\|g_{k}\right\| \rightarrow 0$ monotonically and quadratically by forcing each term in braces to be $O\left(\left\|g_{k}\right\|\right)$. Finally, note that (1.7) can be interpreted in a Banach space, and the extension of our results to such a setting is immediate (see [13], Section 12.1).

Section 2 contains a detailed discussion of our assumptions and global convergence analysis. Section 3 encodes the analysis of Section 2 into a general algorithm. Section 4 presents a further algorithmic discussion and discusses several important examples. We conclude in Section 5 with some numerical results relevant to the solution of semiconductor device partial differential equations.

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2. Parameter Selection and Convergence. Given an arbitrary initial iteration $u_{0}$, we consider here the convergence of the iteration (1.2)-(1.3) where the parameters $t_{k}$ are chosen by the rule

$$
\begin{equation*}
t_{k}=\frac{1}{1+\mathcal{K}_{k}\left\|g_{k}\right\|} \tag{2.1}
\end{equation*}
$$

We make the following assumptions on the mapping $g(u)$ and the sequence $M_{k}$.
Assumption A1: The closed level set

$$
\begin{equation*}
S_{0}=\left\{u \mid\|g(u)\| \leq\left\|g_{0}\right\|\right\} \tag{2.2}
\end{equation*}
$$

is bounded.
Assumption A2: $g$ is differentiable and the Jacobian $g^{\prime}(u)$ is a continuous and nonsingular on $S_{0}$, and the sequence $\left\|M_{k}^{-1}\right\|$ is uniformly bounded, i.e.,

$$
\begin{equation*}
\left\|M_{k}^{-1}\right\| \leq k_{1} \text { on } S_{0} \quad \text { for all } k \geq 0 \tag{2.3}
\end{equation*}
$$

We embed $S_{0}$ in the closed convex ball

$$
\begin{equation*}
S_{1}=\left\{u \mid\|u\| \leq \sup _{v \in S_{0}}\|v\|+k_{1}\left\|g_{0}\right\|\right\} . \tag{2.4}
\end{equation*}
$$

Assumption A3: The Jacobian $g^{\prime}$ is Lipshitz; i.e.,

$$
\begin{equation*}
\left\|g^{\prime}(u)-g^{\prime}(v)\right\| \leq k_{2}\|u-v\| ; \quad u, v \in S_{1} \tag{2.5}
\end{equation*}
$$

Without loss suppose $g_{k} \neq 0$ for all $k$, and let the quantities $\mathcal{A}_{k}, \alpha_{k}, \mathcal{B}_{k}, \beta_{k}$ be defined as follows:

$$
\begin{align*}
\mathcal{A}_{k} & \equiv \frac{g_{k}+g_{k}^{\prime} x_{k}}{\left\|g_{k}\right\|} ;  \tag{2.6}\\
\alpha_{k} & \equiv\left\|\mathcal{A}_{k}\right\|
\end{align*}
$$

and

$$
\begin{align*}
\mathcal{B}_{k} & \equiv \frac{\int_{0}^{1}\left\{g^{\prime}\left(u_{k}+s t_{k} x_{k}\right)-g_{k}^{\prime}\right\} t_{k} x_{k} d s}{\left(t_{k}\left\|g_{k}\right\|\right)^{2}}  \tag{2.7}\\
\beta_{k} & \equiv\left\|\mathcal{B}_{k}\right\| ;
\end{align*}
$$

for $t_{k}, u_{k}$, and $x_{k}$ as in (1.2)-(1.3).
The parameters $\alpha_{k}$ measure the extent to which the $x_{k}$ of (1.2) differ from the Newton correction ( $\alpha_{k} \equiv 0$ ). For $u_{k} \in S_{0}$ note that

$$
\begin{equation*}
\alpha_{k} \leq k_{1}\left\|g_{k}^{\prime}-M_{k}\right\| \tag{2.8}
\end{equation*}
$$

Typically, given the sequence $\left\|g_{k}\right\|$ and $\alpha_{0} \in(0,1)$, we will consider the convergence process when all $\alpha_{k} \leq \alpha_{0}$ and

$$
\begin{equation*}
\alpha_{k} \leq c\left\|g_{k}\right\|^{p} \quad \text { for } p \in(0,1], c>0 \tag{2.9}
\end{equation*}
$$

for example,

$$
\begin{equation*}
\alpha_{k} \leq \alpha_{0}\left(\frac{\left\|g_{k}\right\|}{\left\|g_{0}\right\|}\right)^{p}=\alpha_{k-1}\left(\frac{\left\|g_{k}\right\|}{\left\|g_{k-1}\right\|}\right)^{p} \tag{2.10}
\end{equation*}
$$

For many Newton-like methods the $\alpha_{k}$ can be easily computed, and $\alpha_{0}$ and $p$ can be specified a priori.

The parameters $\beta_{k}$ reflect the size of higher (second) order terms which are ignored in the derivation of a Newton-like method. For $u_{k+1}=u_{k}+t_{k} x_{k} \in \mathcal{S}_{1}$ and $u_{k} \in \mathcal{S}_{0}$, (1.2), (2.3), (2.5), and (2.7) imply

$$
\begin{equation*}
\beta_{k} \leq \frac{k_{2}}{2}\left(\frac{\left\|x_{k}\right\|}{\left\|g_{k}\right\|}\right)^{2} \leq \frac{k_{1}^{2} k_{2}}{2} \tag{2.11}
\end{equation*}
$$

Suppose $u_{k+1} \in S_{1}$ and $u_{k} \in S_{0}$. Taylor's Theorem ([9], Section 3.2) implies

$$
\begin{align*}
g_{k+1}=g_{k}+g_{k}^{\prime}\left\{u_{k+1}-u_{k}\right\} & +\int_{0}^{1}\left\{g^{\prime}\left(u_{k}+s\left(u_{k+1}-u_{k}\right)\right)-g_{k}^{\prime}\right\}\left(u_{k+1}-u_{k}\right) d s \\
& =\left(1-t_{k}\right) g_{k}+\mathcal{A}_{k} t_{k}\left\|g_{k}\right\|+\mathcal{B}_{k} t_{k}^{2}\left\|g_{k}\right\|^{2}, \tag{2.12}
\end{align*}
$$

$\mathcal{A}_{k}$ and $\mathcal{B}_{k}$ as above. Equation (2.12) immediately yields the Taylor inequality

$$
\begin{equation*}
\left\|g_{k+1}\right\| \leq\left\|g_{k}\right\|\left\{\left(1-t_{k}\right)+\alpha_{k} t_{k}+\beta_{k} t_{k}^{2}\left\|g_{k}\right\|\right\} \tag{2.13}
\end{equation*}
$$

We will show the sequence $\left\|g_{k}\right\| \rightarrow 0$ by analyzing the term in braces.
Proposition 2.1. Let $\delta \in\left(0,1-\alpha_{0}\right), \alpha_{0} \in(0,1)$ and $t_{k}$ chosen as in (2.1) where

$$
\begin{equation*}
0 \leq \mathcal{K}_{k} \leq \mathcal{K}_{0} \tag{2.14}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathcal{K}_{k} \geq \frac{k_{1}^{2} k_{2}}{2\left(1-\alpha_{k}-\delta\right)}-\frac{1}{\left\|g_{k}\right\|} . \tag{2.15}
\end{equation*}
$$

Assume A1-A3 and all $\alpha_{k} \leq \alpha_{0}$. Then
(i) all $u_{k} \in S_{0}$, the sequence $\left\|g_{k}\right\|$ is strictly decreasing and $\left\|g_{k}\right\| \rightarrow 0$; furthermore,
(ii) $\left\|g_{k+1}\right\| /\left\|g_{k}\right\| \rightarrow 0$ if and only if $\alpha_{k} \rightarrow 0$, and for any fixed $p \in(0,1]$,

$$
\begin{equation*}
\left\|g_{k+1}\right\| \leq c_{1}\left\|g_{k}\right\|^{1+p} \tag{2.16}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
\alpha_{k} \leq c_{2}\left\|g_{k}\right\|^{p} \tag{2.17}
\end{equation*}
$$

for positive constants $c_{1}$ and $c_{2}$.
Proof. To show (i), suppose $u_{j} \in S_{0}$ and $\left\|g_{j}\right\|<\left\|g_{j-1}\right\|$ for $1 \leq j \leq k$. Since $u_{k+1}=u_{k}+t_{k} M_{k}^{-1} g_{k},\left\|u_{k+1}\right\| \leq\left\|u_{k}\right\|+k_{1}\left\|g_{k}\right\|$ so $u_{k+1} \in S_{1}$. Thus (2.11) and (2.15) imply

$$
\begin{equation*}
\mathcal{K}_{k} \geq \frac{\beta_{k}}{1-\alpha_{k}-\delta}-\frac{1}{\left\|g_{k}\right\|} \tag{2.18}
\end{equation*}
$$

Rearranging (2.18) and using (2.1) shows

$$
\begin{equation*}
\left(1-t_{k}\right)+\alpha_{k} t_{k}+\beta_{k} t_{k}^{2}\left\|g_{k}\right\| \leq 1-\delta t_{k} \tag{2.19}
\end{equation*}
$$

hence

$$
\begin{equation*}
\left\|g_{k+1}\right\| \leq\left(1-\delta t_{k}\right)\left\|g_{k}\right\| \leq\left(1-\delta t_{0}\right)\left\|g_{k}\right\| \tag{2.20}
\end{equation*}
$$

recalling (2.13). Equation (2.20) implies the conclusion (i).
Part (ii) follows from the pair of inequalities

$$
\begin{equation*}
\frac{\left\|g_{k+1}\right\|}{\left\|g_{k}\right\|} \leq\left(\mathcal{K}_{k}+k_{1}^{2} k_{2} / 2\right)\left\|g_{k}\right\|+\alpha_{k} \tag{2.21}
\end{equation*}
$$

and

$$
\begin{equation*}
\alpha_{k} \leq\left(1+\mathcal{K}_{k}\left\|g_{k}\right\|\right) \frac{\left\|g_{k+1}\right\|}{\left\|g_{k}\right\|}+\left(\mathcal{K}_{k}+k_{1}^{2} k_{2} / 2\right)\left\|g_{k}\right\| \tag{2.22}
\end{equation*}
$$

since $\left\|g_{k}\right\| \rightarrow 0$. Recalling $t_{k} \leq 1$, Equations (2.21)-(2.22) are immediate from (2.13) and the analogous inequality derived by transposing (2.12).

Note that (2.15) is satisfied for the constant sequence

$$
\begin{equation*}
\mathcal{K}_{k}=\mathcal{K}_{0} \quad \text { for all } k \tag{2.23}
\end{equation*}
$$

Furthermore (2.15) allows the choice $\mathcal{K}_{k}=0$ when

$$
\begin{equation*}
\left\|g_{k}\right\| \leq \frac{2\left(1-\alpha_{k}-\delta\right)}{k_{1}^{2} k_{2}} \tag{2.24}
\end{equation*}
$$

However, as noted in Section 1, (2.23) can be quite unsatisfactory, and we have chosen to force $t_{k} \rightarrow 1$ by using (2.1) rather than using a test to determine whether the choice $t_{k}=1$ is satisfactory as eventually guaranteed by (2.24); see Sections 3-4.

We will show later that the sequence $u_{k}$ converges to the root $u^{*}$ with $g\left(u^{*}\right)=0$. Recall that convergence is superlinear if

$$
\begin{equation*}
\left\|u_{k+1}-u^{*}\right\| \leq \eta_{k}\left\|u_{k}-u^{*}\right\| \quad \text { and } \eta_{k} \rightarrow 0 \tag{2.25}
\end{equation*}
$$

it is order $Q-(p+1),(p \in(0,1])$ if

$$
\begin{equation*}
\left\|u_{k+1}-u^{*}\right\| \leq c_{p}\left\|u_{k}-u^{*}\right\|^{p+1} \quad c_{p}>0 . \tag{2.26}
\end{equation*}
$$

Convergence is $R$-linear if

$$
\begin{equation*}
\left\|u_{k+1}-u^{*}\right\| \leq \eta_{k+1} \tag{2.27}
\end{equation*}
$$

and if

$$
\begin{equation*}
\eta_{k+1} \leq c \eta_{k}, \quad c \in(0,1) . \tag{2.28}
\end{equation*}
$$

To examine the nature of the convergence of $\left\{u_{k}\right\}$ we consider the relationship between $\left\|g_{k}\right\| \rightarrow 0$ and $\left\|u_{k}-u^{*}\right\| \rightarrow 0$. Consider the Taylor expansion

$$
\begin{aligned}
0 & =g\left(u^{*}\right)=g_{k}+g_{k}^{\prime}\left\{u^{*}-u_{k}\right\}+\int_{0}^{1}\left\{g^{\prime}\left(u_{k}+s\left(u^{*}-u_{k}\right)\right)-g_{k}^{\prime}\right\}\left\{u^{*}-u_{k}\right\} d s \\
(2.29) & =g_{k}+g_{k}^{\prime} t_{k} x_{k}+g_{k}^{\prime}\left\{u^{*}-u_{k+1}\right\}+\int_{0}^{1}\left\{g^{\prime}\left(u_{k}+s\left(u^{*}-u_{k}\right)\right)-g_{k}^{\prime}\right\}\left\{u^{*}-u_{k}\right\} d s
\end{aligned}
$$

Sine $g^{\prime}$ is continuous and invertible on $S_{0},\left\|\left(g_{k}^{\prime}\right)^{-1}\right\| \leq k_{3}$; rearranging the second inequality of (2.29) implies

$$
\begin{equation*}
\left\|u_{k+1}-u^{*}\right\| \leq k_{3}\left\{\left(1-t_{k}\right)\left\|g_{k}\right\|+t_{k} \alpha_{k}\left\|g_{k}\right\|+\left(k_{2} / 2\right)\left\|u_{k}-u^{*}\right\|^{2}\right\} . \tag{2.30}
\end{equation*}
$$

Letting $k_{5}=\sup _{u \in S_{1}}\left\|g^{\prime}(u)\right\|$, note

$$
\begin{equation*}
\left\|g_{k}\right\| \leq k_{5}\left\|u_{k}-u^{*}\right\| ; \tag{2.31}
\end{equation*}
$$

hence,

$$
\begin{equation*}
\left\|u_{k+1}-u^{*}\right\| \leq\left\|u_{k}-u^{*}\right\| k_{3}\left\{k_{5} \alpha_{k}+\left(k_{5}^{2} \mathcal{K}_{0}+k_{2} / 2\right)\left\|u_{k}-u^{*}\right\|\right\} \tag{2.32}
\end{equation*}
$$

using (2.1) and $t_{k} \leq 1$. Equation (2.32) shows that the convergence of $\left\{u_{k}\right\}$ to $u^{*}$ is superliner if $\alpha_{k} \rightarrow 0$ and is of order $Q-(p+1)$ if $\alpha_{k} \leq c_{2}\left\|g_{k}\right\|^{p}$, again using (2.31).

Suppose that in some set $S \subseteq S_{0}$ (2.31) can be extended to

$$
\begin{equation*}
k_{4}\left\|u_{k}-u^{*}\right\| \leq\left\|g_{k}\right\| \leq k_{5}\left\|u_{k}-u^{*}\right\| . \tag{2.33}
\end{equation*}
$$

Then under the conditions of Proposition 2.1, it is immediate from (2.20), (2.21), and (2.33) that the convergence of $\left\{u_{k}\right\} \operatorname{tp} u^{*}$ is:
(i) $R$-linear, and furthermore
(ii) superlinear if and only if $\left\|g_{k+1}\right\| /\left\|g_{k}\right\| \rightarrow 0$, and;
(iii) order $Q-(p+1)$ if and only if $\left\|g_{k+1}\right\| \leq c_{1}\left\|g_{k}\right\|^{p+1}, p \in(0,1]$.

In general (2.33) may not be valid on the entire set $S_{0}$. However (2.29) implies (2.33) for $u_{k}$ sufficiently close to $u^{*}$ as follows. The first inequality in (2.29) leads to

$$
\begin{equation*}
\left\|u_{k}-u^{*}\right\|\left\{1-\left(k_{2} k_{3} / 2\right)\left\|u_{k}-u^{*}\right\|\right\} \leq k_{3}\left\|g_{k}\right\| . \tag{2.34}
\end{equation*}
$$

Thus for any $\rho \in(0,1)$ (say $\rho=1 / 2$ ), (2.34) and (2.31) imply (2.33) for $S=S_{0} \cap S_{\rho}$ where

$$
\begin{equation*}
S_{\rho}=\left\{u \left\lvert\,\left\|u-u^{*}\right\| \leq \frac{2(1-\rho)}{k_{2} k_{3}}\right.\right\} \tag{2.35}
\end{equation*}
$$

and $k_{4}$ of (2.33) is $k_{4}=\rho k_{3}^{-1}$. Summarizing Proposition2.1 and the discussion involving (2.25)-(2.35) we have

Theorem 2.2. Under the conditions of Proposition 2.1
(i) there exists a $u^{*} \in S_{0}$ with $u^{*}=\lim u_{k}$ and $g\left(u^{*}\right)=0$;
(ii) on $S_{0}$ the convergence of $\left\{u_{k}\right\}$ to $u^{*}$ is superlinear or order $O-(p+1)$ if $\alpha_{k} \rightarrow 0$ or $\alpha_{k} \leq c_{2}\left\|g_{k}\right\|^{p}$, respectively;
(iii) on any set $S=S_{0} \cap S_{\rho}$ as in (2.35), the convergence of $\left\{u_{k}\right\}$ to $u^{*}$ is at least $R$-linear; it is superlinear or order $Q-(p+1)$ if and only if $\alpha_{k} \rightarrow 0$ or $\alpha_{k} \leq c_{2}\left\|g_{k}\right\|^{p}$, respectively.
Proof. It remains only to show (i), which we establish by showing that $\left\{u_{k}\right\}$ is a Cauchy sequence. But since

$$
\begin{equation*}
\left\|u_{k+j}-u_{k}\right\|=\left\|\sum_{i=k}^{k+j-1} t_{i} x_{i}\right\| \leq k_{1} \sum_{i}\left\|g_{i}\right\| \tag{2.36}
\end{equation*}
$$

and the $\left\|g_{k}\right\| \rightarrow 0$ with $\left\|g_{k+1}\right\| \leq c\left\|g_{k}\right\|, c<1,\left\{u_{k}\right\}$ is clearly Cauchy with a limit $u^{*}$ in the closed set $S_{0}$. Continuity of $g$ implies that $g\left(u^{*}\right)=\lim g\left(u_{k}\right)=0$.

We now have the following global result.
Theorem 2.3. Let $G: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a homeomorphism. Suppose $g^{\prime}$ is Lipschitz on closed bounded sets and A2 is satisfied. Then, given any $u_{0}$, the sequence $u_{k}$ of (1.2)-(1.3) with $t_{k}$ as in (2.1) and $\mathcal{K}_{k}$ as in (2.15) converges to $u^{*}$ as in Theorem 2.2.

Proof. Since $g^{\prime}$ is Lipschitz on closed bounded sets, $g^{\prime}$ is continuous on $\mathbb{R}^{n}$. Thus $\|g(u)\| \rightarrow \infty$ as $\|u\| \rightarrow \infty$ since $g$ is a homeomorphism ([9], page 137). Hence $S_{0}$ of (2.2) is bounded for any $u_{0}$ and A1 and A3 are satisfied. The result now follows from Theorem 2.2.

As mentioned in Section 1, early investigations by Dennis and Moré examined higher order convergence of approximate Newton methods. In our notation, they characterized convergence by studying the quantity $\left\|\left(M_{k}-g^{\prime}\left(u^{*}\right)\right) x_{k}\right\| /\left\|x_{k}\right\|$ where they chose $x_{k}=u_{k+1}-u_{k}\left(t_{k}=1\right)$. Their results can be recast in the framework of Theorem 2.2 as the following

ThEOREM 2.4. (cf. [6], pages 51-52). In addition to the conditions of Proposition2.1 let $\left\|M_{k}\right\| \leq k_{6}$. Then
(i) on any set $S=S_{0} \cap S_{\rho}, S_{\rho}$ as in (2.35), convergence of $\left\{u_{k}\right\}$ to $u^{*}$ is superlinear if and only if

$$
\begin{equation*}
\frac{\left\|\left(M_{k}-g^{\prime}\left(u^{*}\right)\right) x_{k}\right\|}{\left\|x_{k}\right\|} \rightarrow 0 \tag{2.37}
\end{equation*}
$$

while on $S_{0}$, (2.37) implies superlinear convergence;
(ii) on $S$ convergence is $Q-(p+1)$ if and only if

$$
\begin{equation*}
\frac{\left\|\left(M_{k}-g^{\prime}\left(u^{*}\right)\right) x_{k}\right\|}{\left\|x_{k}\right\|} \leq \mu_{p}\left\|x_{k}\right\|^{p}, \quad p \in(0,1] . \tag{2.38}
\end{equation*}
$$

Proof. The conclusion follows from the pair of inequalities

$$
\begin{equation*}
\alpha_{k} \leq k_{1} \frac{\left\|\left(M_{k}-g^{\prime}\left(u^{*}\right)\right) x_{k}\right\|}{\left\|x_{k}\right\|}+k_{1} k_{2}\left\|u_{k}-u^{*}\right\|, \tag{2.39}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\left\|\left(M_{k}-g^{\prime}\left(u^{*}\right)\right) x_{k}\right\|}{\left\|x_{k}\right\|} \leq \alpha_{k} \frac{\left\|g_{k}\right\|}{\left\|x_{k}\right\|}+\frac{k_{2}}{k_{4}}\left\|g_{k}\right\|, \tag{2.40}
\end{equation*}
$$

which can be derived from the definitions of $\alpha_{k}$ and the constants $k_{i}$. For example, to show the "only if" part of (ii), we first note that $Q-(p+1)$ convergence implies
(Theorem 2.2) that $\alpha_{k} \leq c_{2}\left\|g_{k}\right\|^{p} \leq c_{2} k_{6}^{p}\left\|x_{k}\right\|^{p}$. Hence by (2.40)

$$
\begin{aligned}
\frac{\left\|\left(M_{k}-g^{\prime}\left(u^{*}\right)\right) x_{k}\right\|}{\left\|x_{k}\right\|} & \leq c_{2} k_{6}^{p+1}\left\|x_{k}\right\|^{p}+\frac{k_{2} k_{6}}{k_{4}}\left\|x_{k}\right\| \\
& \leq \mu_{p}\left\|x_{k}\right\|^{p}
\end{aligned}
$$

The other conclusions follow similarly.
We conclude this section with some remarks concerning the generality of the analysis presented here.

Remark R1: If Hölder continuity, i.e.,

$$
\begin{equation*}
\left\|g^{\prime}(u)-g^{\prime}(v)\right\| \leq k_{e}\|u-v\|^{e}, \quad v, u \in S_{1}, \quad e \in(0,1) \tag{2.41}
\end{equation*}
$$

replaces Lipschitz continuity in A 3 , the above analysis remains valid with minor modifications including the restriction of the order exponent $p$ to $p \in(0, e]$. In fact, if $g^{\prime}$ is only uniformly continuous on $S_{1}$ (continuous in $\mathbb{R}^{n}$ ), then

$$
\begin{equation*}
\left\|g^{\prime}(u)-g^{\prime}(v)\right\| \leq w(\|u-v\|) \tag{2.42}
\end{equation*}
$$

where $w(t)$ is the modulus of continuity for $g^{\prime}$ on $S_{1}$ (see [9], page 64). Again much of the analysis remains valid; however, it is now only possible to obtain superlinear convergence. The restriction on $\mathcal{K}_{k}$ analogous to (2.15) in this case is

$$
\begin{equation*}
w\left(\frac{k_{1}\left\|g_{k}\right\|}{1+\mathcal{K}_{k}\left\|g_{k}\right\|}\right) \leq \frac{1-\alpha_{k}-\delta}{k_{1}} \tag{2.43}
\end{equation*}
$$

showing it is possible to choose $0 \leq \mathcal{K}_{k} \leq \mathcal{K}_{0}$ since $w$ is an isotone continuous function with $w(0)=0$. See Daniel ([3], Section 4.2, Chapter 8) and Sanberg [11] for discussions of the role of uniform continuity in similar contexts.

Remark R2: Note that (1.2) and (1.3) can be replaced by any procedure for determining $x_{k}$ such that

$$
\begin{equation*}
0<\left\|x_{k}\right\| \leq k_{1}\left\|g_{k}\right\| \tag{2.44}
\end{equation*}
$$

In all our applications, however, $x_{k}$ can be shown to derive from $g_{k}$ by a linear relationship of the form (1.2). Furthermore, the bound on $\left\|M_{k}^{-1}\right\|$ usually follows from the continuity and invertibility of $g^{\prime}$ on $S_{0}$ in addition to convergence assumptions on the inner process which determines $x_{k}$ (see Section 4).

In the spirit of R1, it is also possible to generalize (2.44) to

$$
\begin{equation*}
0<\left\|x_{k}\right\| \leq k_{1}\left\|g_{k}\right\|^{s}, \quad s \in(1 / 2,1] \tag{2.45}
\end{equation*}
$$

for $g^{\prime}$ Lipschitz on $S_{1}$. If $g^{\prime}$ is less smooth, $s$ must be suitably restricted.
Remark R3: As we have seen in Theorems 2.2 and 2.4, the inability to extend (2.33), in general, to the entire set $S_{0}$ leads to somewhat disquieting technicalities concerning the necessary conditions on the convergence to zero of the sequence $\left\{\alpha_{k}\right\}$. In the important special case that $g(u)$ is uniformly monotine on $S_{0}$, i.e.,

$$
\begin{equation*}
(g(u)-g(v))^{T}(u-v) \geq k_{7}(u-v)^{T}(u-v) \tag{2.46}
\end{equation*}
$$

the Cauchy-Schwarz inequality implies (in the 2-norm) that

$$
\begin{equation*}
k_{7}\left\|u-u^{*}\right\|_{2} \leq\|g(u)\|_{2} \quad \text { on } S_{0} \tag{2.47}
\end{equation*}
$$

Hence (2.33) is valid on $S_{0}$ and the appropriate statements in Theorems 2.2 and 2.4 can be simplified. In an algorithmic setting, however, note that statements (i) and (ii) of Theorem 2.2 are the real content of the result.
3. Algorithm. We now turn to the computational aspects of the analysis described in Section 2. In particular, we consider the problem of determining the $\mathcal{K}_{k}$ of (2.1) such that (2.15), or more importantly (2.18), is satisfied.

Note that inequality (2.20) can be rewritten as

$$
\begin{equation*}
\delta \leq\left(1-\frac{\left\|g_{k+1}\right\|}{\left\|g_{k}\right\|}\right) \frac{1}{t_{k}} \tag{3.1}
\end{equation*}
$$

Since the right-hand side of (3.1) is easily computed, and since we may choose $\delta \in$ $\left(0,1-\alpha_{0}\right)$, equation (3.1) is a convenient test. Failure to satisfy (3.1) implies $\mathcal{K}_{k}$ fails to satisfy (2.18). We then increase $\mathcal{K}_{k}$ and compute new values for $t_{k}$ and $g_{k}$. These increases will eventually lead to $\mathcal{K}_{k}$ satisfying (2.15), (2.18), and (3.1), and this process leads to convergence as in Proposition 2.1.

Consider the choice of $\mathcal{K}_{0}$. Given a guess, say $\mathcal{K}_{0}=0$, each failure of the test (3.1) requires a function evaluation of $g(u)$ to compute a new $g_{1}$. This aspect of the procedure has the flavor of a line search, but with one important difference. Once a value of $\mathcal{K}_{0}$ has been accepted, one might reasonable expect to pass (3.1) for $\mathcal{K}_{k}=\mathcal{K}_{0}$ on almost all subsequent iterations $k$. However, note that as $\left\|g_{k}\right\|$ decreases, the righthand side of (2.18) decreases, suggesting the possibility of taking $\mathcal{K}_{k} \leq \mathcal{K}_{k-1}$. If the $\mathcal{K}_{k}$ decrease in a orderly manner (for example $\mathcal{K}_{k}=\mathcal{K}_{k-1} / 10$ ), we anticipate a process which uses only one function evaluation on most steps. In fact decreasing the $\mathcal{K}_{k}$ can be important; we have found than an excessively large value of $\mathcal{K}_{0}$ will cause the convergence of $t_{k} \rightarrow 1$ to be much slower than necessary, delaying the onset of the observed superlinear convergence, and possibly resulting in many iterations.

The above discussion motivates the following algorithm.

## Algorithm Global

(1) input $u_{0}, \delta \in\left(0,1-\alpha_{0}\right)$
(2) $\mathcal{K} \leftarrow 0, k \leftarrow 0$; compute $g_{0},\left\|g_{0}\right\|$
(3) compute $x_{k}$
(4) $t_{k} \leftarrow\left(1+\mathcal{K}\left\|g_{k}\right\|\right)^{-1}$
(5) compute $u_{k+1}, g_{k+1},\left\|g_{k+1}\right\|$
(6) if $\left(1-\left\|g_{k+1}\right\| /\left\|g_{k}\right\|\right) t_{k}^{-1}<\delta$
(7) then $\{$ if $\mathcal{K}=0$, then $\mathcal{K} \leftarrow 1$; else $\mathcal{K} \leftarrow 10 \mathcal{K}\}$; GOTO (4)
(8) else $\{\mathcal{K} \leftarrow \mathcal{K} / 10 ; k \leftarrow k+1\}$
(9) if converge, then return; else GOTO (3)

In Global, failure to satisfy (3.1) causes $\mathcal{K}$ to be increased in line (7). Each failure requires on additional function evaluation on line (5). On line (8), we take $\mathcal{K} / 10$ as the initial estimate for $\mathcal{K}_{k+1}$. Alternatively, we have considered $\mathcal{K}_{k+1}=\mathcal{K}_{k} 4^{j-k-1}$ where $j$ is the last index resulting in a failure of the test on line (6). In practice, we have found these methods for decreasing $\mathcal{K}$ to be a reasonable compromise between the (possibly) conflicting goals of having $t_{k} \rightarrow 1$ quickly and having (3.1) satisfied on the first function evaluation for most steps.

The procedure for increasing $\mathcal{K}$ is also important, and we have found a procedure other than the relatively simple one given on line (7) to be advantageous. In this scheme, one specifies a priori the maximum number of function evaluations to be allowed on a given step, say $\ell\left(\right.$ typically $\ell=10$ ), The trial values of $\mathcal{K}$ denoted $\mathcal{K}_{k, j}$,

$$
\begin{equation*}
\mathcal{K}_{k, j}=\left(\frac{1}{\left\|g_{k}\right\|}+\frac{\mathcal{K}_{k-1}}{10}\right)\left(\frac{\left\|x_{k}\right\|}{\mu\left\|u_{k}\right\|}\right)^{((j-1) /(\ell-1))^{2}}-\frac{1}{\left\|g_{k}\right\|}, \quad \text { for } \frac{\left\|x_{k}\right\|}{\left\|u_{k}\right\|}>\mu \tag{3.2}
\end{equation*}
$$

This corresponds to the easily implemented formulae

$$
\begin{align*}
& t_{k, 1}=\left(1+\mathcal{K}\left\|g_{k}\right\| / 10\right)^{-1}  \tag{3.3}\\
& t_{k, j}=t_{k, 1}\left(\mu\left\|u_{k}\right\|\left\|x_{k}\right\|\right)^{((j-1) /(\ell-1))^{2}}, \quad 2 \leq j \leq \ell
\end{align*}
$$

We take $\mu$ to be a constant on the order of the machine epsilon $\times 10^{3}$. For small values of $j \geq 2, t_{k, j}$ represents a modest decrease of $t_{k, j-1}$. As $j$ increases, $t_{k, j}$ decreases more rapidly until $t_{k, \ell}\left\|x_{k}\right\|=\left\|u_{k}\right\| \mu t_{k, 1}$. If (3.1) fails for $t_{k, \ell}$, the calculation is terminated and an error flag set. Equations (3.3) represent a compromise between the conflicting goals of increasing $\mathcal{K}$ slowly (so as not to accept a value which is excessively large) and of finding an acceptable value in few function evaluations.

On line (3) we have not detailed the computation involving (1.2). If $M_{k}=g_{k}^{\prime}$ in (1.2), Global is a damped Newton method and $\alpha_{k}=0$ for all $k$ (disregarding round off). Alternatively, (1.2) may represent an iterative process for solving $g_{k}^{\prime} x_{k}+g_{k}=$ 0 terminated when $\alpha_{k}$ satisfies some tolerance such as (2.9)-(2.10). Such damped Newton methods and other approximate Newton methods are outlined in the following section.
4. Applications. In this section we present several applications for the results in the previous sections. In particular, we show how Newton-iterative methods, Newtonapproximate Jacobian methods, and other Newton-like methods fit within our global approximate Newton framework.
4.1. Newton-Iterative Methods. Suppose that $x_{k}$ in line (3) of algorithm Global is computed by using an iterative method to solve the Newton equations

$$
\begin{equation*}
g_{k}^{\prime} w_{k}=-g_{k} \tag{4.1}
\end{equation*}
$$

For example, we might use a standard iterative method such as SOR or a NewtonRichardson method where $g_{k}^{\prime}$ in (4.1) is replaced by a previous Jacobian $g_{k^{\prime}}^{\prime}$. The Newton-Richardson choice is useful when a (possibly sparse) $L U$ factorization of the Jacobian is relatively expensive. Hence the Jacobian is factored infrequently, and in outer iterations where the factorization is not computed, we iterate to approximately solve (4.1) using the last computed factorization. (see [12]).

In all such Newton-iterative methods ([9], Section 7.4), we suppose $g_{k}^{\prime}$ has a uniformly convergent splitting on $S_{0}$; i.e.,

$$
\begin{equation*}
g_{k}^{\prime}=A_{k}-B_{k} \tag{4.2}
\end{equation*}
$$

with $\left\|H_{k}\right\| \leq \rho_{0}<1$ for all $k$, where

$$
\begin{equation*}
H_{k}=A_{k}^{-1} B_{k}=I-A_{k}^{-1} g_{k}^{\prime} \tag{4.3}
\end{equation*}
$$

We then compute $x_{k}$ by computing the inner iteration

$$
\begin{equation*}
A_{k} x_{k, m}=B_{k} x_{k, m-1}-g_{k} \tag{4.4}
\end{equation*}
$$

until $m=m_{k}$, taking $x_{k, 0}=0$ and setting $x_{k}=x_{k, m_{k}}$. Note that (4.4) can be rewritten as

$$
\begin{equation*}
A_{k}\left(x_{k, m}-x_{k, m-1}\right)=-\left(g_{k}^{\prime} x_{k, m-1}+g_{k}\right) \tag{4.5}
\end{equation*}
$$

Using induction and (4.3), it can be shown that the $x_{k, m}$ in (4.4) satisfy

$$
\begin{equation*}
g_{k}^{\prime}\left(I-H_{k}^{m}\right)^{-1} x_{k, m}=-g_{k}, \quad m \geq 1 \tag{4.6}
\end{equation*}
$$

Hence we may identify $M_{k}$ of (1.2) with $M_{k}=g_{k}^{\prime}\left(I-H_{k}^{m_{k}}\right)^{-1}$, and these $M_{k}$ satisfy (2.3) since $\left\|\left(g_{k}^{\prime}\right)^{-1}\right\|$ is bounded on $S_{0}$ and $\left\|I-H_{k}^{m_{k}}\right\| \leq 2$.

Notice that the right hand side of (4.5) contains the Newton residual which suggest defining the quantities

$$
\begin{equation*}
\alpha_{k, m} \equiv \frac{\left\|g_{k}+g_{k}^{\prime} x_{k, m}\right\|}{\left\|g_{k}\right\|} \tag{4.7}
\end{equation*}
$$

in analogy with (2.6) and the quantities $x_{k, m}$. The $\alpha_{k, m}$ are easily computed, certainly when the iteration proceeds as in (4.5) rather then (4.4). Since we have assumed that the $A_{k}$ and $B_{k}$ are a convergent splitting, $\alpha_{k, m} \rightarrow 0$ as $m \rightarrow \infty$. Thus to obtain convergence as discussed in Section 2, we stop the inner iteration when $\alpha_{k, m}$ attains the desired tolerance $\alpha_{k} \equiv \alpha_{k, m_{k}}$. For example, to obtain orer $Q-(p+1)$ superlinear convergence, $p \in(0,1]$, we stop after $m_{k}$ iterations where

$$
\begin{equation*}
\alpha_{k, m} \leq \alpha_{0}\left(\frac{\left\|g_{k}\right\|}{\left\|g_{0}\right\|}\right)^{p}, \quad \alpha_{0} \in(0,1) \tag{4.8}
\end{equation*}
$$

as in (2.9)-(2.10).
Note that

$$
\begin{equation*}
g_{k}+g_{k, m}^{\prime} x_{k, m}=\hat{H}_{k}^{m} g_{k} \tag{4.9}
\end{equation*}
$$

where $H_{k}=g_{k}^{\prime} \hat{H}_{k}\left(g_{k}^{\prime}\right)^{-1}$; this implies

$$
\begin{equation*}
\alpha_{k, m} \leq\left\|\hat{H}_{k}\right\|^{m} \tag{4.10}
\end{equation*}
$$

Assuming (4.1) is an equality with $\left\|\hat{H}_{k}\right\|<1$ and that equation (2.16)-(2.17) are equalities, we see that

$$
\begin{equation*}
m_{k}=\frac{\log c_{2}\left\|g_{k}\right\|^{p}}{\log \left\|\hat{H}_{k}\right\|} . \tag{4.11}
\end{equation*}
$$

Asymptotically (as $k \rightarrow \infty$ ) we expect $\left\|\hat{H}_{k+1}\right\| \approx\left\|\hat{H}_{k}\right\|$; again assuming equality and using (2.16)-(2.17) and (4.11) shows

$$
\begin{equation*}
m_{k+1} \sim(1+p) m_{k} \tag{4.12}
\end{equation*}
$$

Hence for $k$ sufficiently large we can expect the number of inner iterations per outer iteration to approximately increase by a factor of $(1+p)$.

In the preceding general analysis of Newton-iterative methods all $\alpha_{k}=\alpha_{k, m}$ are possibly nonzero. For the special case of a Newton-Richardson method a decision is made at the beginning of the $k$-th outer iteration whether to factor $g_{k}^{\prime}$ thus doing an "exact" damped Newton iteration. Such a factorization implies $\alpha_{k}=0$; otherwise the
inner iteration corresponds to a splitting with $A_{k}=g_{k^{\prime}}^{\prime}, k^{\prime}<k$, where $A_{k}$ has been previously factored.

It is not difficult to decide when to refactor: one should refactor when the total cost of inner iterations using the factored $g_{k^{\prime}}^{\prime}$ just surpasses the cost of a new factorization. For example, using nested dissection on an $n \times n$ mesh cost approximately $10 n^{3}$ operations for a factorization and $5 n^{2} \log _{2} n$ operations for a backsolution. Thus approximately $2 n / \log _{2} n$ inner iterations compared with a new factorization. Note however that initially these $2 n / \log _{2} n$ inner iterations will be part of several outer iterations monitored by the $\alpha_{k, m}$; each time a new such outer iteration is started, $g_{k}^{\prime}$ is computed (but not factored) for use in the right hand side of (4.5). The relative time $T_{1}$ corresponding to $10 n^{3}$ and $T_{2}$ corresponding to $5 n^{2} \log _{2} n$ can often be timed dynamically using a "clock routine" and need not be known a priori. As a final remark, note that superlinear convergence will require an increasing number of inner iterations, and ultimately, the inner iteration time will surpass the cost of a new factorization. However, in practice when only a modest overall accuracy is required Newton-Richardson methods can prove to be highly effective, and we have found such cases.
4.2. Newton-Approximate Jacobian Methods. When the partial derivatives required for the computation of $g_{k}^{\prime}$ are unavailable or expensive to compute, it is common to approximate $g_{k}^{\prime}$, perhaps using finite differences ([6], page 49, [9], pages 185-186). We denote such an approximation to $g_{k}^{\prime}$ by $\tilde{g}_{k}^{\prime}$.

We assume that the $\tilde{g}_{k}^{\prime}$ satisfy

$$
\begin{equation*}
\left\|g_{k}^{\prime}-\tilde{g}_{k}^{\prime}\right\| \leq \frac{\delta_{k}}{k_{1}} \tag{4.13}
\end{equation*}
$$

for $\delta_{k}<1$. Let $\mathcal{A}_{k}$ of (2.6) be written as

$$
\begin{equation*}
\mathcal{A}_{k}=\tilde{\mathcal{A}}_{k}+\Delta_{k} \tag{4.14}
\end{equation*}
$$

where

$$
\begin{equation*}
\tilde{\mathcal{A}}_{k}=\frac{g_{k}+\tilde{g}_{k}^{\prime} x_{k}}{\left\|g_{k}\right\|} ; \quad \Delta_{k}=\frac{\left(g_{k}^{\prime}-\tilde{g}_{k}^{\prime}\right) x_{k}}{\left\|g_{k}\right\|} \tag{4.15}
\end{equation*}
$$

Following (2.6), let $\tilde{\alpha}_{k}=\left\|\tilde{\mathcal{A}}_{k}\right\|$ and note that $\left\|\Delta_{k}\right\| \leq \delta_{k}$; hence

$$
\begin{equation*}
\alpha_{k} \leq \tilde{\alpha}_{k}+\delta_{k} \tag{4.16}
\end{equation*}
$$

If $x_{k}$ is obtained by the linear system

$$
\begin{equation*}
\tilde{g}_{k}^{\prime} x_{k}=-g_{k} \tag{4.17}
\end{equation*}
$$

then $\tilde{\alpha}_{k}=0$, corresponding to $M_{k}=\tilde{g}_{k}^{\prime}$ in (1.2). Alternatively (4.17) can be solved approximately, perhaps by a Newton-iterative method as in Section 4.1; then $\tilde{\alpha}_{k} \neq 0$ in general.

If all $\tilde{\alpha}_{k}=0$ the $\delta_{k}$ play the role of $\alpha_{k}$ in Section2. For example, if all $\tilde{\alpha}_{k}=0$ and

$$
\begin{equation*}
\delta_{k} \leq \delta_{0}\left(\frac{\left\|g_{k}\right\|}{\left\|g_{0}\right\|}\right)^{p} ; \quad p \in(0,1], \delta_{0}<1 \tag{4.18}
\end{equation*}
$$

then the Newton-approximate Jacobian scheme will converge with order $Q-(p+1)$. More generally, let $\tilde{\alpha}_{k} \neq 0$ and $\delta_{k}$ satisfy (4.18) with

$$
\begin{equation*}
\tilde{\alpha}_{k} \leq \max \left\{\tilde{\alpha}_{0}\left(\frac{\left\|g_{k}\right\|}{\left\|g_{o}\right\|}\right)^{q}, \delta_{k}\right\}, \quad \tilde{\alpha}_{0}+\delta_{0}<1 \tag{4.19}
\end{equation*}
$$

In a typical situation we might have $p=0, q=1$, and $\delta_{0}$ small; i.e., we compute the approximation to $g_{k}^{\prime}$ to a fixed accuracy and use an iterative method to solve (4.17) approximately. For the first few outer iterations, relatively few inner iterations will be required. The inner iterations will increase until $\delta_{k}$ becomes the larger of the two terms on the right-hand side of (4.19). From this point onward, approximately a constant number of inner iterations will be used, and the asymptotic outer convergence will be $R$-linear (Theorem 2.2). It may not be easy to estimate or compute $\delta_{k}$; although, if we are computing $\tilde{g}_{k}^{\prime}$ to a fixed accuracy, we expect all $\delta_{k}$ to be approximately equal. Since $\delta_{k}$ enters the computation only through (4.19), we see that its main purpose is to prevent useless inner iterations. If nothing is known about $\delta_{k}$, we can set all $\delta_{k}=\epsilon$ in (4.19), $\epsilon$ a sufficiently small number. Convergence of the outer iteration can loosely be described as superlinear at the beginning and ultimately linear depending on the actual values of $\delta_{k}$. The choice of $\epsilon$ could have a significant effect on the total computation cost and may require some experimentation.
4.3. Two Parameter Damping. In an earlier investigation, [1], we studied the Newton-like method

$$
\begin{align*}
\left(I / s_{k}+g_{k}^{\prime}\right) x_{k} & =-g_{k}  \tag{4.20}\\
u_{k+1} & =u_{k}+x_{k} \tag{4.21}
\end{align*}
$$

motivating the method by considering Euler integration on the autonomous system of ODEs

$$
\begin{equation*}
\frac{d u}{d t}+g(u)=0 ; \quad u(0)=u_{0} \tag{4.22}
\end{equation*}
$$

Under appropriate conditions, including the uniform monotonicity of $g(u)$ on $\mathbb{R}^{n}$ in the form

$$
\begin{equation*}
x^{T} g^{\prime}(u) x \geq k_{7} x^{T} x \tag{4.23}
\end{equation*}
$$

we showed that it is possible to obtain global quadratic convergence by forcing $s_{k}\left\|g_{k}\right\|$ to be a sufficiently small constant for all $k$. We used a norm reducing argument similar to the analysis of Section 2. Here we sketch a more general treatment using the results of Section 2.

Consider the iteration (1.2)-(1.3) with $t_{k}$ as in (2.1), $\lambda_{k} \geq 0$, and

$$
\begin{equation*}
M_{k}=\lambda_{k}\left\|g_{k}\right\| I+g_{k}^{\prime} \tag{4.24}
\end{equation*}
$$

We call such a method two parameter damping because the $\lambda_{k}$ as well as the $t_{k}$ limit the change $\left(u_{k+1}-u_{k}\right)$. It is immediate from (4.24) that the $\alpha_{k}$ of (2.6) satisfy

$$
\begin{equation*}
\alpha_{k}=\lambda_{k}\|x\|_{k} \tag{4.25}
\end{equation*}
$$

If the $\left\|M_{k}\right\|$ are uniformly bounded (for all $\lambda_{k}$ ) as in (2.3), then

$$
\begin{equation*}
\alpha_{k} \leq k_{1} \lambda_{k}\|g\|_{k} \tag{4.26}
\end{equation*}
$$

This shows that the $\lambda_{k}$ can be chosen such that $\alpha_{0}<1$ and $\alpha_{k} \leq \alpha_{0}$ as in Proposition 2.1. Furthermore, the $Q$-quadratic convergence (2.16) is also a consequence of (4.26).

The above discussion can be made more precise by seeking a relation between the uniform bound $k_{1}$ and $\lambda_{k}$. Suppose that $g(u)$ is uniformly monotone as in (4.23) on $S_{0}$. Then, in the 2-norm,

$$
\begin{equation*}
\left\|M_{k}^{-1}\right\|_{2} \leq\left(\lambda_{k}\left\|g_{k}\right\|_{2}+k_{7}\right)^{-1} \tag{4.27}
\end{equation*}
$$

and

$$
\begin{equation*}
\alpha_{k}=\lambda_{k}\left\|x_{k}\right\| \leq \frac{\lambda_{k}\left\|g_{k}\right\|_{2}}{\lambda_{k}\left\|g_{k}\right\|_{2}+k_{7}}<\min \left(1, k_{7}^{-1} \lambda_{k}\left\|g_{k}\right\|_{2}\right) \tag{4.28}
\end{equation*}
$$

Consider the sequence of $\lambda_{k}$ with $0 \leq \lambda_{k} \leq \lambda_{0}$. Note that the $\alpha_{k}$ are easily computable and all $\alpha_{k}<1$, although it may not be the case that all $\alpha_{k} \leq \alpha_{0}$. This requires a minor modification in Proposition 2.1. We will assume that for each $k \geq 0, \delta \in\left(0,1-\hat{\alpha}_{k}\right)$ where $\hat{\alpha}_{k}=\max _{j \leq k} \alpha_{j}$. Note the sup $\hat{\alpha}_{k}<1$ by (4.28) and the induction argument leading to (2.20). Since we do not know $\hat{\alpha}_{k}$ a priori, we may be required to change (decrease) $\delta$ dynamically as the iteration proceeds; that is, if for some $\alpha_{k}, \delta \geq 1-\alpha_{k}$. Such decreases in $\delta$ cause no convergence problems since (3.1) continues to hold if $\delta$ is decreased in subsequent iterations.

It is possible to show that for $\lambda_{k}=\lambda_{0}, k \geq 0$, with $\lambda_{0}$ sufficiently large, $t_{k}$ can be chosen as $t_{k}=1$ for all $k$. One starts with (2.13) and uses (4.28) and $\beta_{k} \leq\left(k_{2} / 2\right)\left(\lambda_{0}\left\|g_{k}\right\|_{2}+k_{7}\right)^{-2}$ as in [1], Section 3. However, such an analysis (without the explicit damping parameter $t_{k}$ ) is more existential than the two parameter analysis sketched above, mainly since the sufficient decrease parameter analogues to $\delta$ of (2.20) depends on the usually unknown constant $k_{7}$. In practice, there seems to be no advantage in using only $\lambda$ damping, whereas two parameter damping may be advantageous when the $g_{k}^{\prime}$ themselves are numerically ill-conditioned.
5. Numerical Remarks. The methods described in this work and our earlier presentation [1], are part of a larger study aimed at solving effectively the coupled partial differential equations arising in semiconductor device modelling. These equations often take the form

$$
\begin{align*}
-\Delta u+e^{u-v}-e^{w-u} & =k(x, y)  \tag{5.1}\\
-\nabla \cdot\left(\mu_{n} e^{u-v} \nabla v\right) & =0  \tag{5.2}\\
-\nabla \cdot\left(\mu_{p} e^{w-u} \nabla w\right) & =0 \tag{5.3}
\end{align*}
$$

Here $u, v$, and $w$ are functions of $(x, y) \in D \subseteq \mathbb{R}^{2}$, as are the known functions $\mu_{n}, \mu_{p}$ and $k(x, y)$, and $D$ is a union of rectangles. The function $k(x, y)$ is the doping profile of the device; (5.1) is a nonlinear Poisson equation and (5.2)-(5.3) are continuity equations.

Equations (5.1)-(5.3) have be attacked numerically on two discretization fronts. Finite differences are used in W. Fichtner's simulation package; the now routine solution of the coupled equations is reported in $[8,7]$. In this package, Newton-Richardson and Newton-block SOR methods (as in Section 4) have proved to be particularly effective.

The device equations, especially (5.1), have also been attacked by a nonlinear multilevel iteration package using piecewise linear elements on triangles. This package has been designed concurrently with developing analysis presented in this paper and
is an extension of the linear package described in [2]. The package presently solves a single nonlinear PDE of the form

$$
\begin{equation*}
-\nabla \cdot(a(x, y) \nabla u)+f\left(u, u_{x}, u_{y}\right)=0 \tag{5.4}
\end{equation*}
$$

on a connected region $\Omega$ in $\mathbb{R}^{2}$ with standard elliptic boundary conditions; the formal generalization to the case

$$
\begin{equation*}
-\nabla \cdot\left(a\left(x, y, u, u_{x}, u_{y}\right) \nabla u\right)+f\left(u, u_{x}, u_{y}\right)=0 \tag{5.5}
\end{equation*}
$$

is straightforward. A special Newton-multilevel iterative method, along the lines discussed in Section 4, is used to solve the discrete equations. Details will be presented elsewhere.

To illustrate the use of the Newton-multilevel iteration package and Algorithm Global, consider, as in [1], Section 4, the $p-n$ junction problem of the form (5.1) above. The functions $v, w$, and $k$ are given, and the domain and boundary conditions are shown in Figure 5.1. Recall that the doping profile $k(x, y)$, and the solution gradient, $\nabla u(x, y)$, vary over several orders of magnitude in a small region near the junction, and there is a notable singularity due to the change in boundary conditions along the upper boundary.


FIG. 5.1. $p-n$ junction problem with boundary conditions
We consider only the level-one nonuniform grid with $n=25$ vertices (unknowns) and a (very poor) initial guess, $u_{0}=0$. (The higher levels are less interesting.) We use Algorithm Global as in Section 4 with the modification (3.2)-(3.3) in line (7). The $x_{k}$ are computed by a sparse $L U$ factorization of $g_{k}^{\prime}$. The convergence trace is presented in Table 5.1.

The relatively large number of iterations necessary in this experiment compared with the experiments reported in [1], Section 4 , is a direct consequence of taking $u_{0}=0$ rather than attempting to even roughly interpolate the boundary values $\rho_{0}$ and $\rho_{1}$ as we did before. This also leads to more searching (evals) than might otherwise be expected.

As a cautionary remark, we report that failing to dynamically change $\mathcal{K}$ (i.e., taking $\mathcal{K}_{k}=\mathcal{K}_{0}$ for all $k$ ) led to time overrun termination after $k=320$, $t_{k}=$ $1.106(-4)$ and $\left\|g_{k}\right\|=4.57(6)$ in the same experiment.

## REFERENCES

[1] R. E. Bank and D. J. Rose, Parameter selection for Newton-like methods applicable to nonlinear partial differential equations, SIAM J. Numer. Anal., 17 (1980), pp. 806-822.

| $\left\\|g_{0}\right\\|=4.738 \cdot 10^{6} ;$ evals $\equiv$ evaluations of $g(u)$ |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $k$ | $t_{k}$ | $\left\\|g_{k}\right\\|$ |  | $\left\\|u_{u}-u_{k-1}\right\\| /\left\\|u_{k}\right\\|$ |
| 1 | $1.07(-4)$ | $4.73(6)$ | $9.28(3)$ | 6 |
| 2 | $7.30(-4)$ | $4.70(6)$ | $8.11(2)$ | 2 |
| 3 | $7.31(-3)$ | $4.56(6)$ | $9.16(1)$ | 1 |
| 4 | $7.05(-2)$ | $3.41(6)$ | $1.12(1)$ | 1 |
| 5 | $2.81(-2)$ | $3.31(6)$ | $1.33(1)$ | 4 |
| 6 | $2.30(-1)$ | $2.21(6)$ | $6.26(-1)$ | 1 |
| 7 | $2.44(-1)$ | $5.27(5)$ | 1.33 | 3 |
| 8 | $9.31(-1)$ | $5.88(4)$ | $4.37(-1)$ | 1 |
| 9 | $6.15(-2)$ | $3.64(4)$ | 2.33 | 4 |
| 10 | $1.52(-1)$ | $3.11(4)$ | 1.56 | 3 |
| 11 | $6.78(-1)$ | $1.29(4)$ | 1.41 | 1 |
| 12 | $9.81(-1)$ | $5.73(3)$ | $1.39(-1)$ | 1 |
| 13 | $9.99(-1)$ | $4.08(3)$ | $1.28(-1)$ | 1 |
| 14 | $7.39(-2)$ | $3.79(3)$ | $4.39(-1)$ | 4 |
| 15 | $3.48(-1)$ | $3.00(3)$ | $2.67(-1)$ | 2 |
| 16 | $8.71(-1)$ | $1.75(3)$ | $1.36(-1)$ | 1 |
| 17 | $3.18(-1)$ | $9.51(2)$ | $2.83(-1)$ | 3 |
| 18 | $8.95(-1)$ | $5.18(2)$ | $3.49(-2)$ | 1 |
| 19 | $7.73(-1)$ | $3.58(2)$ | $1.98(-2)$ | 2 |
| 20 | $7.66(-1)$ | $1.03(2)$ | $1.41(-2)$ | 2 |
| 21 | $9.83(-1)$ | $1.37(1)$ | $3.26(-3)$ | 1 |
| 22 | $1 .^{a}$ | 1.44 | $2.24(-4)$ | 1 |
| 23 | $1 .^{a}$ | $1.57(-1)$ | $2.36(-5)$ | 1 |
| 24 | $1 .^{a}$ | $7.67(-6)$ | $2.90(-6)$ | 1 |
| 25 | $1 .^{a}$ | $4.27(-7)$ | $1.40(-10)$ | 1 |
| $a$ | place rounding |  |  |  |
|  |  |  | TABLE | 5.1 |

[2] R. E. Bank and A. H. Sherman, PLTMG user's guide, Tech. Rep. CNA-152, Center for Numerical Analysis, University of Texas at Austin, 1979.
[3] J. W. Daniel, The Approximate Minimization of Functionals, Prentice-Hall, Englewood Cliffs, NJ, 1971.
[4] R. S. Dembo, S. C. Eisenstat, and T. Steihaug, Inexact-Newton methods, Tech. Rep. SOM Working Paper Series number 47, Yale University, 1981.
[5] J. E. Dennis and J. J. Moré, A characterization of superlinear convergence and its application to quasi-Newton methods, Math. Comp., 28 (1974), pp. 549-560.
[6] —, Quasi-Newton methods: motivation and theory, SIAM Rev., 19 (1977), pp. 46-89.
[7] W. Fichtner and D. J. Rose, On the numerical solution of nonlinear PDEs arising from semiconductor device modelling, in Elliptic Problem Solvers, (M. H. Schultz, ed.), Academic Press, New York, 1980.
[8] - Numerical semiconductor device simulation, Tech. Rep. Technical memo, Bell Laboratories, 1981.
[9] J. M. Ortega and W. C. Rheinboldt, Iterative Solution on Nonlinear Equations in Several Variables, Academic Press, New York, 1970.
[10] I. W. Sandberg, One Newton-direction algorithms and diffeomorphisms, Tech. Rep. Technical memo, Bell Laboratories, 1980.
[11] ——, Diffeomorphisms and Newton-direction algorithms, Bell Sys. Tech. J., (1981).
[12] A. H. Sherman, On Newton-iterative methods for the solution of systems of nonlinear equations, SIAM J. Numer. Anal., 15 (1978), pp. 755-771.
[13] A. Wouk, A Course of Applied Functional Analysis, J. Wiley and Sons, New York, 1979.


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