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# DYNAMICS OF THE GINZBURG-LANDAU EQUATIONS OF SUPERCONDUCTIVITY 

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#### Abstract

This article is concerned with the dynamical properties of tolutions of the time-dependent Ginzburg-Landau (TDGL) equations of superconductivity. It is shown that the TDGL equations define a dynamical process when the applied magnetic field varies with time, and a dynamical system when the applied magnetic field is stationary. The dynamical system describes the large-time asymptotic behavior: Every solution of the TDGL equations is attracted to a set of stationary solutions, which are divergence free. These results are obtained in the " $\phi=-\omega(\nabla \cdot A)$ " gauge, which reduces to the standard " $\phi=-\nabla \cdot \boldsymbol{A}$ " gauge if $\omega=1$ and to the zero-electric potential gauge if $\omega=0$; the treatment captures both in a unified framework. This gauge forces the London gauge, $\nabla \cdot \boldsymbol{A}=0$, for any stationary solution of the TDGL equations.


Keywords. Ginzburg-Landau equations, superconductivity, gauge, weak solotions, global existence, uniqueness, dynamical process, global attractor.

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## 1 Introduction

In this article, we are concerned with the dynamical properties of solutions of the time-dependent Ginzburg-Landau (TDGL) equations of superconductivity. While the emphasis is on the formal mathematical aspects of the equations, we make every effort to comply with the physical nature of the problem. We make no simplifications for the convenience of mathematics, and our rigorous treatment is motivated by known facts from physics. We show that the TDGL equations define a dynamical process when the applied magnetic field varies with time, and a dynamical system when the applied magnetic field is stationary. We work consistently in the " $\dot{\phi}=-\omega(\nabla \cdot \boldsymbol{A})$ " gauge introduced in [1] and [2] and deduce by logical arguments the ramifications for the zero-electric potential gauge $\phi=0$. The " $\phi=-\omega(\nabla \cdot \boldsymbol{A})^{\prime}$ gauge enables us to rigorously establish the large-time asymptotic behavior and make the connection with solutions of the time-independent GL equations of superconductivity.

### 1.1 Ginzburg-Landau Model of Superconductivity

In the Ginzburg-Landau theory of phase transitions [3], the state of a superconducting material near the critical temperature is described by a complexvalued order parameter $\psi$, a real vector-valued vector potential $\boldsymbol{A}$, and, when the system changes with time, a real-valued scalar potential $\phi$. The latter is a diagnostic variable; $\psi$ and $\boldsymbol{A}$ are prognostic variables, whose evolution is governed by a system of coupled differential equations,

$$
\begin{align*}
\eta\left(\frac{\partial}{\partial t}+i \kappa \phi\right) \psi & =-\left(\frac{i}{\kappa} \nabla+\boldsymbol{A}\right)^{2} \psi+\left(1-|\psi|^{2}\right) \psi  \tag{1.1}\\
\frac{\partial \boldsymbol{A}}{\partial t}+\nabla \phi & =-\nabla \times \nabla \times \boldsymbol{A}+J_{s}+\nabla \times H \tag{1.2}
\end{align*}
$$

where $J_{s}$, the supercurrent density, is a nonlinear function of $\psi$ and $\boldsymbol{A}$,

$$
\begin{equation*}
J_{s} \equiv J_{s}(\psi, \boldsymbol{A})=\frac{1}{2 i \kappa}\left(\psi^{*} \nabla \psi-\psi \nabla \psi^{*}\right)-|\psi|^{2} \boldsymbol{A}=-\operatorname{Re}\left[\psi^{*}\left(\frac{i}{\kappa} \nabla+\boldsymbol{A}\right) \psi\right] . \tag{1.3}
\end{equation*}
$$

The order parameter can be thought of as the wave function of the center-ofmass motion of the "superelectrons" (Cooper pairs), whose density is $n_{s}=|\psi|^{2}$ and whose flux is $J_{s}$.

The system of Eqs. (1.1)-(1.3), henceforth referred to as the "TDGL equations," must be satisfied everywhere in $\Omega$, the region occupied by the superconducting material, and at all times $t>0$. We assume that $\Omega$ is a bounded
domain in $\mathbf{R}^{n}$ with a boundary $\partial \Omega$ of class $C^{1.1}$. That is, $\Omega$ is an open and connected set whose boundary $\partial \Omega$ is a compact ( $n-1$ )-manifold described by Lipschitz-continuously differentiable charts. We consider two- and threedimensional problems ( $n=2$ and $n=3$, respectively). The vector potential $\boldsymbol{A}$ takes its values in $\mathbf{R}^{n}$. The vector $H$ represents the (externally) applied magnetic field, which is a given function of space and time; like $A$, it takes its values in $\mathbf{R}^{\boldsymbol{n}}$.

The parameters in the TDGL equations are $\eta$. a (dimensionless) friction coefficient, and $\kappa$, the (dimensionless) Ginzburg-Landau parameter. The former measures the temporal, the latter the spatial rate of change of the order parameter relative to the vector potential. As usual, $\nabla \equiv \operatorname{grad}, \nabla \times \equiv$ curl, $\nabla \cdot \equiv \operatorname{div}$, and $\nabla^{2}=\nabla \cdot \nabla \equiv \Delta ; i$ is the imaginary unit, and a superscript * denotes complex conjugation. Sometimes, we use the symbol $\partial_{t}$ to denote the partial derivative $\partial / \partial t$.

The boundary conditions associated with the TDGL equations are

$$
\begin{equation*}
\boldsymbol{n} \cdot\left(\frac{i}{\kappa} \nabla+\boldsymbol{A}\right) \psi+\frac{i}{\kappa} \gamma \psi=0 \quad \text { and } \quad n \times(\nabla \times \boldsymbol{A}-H)=0 \quad \text { on } \partial \Omega, \tag{1.4}
\end{equation*}
$$

where $\boldsymbol{n}$ is the local outer unit normal to $\partial \Omega$. They must be satisfied at all times $t>0$. The function $\gamma$ is Lipschitz continuous on $\partial \Omega$, and $\gamma(x) \geq 0$ for $x \in \partial \Omega$.

The vector $\boldsymbol{E}=-\partial_{t} \boldsymbol{A}-\nabla \phi$ is the electric field and $\boldsymbol{B}=\nabla \times \boldsymbol{A}$ the magnetic induction. Therefore, Eq. (1.2) is Ampère's law, $\nabla \times B=J$, where the total current $J$ is the sum of a "normal" current $J_{n}=E$, the supercurrent $J_{s}$, and the transport current $J_{t}=\nabla \times H$. The normal current obeys Ohm 's law $J_{n}=\sigma_{n} E$; the "normal conductivity" coefficient $\sigma_{n}$ is equal to one in the adopted system of units. The difference $\boldsymbol{M}=\boldsymbol{B}-\boldsymbol{H}$ is known as the magnetization. The trivial solution ( $\dot{\psi}=0, B=H, E=0$ ) represents the normal state, where all superconducting properties have been lost.

The TDGL equations generalize the original GL equations to the timedependent case. The GL equations themselves embody in a most simple way the macroscopic quantummechanical nature of the superconducting state. The generalization, first proposed by Schmid [4], was analyzed by Gor'кov and Eliashberg [5] in the context of the microscopic Bardeen-Cooper-Schrieffer (BCS) theory of superconductivity. Although the validity of the TDGL equations seems to be limited to a narrow range of temperatures near the critical temperature, $T_{c}$, the equations are being used extensively and successfully in large-scale numerical simulations to study vortex dynamics in type-II superconductors at Argonne National Laboratory [6, 7, 8]. We refer the reader to the physics literature $[9,10,11]$ for further details.

### 1.2 Previous Work

The TDGL equations have been the object of several recent mathematical studies. Elliott and Tang [12] proved the existence and miqueness of solutions of the TDGL equations in two-dimensional domains under some complicated mathematical boundary conditions, using a time-discretization procedure. Their article was followed by another article by Ta.vg [13], who applied the same methods to the TDGL equations with fixed total magnetic flux. Du [14]. using a finite-element approach. established the existence and uniqueness of weak solutions in two- and three-dimensional domains under the assumption that the order parameter is initially bounded in $L^{\infty}(\Omega)$. The same results were obtained independently and at about the same time by CHEN, Hoffmann, and Liang [15], who used the Leray-Schauder fixed-point theorem. Du adopted the zero-electric potential gauge ( $\dot{\phi}=0$ ), Chen, Hoffmann, and Liang the " $\phi=-\nabla \cdot A^{\text {" }}$ gauge for their analysis.

In [16], Llang and TANg considered the dynamics of the TDGL equations in bounded domains in $\mathbf{R}^{3}$, assuming the London gauge, $\nabla \cdot \boldsymbol{A}=0$, at all times. They claimed to prove the existence of a dynamical system, but failed to verify the continuous dependence of the solution operator on the initial data. Therefore, it is not evident that the solution operator actually defines a dynamical system. Moreover, the limiting relation displayed in the proof of [16, Theorem 6.1] does not follow from [17, Theorem 4.3.4], as claimed by the authors.

Finally, Ta.vg and Wang in their recent article [18] observed that the TDGL equations have the same type of nonlinearity as the Navier-Stokes equations for an incompressible fluid. They adapted the methods developed for the Navier-Stokes equations to prove the existence of strong solutions in two and three dimensions, with no boundedness assumption on the initial data, and the existence of weak solutions in two dimensions. However, they too failed to verify the continuous dependence of the solution operator on the initial data.

### 1.3 Outline of Present Work

In the present article, we continue work begun in [1] and [2]. As noted by TAKÁC [2], the natural gauge for the study of the dynamical properties of the TDGL equations is the " $\phi=-\nabla \cdot A^{\prime \prime}$ gauge. In this gauge, the TDGL equations generate a dynamical system, and every stationary solution satisfies the London gauge, $\nabla \cdot \boldsymbol{A}=0$. A generalization of this gauge was introduced by Fleckinger-Pellé and Kaper in [1]. The generalization replaces the
standard " $\dot{\phi}=-\nabla \cdot \boldsymbol{A}$ " gauge by the " $\sigma=-\omega(\nabla \cdot \boldsymbol{A})$ " gauge, where $\omega$ is any nonnegative number. This gauge captures the standard " $\dot{\phi}=-\nabla \cdot A$ " gauge and the zero-electric potential gauge in a unified framework. Applying the methods developed by TAKÁc in [2], we are able to establish rigorously the existence of a dynamical process for the TDGL equations if the applied magnetic field is time dependent, and the existence of a dynamical system if it is time independent. In the latter case. we prove that every stationary solution of the TDGL equations satisfies the London gauge. This result indicates bow the stationary solutions of the TDGL equations can be connected to the solutions of the time-independent GL equations.

Following is an outline of the article.
Section 2 contains preliminary material. In Section 2.1 we derive some auxiliary identities from the TDGL equations; in Section 2.2 we introduce the " $\phi=-\omega(\nabla \cdot \boldsymbol{A})$ " gauge. Unless mentioned otherwise, all further arguments refer to the TDGL equations in this gauge. In Section 2.3 we give various estimates that follow from an energy-type functional.

Section 3 defines the abstract initial-value problem for the TDGL equations. Section 3.1 presents the notation. In Section 3.2 we introduce the applied vector potential to homogenize the boundary conditions and in Section 3.3 give the abstract initial-value problem. In Section 3.4 we prove a regularity result involving the applied vector potential, which eventually determines the regularity of a mild solution of the abstract initial-value problem.

Section 4 contains our results in three theorems, each with a corollary. Theorem 1 (Section 4.1) gives an existence and uniqueness result, Theorem 2 (Section 4.2) a regularity result. Both theorems hold when the applied magnetic field varies with time. A corollary of Theorem 2 is the existence of a dynamical process. Specializing to the case of a time-independent magnetic field, we obtain a dynamical system whose properties are given in Theorem 3 (Section 4.3).

The proofs of the theorems are given in Section 5.

## 2 Preliminaries

In this section we establish several auxiliary identities, which follow from the TDGL equations (1.1)-(1.4). We also introduce the gauge choice and define an energy-type functional for the TDGL equations.

### 2.1 Auxiliary Identities

The TDGL model of superconductivity is basically a system of semilinear parabolic equations. This is most easily seen if. in Eqs. (1.1) and (!.2). one uses the identities

$$
\begin{equation*}
-\left(\frac{i}{\kappa} \nabla+\boldsymbol{A}\right)^{2} \psi=\frac{1}{\kappa^{2}} \Delta \psi-\frac{2 i}{\kappa}(\nabla \dot{\psi}) \cdot \boldsymbol{A}-\frac{i}{\kappa} \dot{\psi}(\nabla \cdot \boldsymbol{A})-\psi|\boldsymbol{A}|^{2} \tag{2.1}
\end{equation*}
$$

and

$$
\begin{equation*}
-\nabla \times \nabla \times A=\Delta A-\nabla(\nabla \cdot A) \tag{2.2}
\end{equation*}
$$

Many of the methods developed for such systems are indeed applicable to the TDGL equations. But, as we will see in the following analysis, the TDGL equations have several distinct features that make them mathematically interesting in their own right and different from, say, the Navier-Stokes equations.

The curl of a gradient vanishes, so the TDGL equations do not change if we replace $H$ by $\boldsymbol{H}^{\prime}=H+\nabla \Phi$, for any (sufficiently smooth) real scalar-valued function $\Phi$ of position and time. If $\Phi=0$ on $\partial \Omega$, we also have $n \times H=n \times H^{\prime}$ on $\partial \Omega$, so the boundary conditions do not change either. In particular, if we take $\Phi$ at any time as the (unique) solution of the Dirichlet problem for Poisson's equation $\Delta \Phi=-\nabla \cdot H$, we have $\nabla \cdot H^{\prime}=0$ at all times. Hence, there is no loss of generality if, from now on, we assume that the applied magnetic field $\boldsymbol{H}$ is divergence free,

$$
\begin{equation*}
\nabla \cdot H=0 \quad \text { in } \Omega . \tag{2.3}
\end{equation*}
$$

The quantity $n_{s}=|\psi|^{2}$ corresponds to the superelectron density. Its evolution is governed by the equation

$$
\begin{equation*}
\eta \frac{\partial|\psi|^{2}}{\partial t}=-2 \operatorname{Re}\left[\psi^{*}\left(\frac{i}{\kappa} \nabla+\boldsymbol{A}\right)^{2} \psi\right]+2\left(1-|\psi|^{2}\right)|\dot{\psi}|^{2} \tag{2.4}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
\eta \frac{\partial|\psi|^{2}}{\partial t}=\frac{1}{\kappa^{2}} \Delta|\psi|^{2}-2\left|\left(\frac{i}{\kappa} \nabla+A\right) \psi\right|^{2}+2\left(1-|\psi|^{2}\right)|\psi|^{2} . \tag{2.5}
\end{equation*}
$$

Clearly, if the inequality $|\psi| \leq 1$ is satisfied on $\Omega$ at $t=0$, it is satisfied at all later times. Note that the scalar potential $\phi$ does not figure in Eq. (2.4).

The divergence of a curl vanishes, so Eq. (1.2) implies the identity

$$
\begin{equation*}
\nabla \cdot\left(\frac{\partial A}{\partial t}+\nabla \phi\right)=\nabla \cdot J_{s} \quad \text { in } \Omega \tag{2.6}
\end{equation*}
$$

An expression for $\nabla \cdot J_{s}$ is easily obtained by taking the divergence of Eq. (1.3).

$$
\begin{equation*}
\nabla \cdot J_{s}=-\kappa \operatorname{Im}\left[\psi^{*}\left(\frac{i}{\kappa} \nabla+A\right)^{2} \psi\right] \tag{2.7}
\end{equation*}
$$

From this expression and Eq. (1.1) we obtain the more interesting expression

$$
\begin{equation*}
\nabla \cdot J_{s}=\eta \kappa^{2}\left[\frac{1}{2 i \kappa}\left(\psi^{*} \frac{\partial \psi}{\partial t}-\psi \frac{\partial \psi^{*}}{\partial t}\right) \div \phi|v|^{2}\right] . \tag{2.8}
\end{equation*}
$$

An immediate consequence of the definition (1.3) of $J_{s}$ and the first boundary condition in Eq. (1.4) is that $n \cdot J_{s}=0$ on $\partial \Omega$.

By assumption, $\partial \Omega$ is locally the level surface (or curve) of a $C^{1.1}$-function $\Phi: \mathbf{R}^{n} \rightarrow \mathbf{R}$. Hence, the unit normal vector is given by $n=|\nabla \Phi|^{-1} \nabla \Phi$, where $\nabla \Phi$ is nonvanishing and Lipschitz continuous near every point of $\partial \Omega$. Consequently, $\boldsymbol{n} \cdot(\nabla \times \boldsymbol{n})=0$. According to the second boudary condition in Eq. (1.4), $\boldsymbol{\nabla} \times \boldsymbol{A}-\boldsymbol{H}$ and $\boldsymbol{n}$ are colinear on $\partial \Omega$. Therefore, it must be the case that $\boldsymbol{n} \cdot \nabla \times(\nabla \times \boldsymbol{A}-\boldsymbol{H})=0$ on $\partial \Omega$.

When we combine the two identities $\boldsymbol{n} \cdot \boldsymbol{J}_{s}=0$ and $\boldsymbol{n} \cdot \nabla \times(\nabla \times \boldsymbol{A}-\boldsymbol{H})=0$ with Eq. (1.2), we see that $\boldsymbol{n} \cdot\left(\partial_{t} \boldsymbol{A}+\nabla \phi\right)=0$ on $\partial \Omega$. Therefore, any solution of the TDGL equations is such that

$$
\begin{equation*}
n \cdot\left(\frac{\partial A}{\partial t}+\nabla \phi\right)=0 \quad \text { and } \quad n \cdot J_{s}=0 \quad \text { on } \partial \Omega \tag{2.9}
\end{equation*}
$$

These identities express the physical fact that the electric field and the supercurrent are always tangential to the surface of the superconductor.

### 2.2 Gauge Choice

The TDGL equations are invariant under the gauge transformation

$$
\begin{equation*}
\mathcal{G}_{\chi}:(\psi, \boldsymbol{A}, \phi) \mapsto\left(\psi e^{i \kappa \chi}, \boldsymbol{A}+\nabla \chi, \dot{\phi}-\partial_{t} \chi\right) \tag{2.10}
\end{equation*}
$$

The gauge $\chi$ can be any (sufficiently smooth) real scalar-valued function of position and time. For the present investigation we adopt the " $\phi=-\omega(\nabla \cdot A)$ " gauge introduced in [1]. This gauge is determined by taking $\chi \equiv \chi_{\omega}(x, t)$ as the (unique) solution of the boundary-value problem

$$
\begin{align*}
\left(\partial_{t}-\omega \Delta\right) \chi & =\phi+\omega(\nabla \cdot \boldsymbol{A}) \quad \text { in } \Omega \times(0, \infty)  \tag{2.11}\\
\omega(\boldsymbol{n} \cdot \nabla \chi) & =-\omega(\boldsymbol{n} \cdot \boldsymbol{A}) \quad \text { on } \partial \Omega \times(0, \infty) \tag{2.12}
\end{align*}
$$

subject to a suitable initial condition, $\gamma(\cdot, 0)=\chi_{0}$ in $\Omega$. Here. $\omega$ is a real nonnegative parameter. In this gauge we have, at all times $t>0$, the identities

$$
\begin{equation*}
\dot{\phi}+\omega(\nabla \cdot \boldsymbol{A})=0 \quad \text { in } \Omega, \quad \omega(\boldsymbol{n} \cdot \boldsymbol{A})=0 \quad \text { on } \partial \Omega \tag{2.13}
\end{equation*}
$$

The gauge choice fixes the potential $\varphi$ in terms of the order parameter $\psi$ and the vector potential $\boldsymbol{A}$. The differential equations (1.1) and (1.2) reduce to

$$
\begin{align*}
\eta \frac{\partial v}{\partial t} & =-\left(\frac{i}{\kappa} \nabla+\boldsymbol{A}\right)^{2} \psi+i \eta \kappa \omega \dot{w}(\nabla \cdot \boldsymbol{A})+\left(1-|\dot{w}|^{2}\right) w \text { in } \Omega \times(0, \infty)  \tag{2.14}\\
\frac{\partial \boldsymbol{A}}{\partial t} & =-\nabla \times \nabla \times \boldsymbol{A}+\omega \nabla(\nabla \cdot \boldsymbol{A})+J_{s}+\nabla \times H \quad \text { in } \Omega \times(0, \infty) \tag{2.15}
\end{align*}
$$

where $J_{s}$ is again given by Eq. (1.3), and the boundary conditions (1.4) to $\boldsymbol{n} \cdot \nabla \psi+\gamma \psi=0, \quad \omega(\boldsymbol{n} \cdot \boldsymbol{A})=0, \quad \boldsymbol{n} \times(\nabla \times \boldsymbol{A}-\boldsymbol{H})=\mathbf{0} \quad$ on $\partial \Omega \times(0, \infty)$.

The boundary-value problem (2.14)-(2.16) is strongly parabolic for $\omega>0$. It becomes degenerate for $\omega=0$.

The scalar potential $\phi$ does not figure in the evolution equation (2.4), so the gauge choice does not affect the observation that $|\psi| \leq 1$ on $\Omega$ at all times $t>0$ if the inequality is satisfied at $t=0$. (Cf. the "maximum modulus principle" in Section 4.1, Theorem 1.)

The auxiliary identity (2.6) and the expression (2.8) involve $\phi$. In the $" \phi=-\omega(\nabla \cdot \boldsymbol{A})$ " gauge, Eq. (2.6) assumes the form

$$
\begin{equation*}
\left(\partial_{t}-\omega \Delta\right)(\nabla \cdot A)=\nabla \cdot J_{s} \tag{2.17}
\end{equation*}
$$

where

$$
\begin{equation*}
\nabla \cdot J_{s}=\eta \kappa^{2}\left[\frac{1}{2 i \kappa}\left(\psi^{*} \frac{\partial \psi}{\partial t}-\psi \frac{\partial \psi^{*}}{\partial t}\right)-\omega|\psi|^{2}(\nabla \cdot A)\right] \tag{2.18}
\end{equation*}
$$

Finally, the identities (2.9) reduce to

$$
\begin{equation*}
\boldsymbol{n} \cdot \nabla(\nabla \cdot \boldsymbol{A})=0 \quad \text { and } \quad n \cdot J_{s}=0 \quad \text { on } \partial \Omega \tag{2.19}
\end{equation*}
$$

Unless mentioned otherwise, all further arguments refer to the TDGL equations in the " $\dot{\Phi}=-\omega(\nabla \cdot \boldsymbol{A})$ " gauge.

### 2.3 Energy-Type Functionals

Consider the functional $E_{\omega} \equiv E_{\omega}[\psi, A]$,

$$
E_{\omega}[\psi, \boldsymbol{A}]=\int_{\Omega}\left[\left|\left(\frac{i}{\kappa} \nabla+\boldsymbol{A}\right) \psi\right|^{2}+\frac{1}{2}\left(1-|\psi|^{2}\right)^{2}+2 \omega(\nabla \cdot \boldsymbol{A})^{2}\right.
$$

$$
\begin{equation*}
\left.+|\nabla \times A-H|^{2}\right] \mathrm{d} x+\int_{\ni \Omega} \gamma\left|\frac{i}{\kappa} \psi\right|^{2} \mathrm{~d} \sigma(x) \tag{2.20}
\end{equation*}
$$

If $\psi$ and $\boldsymbol{A}$ satisfy Eqs. (2.14)-(2.16), then the time derivative of $E_{\omega}$ is

$$
\begin{gather*}
\frac{\mathrm{d} E_{\omega}}{\mathrm{d} t}=-2 \int_{\Omega}\left[\eta\left|\frac{\partial \psi}{\partial t}-i \kappa \omega \psi(\nabla \cdot \boldsymbol{A})\right|^{2}+\left|\frac{\partial A}{\partial t}\right|^{2}+\omega^{2}|\nabla(\nabla \cdot \boldsymbol{A})|^{2}\right] \mathrm{d} x \\
-2 \int_{\Omega} \frac{\partial H}{\partial t} \cdot(\nabla \times \boldsymbol{A}-H) \mathrm{d} x \tag{2.21}
\end{gather*}
$$

If $\partial_{t} H=0$ (stationary applied magnetic field), the expression in the right member is negative semidefinite, $E_{i \nu}$ is a Liapunov functional, and $E_{\omega}(t) \leq$ $E_{\omega}(0)$ for all $t \geq 0$. In general, the applied magnetic field is not stationary, and $E_{\omega}$ is not necessarily bounded by a constant. However, as the following lemma shows, $E_{\omega}(t)$ can still be estimated in terms of the quantity $P(t)$,

$$
\begin{equation*}
P(t)=\int_{0}^{t}\left(\int_{\Omega}\left|\partial_{t} \boldsymbol{H}(x, s)\right|^{2} \mathrm{~d} x\right)^{1 / 2} \mathrm{~d} s \tag{2.22}
\end{equation*}
$$

Lemma 1 If $E_{\omega} \equiv E_{\omega}(t)$ exists and is finite, and $P(T)<\infty$ for some $T>0$, then

$$
\begin{gather*}
E_{\omega}(t)+2 \int_{0}^{t} \int_{\Omega}\left[\eta\left|\frac{\partial \psi}{\partial t}-i \kappa \omega \psi(\nabla \cdot \boldsymbol{A})\right|^{2}+\left|\frac{\partial \boldsymbol{A}}{\partial t}\right|^{2}+\omega^{2}|\nabla(\nabla \cdot \boldsymbol{A})|^{2}\right] \mathrm{d} x \mathrm{~d} t^{\prime} \\
\leq\left(\left(E_{\omega}(0)\right)^{1 / 2}+P(t)\right)^{2}, \quad t \in(0, T) \tag{2.23}
\end{gather*}
$$

Proof. It follows from Eq. (2.21) and the Cauchy-Schwarz inequality that

$$
\begin{equation*}
\frac{\mathrm{d} E_{\omega}}{\mathrm{d} t} \leq 2\left(\int_{\Omega}\left|\frac{\partial H}{\partial t}\right|^{2} \mathrm{~d} x\right)^{1 / 2}\left(\int_{\Omega}|\nabla \times A-H|^{2} \mathrm{~d} x\right)^{1 / 2} \leq 2 \frac{\mathrm{~d} P}{\mathrm{~d} t}\left(E_{\omega}(t)\right)^{1 / 2} \tag{2.24}
\end{equation*}
$$

Hence, $\mathrm{d} E_{\omega}^{1 / 2} / \mathrm{d} t \leq \mathrm{d} P / \mathrm{d} t$. Upon integration, we obtain

$$
\begin{equation*}
E_{\omega}(t) \leq\left(\left(E_{\omega}(0)\right)^{1 / 2}+P(t)\right)^{2}, \quad t \in(0, T) \tag{2.25}
\end{equation*}
$$

To obtain the inequality (2.23), we use Eq. (2.21) again, this time including the first integral, and apply the estimate (2.25),

$$
\begin{gather*}
\frac{\mathrm{d} E_{\omega}}{\mathrm{d} t}+2 \int_{\Omega}\left[\eta\left|\frac{\partial \psi}{\partial t}-i \kappa \omega \psi(\nabla \cdot \boldsymbol{A})\right|^{2}+\left|\frac{\partial \boldsymbol{A}}{\partial t}\right|^{2}+\omega^{2}|\nabla(\nabla \cdot \boldsymbol{A})|^{2}\right] \mathrm{d} x \\
\leq 2 \frac{\mathrm{~d} P}{\mathrm{~d} t}\left(E_{\omega}(t)\right)^{1 / 2} \leq 2 \frac{\mathrm{~d} P}{\mathrm{~d} t}\left(\left(E_{\omega}(0)\right)^{1 / 2}+P(t)\right) \tag{2.26}
\end{gather*}
$$

The inequality (2.23) follows upon integration.

Lemma 2 Assume that $M=$ ess $\sup \{|\dot{\psi}(x, t)|:(x, t) \in \Omega \times(0, T)\}<\infty$. Then

$$
\begin{align*}
& 2 \int_{0}^{t} \int_{\Omega}\left[\eta\left|\frac{\partial \psi}{\partial t}\right|^{2}+\left|\frac{\partial A}{\partial t}\right|^{2}+\omega^{2}|\nabla(\nabla \cdot \boldsymbol{A})|^{2}\right] \mathrm{d} x \mathrm{~d} t^{\prime} \\
& \leq\left(3+\eta \kappa^{2} \omega M^{2} t\right)\left(\left(E_{\omega}(0)\right)^{1 / 2}+P(t)\right)^{2}, \quad t \in(0, T) \tag{2.27}
\end{align*}
$$

whenever the terms in the inequality are well defined.

Proof. Using the elementary inequality $|a|^{2} \leq 2\left(a-\left.b\right|^{2}+|b|^{2}\right)$ and the inequality ( 2.23 ), we obtain
$\int_{0}^{t} \int_{\Omega} \eta\left|\frac{\partial \psi}{\partial t}\right|^{2} \mathrm{~d} x \mathrm{~d} t^{\prime} \leq\left(\left(E_{\omega}(0)\right)^{1 / 2}+P(t)\right)^{2}+\pi \kappa^{2} \omega_{\cdot} M^{2} \int_{0}^{t} \int_{\Omega} 2 \omega(\nabla \cdot \boldsymbol{A})^{2} \mathrm{~d} x \mathrm{~d} t^{\prime}$,
where

$$
\int_{0}^{t} \int_{\Omega} 2 \omega(\nabla \cdot \boldsymbol{A})^{2} \mathrm{~d} x \mathrm{~d} t^{\prime} \leq \int_{0}^{t} E_{\omega}\left(t^{\prime}\right) \mathrm{d} t^{\prime} \leq t\left(\left(E_{\omega}(0)\right)^{1 / 2}+P(t)\right)^{2}
$$

The remaining terms in the left member of the inequality (2.27) have already been estimated by $\left(\left(E_{\omega}(0)\right)^{1 / 2}+P(t)\right)^{2}$ in Lemma 1, inequality (2.23).

The additional term $2 \omega(\nabla \cdot A)^{2}$ in the functional $E_{\omega}$ has no basis in physics. Indeed, $E_{\omega}$ is not an energy functional, unless $\omega=0$. If $\omega=0, E_{\omega}$ reduces to the Ginzburg-Landau free-energy functional,

$$
\begin{align*}
E_{0}[\psi, \boldsymbol{A}]=\int_{\Omega}\left[\left\lvert\,\left(\frac{i}{\kappa} \nabla+\right.\right.\right. & \left.\boldsymbol{A})\left.\psi\right|^{2}+\frac{1}{2}\left(1-|\psi|^{2}\right)^{2}+|\nabla \times \boldsymbol{A}-H|^{2}\right] \mathrm{d} x \\
& +\int_{\partial \Omega} \gamma\left|\frac{i}{\kappa} \psi\right|^{2} \mathrm{~d} \sigma(x) \tag{2.28}
\end{align*}
$$

The gauge restriction (2.13) reduces to $\phi=0$ in $\Omega$, and the Euler equations and natural boundary conditions associated with $E_{0}$ are

$$
\begin{gather*}
-\left(\frac{i}{\kappa} \nabla+\boldsymbol{A}\right)^{2} \psi+\left(1-|\psi|^{2}\right) \psi=0 \quad \text { in } \Omega,  \tag{2.29}\\
-\nabla \times \nabla \times \boldsymbol{A}+J_{s}+\nabla \times \boldsymbol{H}=0 \quad \text { in } \Omega,  \tag{2.30}\\
\boldsymbol{n} \cdot\left(\frac{i}{\kappa} \nabla+\boldsymbol{A}\right) \psi+\gamma \frac{i}{\kappa} \psi=0 \quad \text { and } \quad \boldsymbol{n} \times(\nabla \times \boldsymbol{A}-\boldsymbol{H})=0 \quad \text { on } \partial \Omega . \tag{2.31}
\end{gather*}
$$

This is the time-independent Ginzburg-Landau model. A manuscript documenting the relationship between this and the TDGL model is in preparation.

## 3 Functional Formulation

In this section, we formulate the TDGL equations as an abstract evolution equation in a Hilbert space. The formulation requires a reduction of tize boundary conditions to homogeneous form, which is accomplished by the introduction of an applied vector potential.

### 3.1 Notation

The symbol $C$ denotes a generic positive constant, not necessarily the same at different instances. All Banach spaces are real; the (real) dual of a Banach space $X$ is denoted by $X^{\prime}$.
$L^{p}(\Omega)$, for $1 \leq p \leq \infty$, is the usual Lebesgue space, with norm $\|\cdot\|_{L^{p}}$; $(\cdot, \cdot)$ is the inner product in $L^{2}(\Omega)$. $W^{m, 2}(\Omega)$, for nonnegative integer $m$, is the usual Sobolev space, with norm $\|\cdot\| w_{m, 2} ; W^{m, 2}(\Omega)$ is a Hilbert space for the inner product $(\cdot, \cdot)_{m, 2}$,

$$
(u, v)_{m, 2}=\sum_{|\alpha| \leq m}\left(\partial^{\alpha} u . \partial^{\alpha} v\right) \quad \text { for } u, v \in W^{m, 2}(\Omega)
$$

Fractional Sobolev spaces $W^{s, 2}(\Omega)$, with noninteger $s$, are defined by interpolation [19, Chapter VII].
$C^{\nu}(\Omega)$, for $\nu \geq 0, \nu=m+\lambda$ with $0 \leq \lambda<1$, is the space of $m$ times continuously differentiable functions on $\Omega$. whose $m$ th order derivatives satisfy a Hölder condition with exponent $\lambda$ if $\nu$ is not an integer; the norm $\|\cdot\|_{C^{\nu}}$ is defined in the usual way.

The definitions extend to spaces of vector-valued functions in the standard way, with the caveat that the inner product in $\left[L^{2}(\Omega)\right]^{n}$ is defined by $(u, v)=$ $\int_{\Omega} u \cdot v$, where - indicates the scalar product in $\mathrm{R}^{n}$. Complex-valued functions are interpreted as vector-valued functions with two real components.

Functions that vary in space and time, like the order parameter and the vector potential, are considered as mappings from the time domain, which is a subinterval of $\mathbf{R}_{+}$, into spaces of complex- or vector-valued functions defined in $\Omega$. Let $X=\left(X,\|\cdot\|_{X}\right)$ be a Banach space of functions defined in $\Omega$. Then functions of space and time defined on $\Omega \times(0, T)$, for $T>0$, may be considered as elements of $L^{p}(0, T ; X)$, for $1 \leq p \leq \infty$, or $W^{m, 2}(0, T ; X)$, for nonnegative integer $m$, or $C^{\nu}(0, T ; X)$, for $\nu \geq 0, \nu=m+\lambda$ with $0 \leq \lambda<1$. Detailed definitions can be found, for example, in [17].

Obviously, function spaces of ordered pairs $(\dot{w}, A)$, where $\dot{\psi}: \Omega \rightarrow \mathrm{R}^{2}$ and $\boldsymbol{A}: \Omega \rightarrow \mathbf{R}^{n}(n=2,3)$ play an important role in the study of the TDGL equations. We therefore adopt the following special notation: $X=[X(\Omega)]^{2} \times$ $[X(\Omega)]^{n}$ for any Banach space $X(\Omega)$ of real-valued functions defined in $\Omega$. Here, $[X(\Omega)]^{2}$ and $[X(\Omega)]^{n}$ are the underiying Banach spaces for the order parameter $\psi$ and the vector potential $A$, respectively. A suitable framework for the functional analysis of the TDGL equations is the Cartesian product

$$
W^{1+\alpha \cdot 2}=\left[W^{1+\alpha \cdot 2}(\Omega)\right]^{2} \times\left[W^{-1+\alpha .2}(\Omega)\right]^{n} .
$$

This space is densely and continuously imbedded in $W^{1.2} \cap L^{\infty}$ if $\frac{1}{2}<\alpha<1$.

### 3.2 Reduction to Homogeneous Form

We first reduce the boundary conditions in the TDGL equations to homogeneous form. Assume $\boldsymbol{H} \in\left[L^{2}(\Omega)\right]^{n}$. Let $\boldsymbol{A}_{\mathrm{H}}$ be a minimizer of the convex quadratic form $J_{\omega} \equiv J_{\omega}[A]$,

$$
\begin{equation*}
J_{\omega}[\boldsymbol{A}]=\int_{\Omega}\left[\omega(\nabla \cdot \boldsymbol{A})^{2}+|\nabla \times \boldsymbol{A}-\boldsymbol{H}|^{2}\right] \mathrm{d} x \tag{3.1}
\end{equation*}
$$

on the domain

$$
\mathcal{D}\left(J_{\omega}\right)=\mathcal{D}\left(\mathcal{A}_{\mathbf{A}}^{1 / 2}\right)=\left\{\boldsymbol{A} \in\left[W^{1,2}(\Omega)\right]^{n}: \omega(\boldsymbol{n} \cdot \boldsymbol{A})=0 \text { on } \partial \Omega\right\} .
$$

Lemma 3 The functional $J_{\omega}$ has a unique minimizer $\boldsymbol{A}_{\mathrm{H}}$ on $\mathcal{D}\left(J_{\omega}\right)$ if $\omega>0$, and this minimizer has the property $\nabla \cdot \boldsymbol{A}_{\mathrm{H}}=0$ in $\Omega$. The functional $J_{0}$ has a unique minimizer $A_{H}$ on the closed linear subspace

$$
\left\{\boldsymbol{A} \in\left[W^{1,2}(\Omega)\right]^{n}: \nabla \cdot \boldsymbol{A}=0 \text { in } \Omega, \boldsymbol{n} \cdot \boldsymbol{A}=0 \text { on } \partial \Omega\right\}
$$

of $\mathcal{D}\left(J_{0}\right)$.

Proof. Assume $\omega>0$. Then $J_{\omega}[A] \rightarrow \infty$ as $\|A\|_{W^{1,2}} \rightarrow \infty$; see [21. Chapter I, Eq. (5.45)]. Also, $J_{\omega}$ is strictly convex and continuous. Standard methods of the calculus of variations yield the existence of a unique minimizer. This minimizer, $\boldsymbol{A}_{\mathbf{H}}$, is necessarily divergence free. Otherwise, we could replace it by $\boldsymbol{A}_{\mathrm{H}}+\nabla \Phi$ without changing the term $\nabla \times \boldsymbol{A}-\boldsymbol{H}$ and, by taking $\Phi$ as the solution of the Neumann problem for Poisson's equation $\Delta \Phi=-\nabla \cdot \boldsymbol{A}_{\mathrm{H}}$ in $\Omega$, reduce the value of the functional to $J_{\omega}\left[\boldsymbol{A}_{\mathbf{H}}+\nabla \Phi\right]=\int_{\Omega}\left|\nabla \times \boldsymbol{A}_{\mathrm{H}}-\boldsymbol{H}\right|^{2} \mathrm{~d} x$, which is strictly less than $J_{\omega}\left[\boldsymbol{A}_{\mathrm{H}}\right]$. The case $\omega=0$ is similar.

The lemma shows that the property $\nabla \cdot \boldsymbol{A}_{\mathrm{H}}=0$ in $\Omega$ is a consequence of the fact that $\boldsymbol{A}_{\mathbf{H}}$ minimizes the functional $J_{\nu}$ if $\omega>0$. If $\omega=0$. we impose the condition $\boldsymbol{A}_{\mathbf{H}}=0$. In either case, $\boldsymbol{A}_{\mathbf{H}}$ is uniquely determined. and

$$
\begin{equation*}
\nabla \cdot A_{H}=0 \text { in } \Omega . \tag{3.2}
\end{equation*}
$$

We refer to this minimizer $\boldsymbol{A}_{\mathrm{H}}$ as the applied vector potential. It is the solution of the strongly elliptic boundary-value problem

$$
\begin{gather*}
\nabla \times \nabla \times \boldsymbol{A}_{\mathrm{H}}-\omega \nabla\left(\nabla \cdot \boldsymbol{A}_{\mathrm{H}}\right)=\nabla \times H \quad \text { in } \Omega .  \tag{3.3}\\
\omega\left(\boldsymbol{n} \cdot \boldsymbol{A}_{\mathrm{H}}\right)=0 \quad \text { and } \quad \boldsymbol{n} \times\left(\nabla \times \boldsymbol{A}_{\mathrm{H}}-H\right)=0 \quad \text { on } \partial \Omega . \tag{3.4}
\end{gather*}
$$

Lemma 4 If $\boldsymbol{H} \in\left[L^{2}(\Omega)\right]^{n}$, then $\boldsymbol{A}_{\mathbf{H}} \in \mathcal{D}\left(J_{\omega}\right)$. The mapping $\boldsymbol{H} \mapsto \boldsymbol{A}_{H}$ is time independent and continuous from $\left[W^{\theta, 2}(\Omega)\right]^{n}$ to $\left[W^{1+\theta, 2}(\Omega)\right]^{n}$, for $0 \leq$ $\theta \leq 1$.

Proof. The continuity of the mapping $H \mapsto \boldsymbol{A}_{\mathbf{H}}$ follows from the regularity results in Georgescu [20].

We now introduce the reduced vector potential $\boldsymbol{A}^{\prime}$,

$$
\begin{equation*}
A^{\prime}=A-A_{H} \tag{3.5}
\end{equation*}
$$

In terms of $\psi$ and $\boldsymbol{A}^{\prime}$, the TDGL equations assume the form

$$
\begin{gather*}
\frac{\partial v}{\partial t}-\frac{1}{\eta \kappa^{2}} \Delta v=\varphi \quad \text { in } \Omega \times(0, \infty),  \tag{3.6}\\
\frac{\partial \boldsymbol{A}^{\prime}}{\partial t}+\nabla \times \nabla \times \boldsymbol{A}^{\prime}-\omega \nabla\left(\nabla \cdot \boldsymbol{A}^{\prime}\right)=\boldsymbol{F} \quad \text { in } \Omega \times(0, \infty),  \tag{3.7}\\
\boldsymbol{n} \cdot \nabla v+\gamma \psi=0, \quad \omega\left(\boldsymbol{n} \cdot \boldsymbol{A}^{\prime}\right)=0, \quad n \times\left(\nabla \times \boldsymbol{A}^{\prime}\right)=0 \quad \text { on } \partial \Omega \times(0, \infty), \tag{3.8}
\end{gather*}
$$

where $\varphi$ and $\boldsymbol{F}$ are nonlinear functions of $t$ and $u$,

$$
\begin{align*}
& \varphi \equiv \varphi(t, u)=\frac{1}{\eta}\left[-\frac{2 i}{\kappa}(\nabla \psi) \cdot\left(A^{\prime}+A_{\mathrm{H}}\right)-\frac{i}{\kappa}\left(1-\eta \kappa^{2} \omega\right) \psi\left(\nabla \cdot A^{\prime}\right)\right. \\
&\left.-\psi\left|\boldsymbol{A}^{\prime}+\boldsymbol{A}_{\mathrm{H}}\right|^{2}+\left(1-|\psi|^{2}\right) \psi\right]  \tag{3.9}\\
& \boldsymbol{F} \equiv \boldsymbol{F}(t, u)= \frac{1}{2 i \kappa}\left(\psi^{*} \nabla \psi-\psi \nabla \psi^{*}\right)-|\psi|^{2}\left(\boldsymbol{A}^{\prime}+\boldsymbol{A}_{\mathrm{H}}\right)-\frac{\partial \boldsymbol{A}_{\mathrm{H}}}{\partial t} \tag{3.10}
\end{align*}
$$

(The explicit dependence on $t$ is through the applied vector potential $\boldsymbol{A}_{\mathbf{H}}$.) The boundary-value problem (3.6)-(3.8) must be supplemented by appropriate
initial conditions on $\psi$ and $\boldsymbol{A}^{\prime}$. If the initial conditions for the TDGL equations are of the type $\psi=\psi_{0}, \boldsymbol{A}=\boldsymbol{A}_{0}$ on $\Omega \times\{0\}$, then

$$
\begin{equation*}
\psi=\psi_{0} \quad \text { and } \boldsymbol{A}^{\prime}=\boldsymbol{A}_{0}-\boldsymbol{A}_{\mathrm{H}}(0) \quad \text { on } \Omega \times\{0\} \tag{3.11}
\end{equation*}
$$

In the next section we show how the evolution of the solution $\left(\psi, A^{\prime}\right)$ of Eqs. (3.6)-(3.8) is connected with the dynamics of a vector $u$ in the Hilbert space $L^{2}$.

### 3.3 Functional Formulation of the TDGL Equations

Let the vector $u: \mathbf{R}_{+} \rightarrow L^{2}$ represent the pair $\left(\psi, A^{\prime}\right)$,

$$
\begin{equation*}
u=\left(\psi, \boldsymbol{A}^{\prime}\right) \equiv\left(\psi, \boldsymbol{A}-\boldsymbol{A}_{\boldsymbol{H}}\right) \tag{3.12}
\end{equation*}
$$

and let $\mathcal{A}$ be the linear selfadjoint operator in $L^{2}$ associated with the quadratic form $Q_{\omega} \equiv Q_{\omega}[u]$,

$$
\begin{equation*}
Q_{\omega}[u]=\int_{\Omega}\left[\frac{1}{\eta \kappa^{2}}|\nabla \psi|^{2}+\omega\left(\nabla \cdot A^{\prime}\right)^{2}+\left|\nabla \times A^{\prime}\right|^{2}\right] \mathrm{d} x+\int_{\partial \Omega} \frac{\gamma}{\eta \kappa^{2}}|\psi|^{2} \mathrm{~d} \sigma(x) \tag{3.13}
\end{equation*}
$$

with the domain

$$
\mathcal{D}\left(Q_{\mu}\right)=\mathcal{D}\left(\mathcal{A}^{1 / 2}\right)=\left\{u=\left(\psi, \boldsymbol{A}^{\prime}\right) \in W^{1,2}: \omega\left(\boldsymbol{n} \cdot \boldsymbol{A}^{\prime}\right)=0 \text { on } \partial \Omega\right\}
$$

If no confusion is possible, we use the same symbol $\mathcal{A}$ for the restrictions $\mathcal{A}_{\psi}$ and $\mathcal{A}_{\mathrm{A}}$ of $\mathcal{A}$ to the respective linear subspaces $\left[L^{2}(\Omega)\right]^{2} \equiv\left[L^{2}(\Omega)\right]^{2} \times\{0\}$ (for $\psi)$ and $\left[L^{2}(\Omega)\right]^{n} \equiv\{0\} \times\left[L^{2}(\Omega)\right]^{n}($ for $A)$ of $L^{2}$.

Now, consider the initial-value problem

$$
\begin{equation*}
\frac{\mathrm{d} u}{\mathrm{~d} t}+\mathcal{A} u=\mathcal{F}(t, u(t)) \quad \text { for } t>0: \quad u(0)=u_{0} \tag{3.14}
\end{equation*}
$$

in $L^{2}$, where $\mathcal{F}(t, u)=(\varphi, F)$, with $\varphi$ and $F$ given by Eqs. (3.9) and (3.10), and $u_{0}=\left(\psi_{0}, \boldsymbol{A}_{0}-\boldsymbol{A}_{\mathbf{H}}(0)\right)$.

With $\frac{1}{2}<\alpha<1$ and $u_{0} \in W^{1+\alpha, 2}$, we say that $u$ is a mild solution of Eq. (3.14) on the interval $(0, T)$, for some $T>0$, if $u:(0, T) \rightarrow W^{1+\alpha, 2}$ is continuous and

$$
\begin{equation*}
u(t)=e^{-\mathcal{A} t} u_{0}+\int_{0}^{t} e^{-\mathcal{A}(t-s)} \mathcal{F}(s, u(s)) \mathrm{d} s \quad \text { for } 0 \leq t \leq T \tag{3.15}
\end{equation*}
$$

in $L^{2}$. A mild solution of the initial-value problem (3.14) defines a weak solution ( $\boldsymbol{\psi}, \boldsymbol{A}^{\prime}$ ) of the boundary-value problem (3.6)-(3.3). which in turn defines a weak solution ( $w, \boldsymbol{A}$ ) of the TDGL equations, provided $\boldsymbol{A}_{\mathrm{H}}$ is sufficiently regular.

Given any $f=(\varphi, \boldsymbol{F}) \in L^{2}$, the equation $\mathcal{A} u=f$ in $L^{2}$ is equivalent with the system of uncoupled boundary-value problems

$$
\begin{gather*}
-\frac{1}{\eta \kappa^{2}} \Delta w=\because \text { in } \Omega, \quad n \cdot \nabla \because+\gamma=0 \quad \text { on } \partial \Omega:  \tag{3.16}\\
\nabla \times \nabla \times A^{\prime}-w \nabla\left(\nabla \cdot A^{\prime}\right)=F \quad \text { in } \Omega, \quad w\left(n \cdot A^{\prime}\right)=0, n \times\left(\nabla \times A^{\prime}\right)=0 \quad \text { on } \partial \Omega . \tag{3.17}
\end{gather*}
$$

(More precisely, the system of Eqs. (3.16), (3.17) holds in the dual space $\mathcal{D}\left(Q_{\omega}\right)^{\prime}$ of $\mathcal{D}\left(Q_{\omega}\right)$ with respect to the canonically extended inner product $(\cdot, \cdot)$ in $\left[L^{2}(\Omega)\right]^{n}$.) Boundary-value problems of this type have been studied by GEORGESCU [20]. Applying his results, we see that $\mathcal{D}(\mathcal{A})$ is a closed linear subspace of $\boldsymbol{W}^{2,2}$. Since $\mathcal{A}$ is positive definite on $L^{2}$. its fractional powers $\mathcal{A}^{\theta}$ are well defined for all $\theta \in \mathbf{R}$; they are unbounded for $\theta>0$. Interpolation theory shows that $\mathcal{D}\left(\mathcal{A}^{\theta}\right)$ is a closed linear subspace of $W^{2 \theta, 2}$ for $0<\theta<1$.

### 3.4 Smoothing of the Applied Vector Potential

The term $\partial_{t} \boldsymbol{A}_{\mathrm{H}}$ in Eq. (3.10) introduces an integral $\mathcal{J}_{\mathrm{H}}(t)$ in Eq. (3.15),

$$
\begin{equation*}
\mathcal{J}_{\mathrm{H}}(t)=\int_{0}^{t} e^{-\mathcal{A}(t-s)} \frac{\partial \boldsymbol{A}_{\mathrm{H}}}{\partial t}(s) \mathrm{d} s . \tag{3.18}
\end{equation*}
$$

where $\mathcal{J}_{\mathrm{H}}(t) \subseteq\left[L^{2}(\Omega)\right]^{n} \equiv\{0\} \times\left[L^{2}(\Omega)\right]^{n} \subset L^{2}$ for $t \in(0, T)$. The regularity of this integral determines the regularity of the solution $u$ of Eq. (3.14).

Lemma 5 If $H \in W^{1,2}\left(0, T ;\left[L^{2}(\Omega)\right]^{n}\right)$, then $\mathcal{J}_{H}(t) \in \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$ for $0 \leq$ $\alpha<1$, for every $t \in(0, T)$, and $\mathcal{J}_{\mathrm{H}} \in C^{\beta}\left(0, T ;\left[W^{1+\alpha, 2}(\Omega)\right]^{n}\right)$ for $0 \leq \beta<$ $\frac{1}{2}(1-\alpha)$.

Proof. Assume that $0 \leq \alpha<1$ and $0 \leq \beta<\frac{1}{2}(1-\alpha)$. The proof of the lemma uses the inequalities

$$
\begin{gather*}
\left\|\mathcal{A}^{\alpha / 2} e^{-\mathcal{A} s}\right\|_{L^{2} \leq C s^{-\alpha / 2} \quad \text { for } 0<s \leq T}\left\|\left(e^{-\mathcal{A} s}-I\right) \mathcal{A}^{-\beta}\right\|_{L^{2} \leq C s^{\beta}} \quad \text { for } 0 \leq s \leq T \tag{3.19}
\end{gather*}
$$

where the positive constants $C$ do not depend on $s$ [17, Theorem 1.4.3].
Because $\partial_{t} H \in L^{2}\left(0, T ;\left[L^{2}(\Omega)\right]^{n}\right)$, it follows immediately from Lemma $t$ that $\partial_{t} \boldsymbol{A}_{\mathrm{H}} \in L^{2}\left(0, T ;\left[W^{1,2}(\Omega)\right]^{n}\right)$. Standard arguments then lead to the con-
 $\mathcal{A}^{1 / 2} \partial_{t} A_{\mathrm{H}} \in L^{2}\left(0, T ;\left[L^{2}(\Omega)\right]^{n}\right)$ and

$$
\begin{equation*}
\mathcal{A}^{(1+\alpha) / 2} \mathcal{J}_{\mathrm{H}}(t)=\int_{0}^{t} \mathcal{A}^{\alpha / 2} e^{-\mathcal{A}(t-s)} \mathcal{A}^{1 / 2} \frac{\partial \mathcal{A}_{\mathrm{H}}}{\partial t}(s \mid \mathrm{d} s \quad \text { for } 0 \leq t \leq T \tag{3.21}
\end{equation*}
$$

in $\left[L^{2}(\Omega)\right]^{n}$. Applying the estimate (3.19). we obtain

$$
\begin{gather*}
\left\|\mathcal{A}^{(1+\alpha) / 2} \mathcal{J}_{\mathrm{H}}(t)\right\|_{L^{2}} \leq \int_{0}^{t}\left\|\mathcal{A}^{\alpha / 2} e^{-\mathcal{A} s}\right\|_{L^{2}}\left\|_{\mathcal{A}^{1 / 2}} \frac{\partial A_{\mathrm{H}}}{\partial t}(t-s)\right\|_{L^{2}} \mathrm{~d} s \\
\leq C \int_{0}^{t}\left\|\mathcal{A}^{1 / 2} \frac{\partial \mathcal{A}_{\mathrm{H}}}{\partial t}(t-s)\right\|_{L^{2}} s^{-\alpha / 2} \mathrm{~d} s \\
\leq C\left(\int_{0}^{t}\left\|\mathcal{A}^{1 / 2} \frac{\partial A_{\mathrm{H}}}{\partial t}(t-s)\right\|_{L^{2}}^{2} \mathrm{~d} s\right)^{1 / 2}\left(\int_{0}^{t} s^{-\alpha} \mathrm{d} s\right)^{1 / 2} \\
=\frac{C}{(1-\alpha)^{1 / 2}} t^{(1-\alpha) / 2}\left(\int_{0}^{t}\left\|\mathcal{A}^{1 / 2} \frac{\partial \mathcal{A}_{\mathrm{H}}}{\partial t}(s)\right\|_{L^{2}}^{2} \mathrm{~d} s\right)^{1 / 2}, \tag{3.22}
\end{gather*}
$$

so $\mathcal{J}_{\mathrm{H}}(t) \in \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$, a closed subspace of $\left[W^{1+\alpha, 2}(\Omega)\right]^{n}$, for every $t \in(0, T)$.
To prove the Holder continuity of $\mathcal{J}_{\mathrm{H}}$, we take $0<t<t^{\prime}<T$ and use the following identity in $\left[L^{2}(\Omega)\right]^{n}$, which follows immediately from the definition (3.18).

$$
\begin{gather*}
\mathcal{A}^{(1+\alpha) / 2}\left(\mathcal{J}_{\mathrm{H}}\left(t^{\prime}\right)-\mathcal{J}_{\mathrm{H}}(t)\right) \\
=\mathcal{A}^{\alpha / 2}\left[\int_{0}^{t^{\prime}} e^{-\mathcal{A}\left(t^{\prime}-s\right)} \mathcal{A}^{1 / 2} \frac{\partial \boldsymbol{A}_{\mathrm{H}}}{\partial t}(s) \mathrm{d} s-\int_{0}^{t} e^{-\mathcal{A}(t-s)} \mathcal{A}^{1 / 2} \frac{\partial \boldsymbol{A}_{\mathrm{H}}}{\partial t}(s) \mathrm{d} s\right] \\
=\mathcal{J}_{1}\left(t, t^{\prime}\right)+\mathcal{J}_{2}\left(t, t^{\prime}\right) . \tag{3.23}
\end{gather*}
$$

where

$$
\begin{gathered}
\mathcal{J}_{1}\left(t, t^{\prime}\right)=\int_{0}^{t^{\prime}-t} \mathcal{A}^{\alpha / 2} e^{-\mathcal{A} s} \mathcal{A}^{1 / 2} \frac{\partial \mathcal{A}_{\mathrm{H}}}{\partial t}\left(t^{\prime}-s\right) \mathrm{d} s . \\
\mathcal{J}_{2}\left(t, t^{\prime}\right)=\left(e^{-\mathcal{A}\left(t^{\prime}-t\right)}-I\right) \int_{0}^{t} \mathcal{A}^{\alpha / 2} e^{-\mathcal{A}(t-s)} \mathcal{A}^{1 / 2} \frac{\partial \boldsymbol{A}_{\mathrm{H}}}{\partial t}(s) \mathrm{d} s .
\end{gathered}
$$

We estimate the $\left[L^{2}(\Omega)\right]^{n}$-norms of $\mathcal{J}_{1}\left(t, t^{\prime}\right)$ and $\mathcal{J}_{2}\left(t, t^{\prime}\right)$ as in (3.22), making use of the inequalities (3.19) and (3.20),

$$
\begin{equation*}
\left\|\mathcal{J}_{1}\left(t, t^{\prime}\right)\right\|_{L^{2}} \leq \frac{C}{(1-\alpha)^{1 / 2}}\left|t^{\prime}-t\right|^{(1-\alpha) / 2}\left(\int_{t}^{t^{\prime}}\left\|\mathcal{A}^{1 / 2} \frac{\partial A_{\mathrm{H}}}{\partial t}(s)\right\|_{L^{2}}^{2} \mathrm{~d} s\right)^{1 / 2} \tag{3.24}
\end{equation*}
$$

$$
\begin{aligned}
& \left\|\mathcal{J}_{2}\left(t, t^{\prime}\right)\right\|_{L^{2}}=\left\|\left(e^{-\mathcal{A}\left(t^{\prime}-t\right)}-I\right) \mathcal{A}^{-3} \int_{0}^{t} \mathcal{A}^{3+(\alpha / 2)} e^{-\mathcal{A}(t-s)} \mathcal{A}^{1 / 2} \frac{\partial \mathcal{A}_{H}}{\partial t}(s) \mathrm{d} s\right\|_{L^{2}}
\end{aligned}
$$

Here, the positive constants $C$ depend only on $\mathcal{A} . \alpha$. and 3. The statement of the lemma follows.

## 4 Results

In this section we present our results in three theorems, each with a corollary. The proofs are deferred until Section 5. Unless indicated otherwise, we assume that the data entering the equations satisfy the following hypotheses:
(H1) $\Omega \subset \mathbf{R}^{n}(n=2$ or 3$)$ is bounded, with $\partial \Omega$ of class $C^{1,1}$. (That is, $\partial \Omega$ is a compact ( $n-1$ )-manifold described by Lipschitz-continuously differentiable charts.)
(H2) $\gamma: \partial \Omega \rightarrow \mathbf{R}$ is Lipschitz continuous, with $\gamma(x) \geq 0$ for all $x \in \partial \Omega$.
(H3) $\omega . T, \alpha .3 \in \mathbf{R}$ are constants. such that $0 \leq \omega<\infty, 0 \leq T<\infty$, $\frac{1}{2}<\alpha<1$, and $0 \leq 3<\frac{1}{2}(1-\alpha)$.
(H4) $\boldsymbol{H} \in L^{\infty}\left(0 . T ;\left[W^{\alpha, 2}(\Omega)\right]^{n}\right) \cap W^{1,2}\left(0 . T ;\left[L^{2}(\Omega)\right]^{n}\right)$.

The applied vector potential $\boldsymbol{A}_{\mathrm{H}}$ is defined by Eqs. (3.3), (3.4).

### 4.1 Existence and Uniqueness

Our first theorem gives the existence and uniqueness of a mild solution of the initial-value problem (3.14).

Theorem 1 Let the initial data $\left(\psi_{0}, \boldsymbol{A}_{0}\right)$ be such that $u_{0}=\left(\psi_{0}, \boldsymbol{A}_{0}^{\prime}\right) \equiv\left(\psi_{0}, \boldsymbol{A}_{0}\right.$ $\left.-A_{H}(0)\right)$ is in $\mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$. Then the initial-value problem (3.14) has a unique
mild solution $u=\left(\psi, \boldsymbol{A}^{\prime}\right) \equiv\left(\psi, \boldsymbol{A}-\boldsymbol{A}_{H}\right)$. such that $u \in C\left(0 . T: W^{1+\alpha \cdot 2}\right)$. The order parameter $w$ of this solution satisfies the "maximum modulus principle,"

$$
\begin{equation*}
|\psi(x, t)| \leq \max \left\{1,\left\|\psi_{0}\right\|_{L^{\infty}(\Omega)}\right\} \quad \text { for all }(x, t) \in \bar{\Omega} \times(0 . T) \tag{4.1}
\end{equation*}
$$

Also, $(\dot{\psi}, \boldsymbol{A}) \in W^{1,2}\left(0, T ; L^{2}\right)$ and $\nabla \cdot \boldsymbol{A} \in L^{2}\left(0, T ;\left[W^{-1.2}(\Omega)\right]^{n}\right)$.

The proof of Theorem 1 is given in Section 5.1.
Observe that the theorem states that $\left.\left(\boldsymbol{\psi} \cdot \mathcal{A}^{\prime}\right) \subseteq C i 0 . T: W^{-1+x \cdot 2}\right)$. To obtain a comparable result for ( $\boldsymbol{\psi}, \boldsymbol{A}$ ), we need the continuity $\boldsymbol{A}_{H}$ in time. which. by Lemma 4 , is controlled by the continuity of $H$ in time. In the hyporhesis (H4), we have imposed a minimum condition on $H$. If the hypothesis (H4) is strengthened to $\boldsymbol{H} \in C\left(0, T ;\left[W^{\alpha, 2}(\Omega)\right]^{n}\right)$, then $(\psi, \boldsymbol{A}) \in C\left(0, T ; W^{1+\alpha, 2}\right)$.

Theorem 1 defines a solution map $S_{0}: \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right) \rightarrow C\left(0, T ; W^{1+\alpha, 2}\right)$,

$$
\begin{equation*}
u(t)=S_{0}(t) u_{0}, \quad u_{0} \in \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right), t \in(0, T) \tag{4.2}
\end{equation*}
$$

The properties of $S_{0}$ are considered in more detail in the following section.
Theorem 1 implies the existence and uniqueness of a weak solution of the TDGL equations.

Corollary 1 The pair $(\psi, \boldsymbol{A})$ obtained in Theorem 1 is a weak solution of the boundary-value problem (3.6)-(3.8): Eqs. (3.6) and (3.7) are satisfied in the $L^{2}(\Omega \times(0 . T))$-sense, Eq. (3.8) in the sense of traces in $L^{\infty}\left(0 . T: W^{\alpha-1 / 2.2}(\partial \Omega)\right)$.

### 4.2 Regularity

The following theorem improves the continuous dependence of the solution $u$ on the initial data $u_{0}$. Let the map $S_{\mathcal{B}^{\prime}}: \mathcal{D}\left(\mathcal{A}^{\left(1+\alpha^{\prime}\right) / 2}\right) \rightarrow C^{\mathcal{B}}\left(0, T ; W^{1+\alpha}\right)$ be defined by the identity

$$
\begin{equation*}
t^{\beta^{\prime}} u(t)=S_{\beta^{\prime}}(t) u_{0}, \quad u_{0} \in \mathcal{D}\left(\mathcal{A}^{\left(1+\alpha^{\prime}\right) / 2}\right), t \in(0, T) \tag{4.3}
\end{equation*}
$$

for suitable exponents $\alpha, \alpha^{\prime}, \beta$, and $\beta^{\prime}$.

Theorem 2 Assume that $\frac{1}{2}<\alpha^{\prime} \leq \alpha<1,0 \leq \beta<\frac{1}{2}(1-\alpha)$, and $\beta^{\prime}=$ $\beta+\frac{1}{2}\left(\alpha-\alpha^{\prime}\right)$. Then the mapping $S_{\beta^{\prime}}$ defined in $E q$. (4.3) is uniformly Lipschitz continuous on bounded subsets of $\mathcal{D}\left(\mathcal{A}^{\left(1+\alpha^{\prime}\right) / 2}\right)$.

The proof of Theorem 2 is given in Section 5.2.
Theorem 2 implies the existence of a dynamical process for the TDGL equations, even when the applied magnetic field is time dependent.

Corollary 2 The mild solutions $u(t)(0<t<T)$ of Eq. (3.14) obtained in Theorem 1 generate a dynamical process $\{U(t, s): 0<s \leq t<T\}$ on $\mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$. so $u(t)=L^{-}(t, s) u(s)$ for all $0<s \leq t<T$. Woreover. for $0<s<t<T$. each $C^{-}(t . s): \mathcal{D}\left(\mathcal{A}^{(1+x) / 2}\right)-\mathcal{D}\left(\mathcal{A}^{(1+x) / 2}\right)$ maps bounded sets into relatively compact sets.

Observe that the theorem states a regularity result for ( $\psi, \boldsymbol{A}^{\prime}$ ). To obtain a comparable result for $(\psi, \boldsymbol{A})$, we need sufficient regularity of $\boldsymbol{A}_{\mathrm{H}}$. The regularity of $\boldsymbol{A}_{\mathbf{H}}$ is, by Lemma 4, controlled by the regularity of $H$. In the hypothesis (H4), we have imposed minimum regularity on $H$. If the hypothesis (H4) is strengthened to $\boldsymbol{H} \in C^{\beta}\left(0, T ;\left[W^{\alpha, 2}(\Omega)\right]^{n}\right)$, then $(\psi, \boldsymbol{A}) \in C^{\beta}\left(0, T ; W^{1+\alpha, 2}\right)$.

### 4.3 Large-Time Asymptotic Behavior

Next, we investigate the asymptotic behavior of the mild solution $u(t) \in$ $\mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$ of Eq. (3.14) as $t \rightarrow \infty$. We restrict ourselves to the case of a stationary applied magnetic field $H$.

If $\partial_{t} H=0$. the hypothesis (H4) reduces to $H \in\left[W^{\alpha, 2}(\Omega)\right]^{n}$, the quantity $P$ defined in Eq. (2.22) is zero, and the inequality (2.23) simplifies to

$$
\begin{gather*}
E_{\mu}(t)+2 \int_{0}^{t} \int_{\Omega}\left[\eta\left|\frac{\partial v}{\partial t}-i \kappa \dot{\omega} w(\nabla \cdot \boldsymbol{A})\right|^{2}+\left|\frac{\partial \boldsymbol{A}}{\partial t}\right|^{2}+\omega^{2}|\nabla(\nabla \cdot A)|^{2}\right] \mathrm{d} x \mathrm{~d} t^{\prime} \\
\leq E_{\nu}(0) . \quad t \in(0, T) \tag{4.4}
\end{gather*}
$$

The dynamical process $\{U(t, s): 0<s \leq t<T\}$ on $\mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$ introduced in Corollary 2 is defined for every $T>0$ (see Lemma 1) and becomes a dynamical system,

$$
\begin{equation*}
S=\{S(t): t \geq 0\} \quad \text { on } \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right) \tag{4.5}
\end{equation*}
$$

where

$$
\begin{equation*}
S(t-s)=U(t, s), \quad t \geq s \geq 0 \tag{4.6}
\end{equation*}
$$

The set $\left\{S(t) u_{0}: t \geq 0\right\}$ is called the (forward) orbit of $u_{0} \in \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$ under $S$. We denote the set of all limit points (as $t \rightarrow \infty$ ) of the orbit of $u_{0}$ by $\omega\left(u_{0}\right)$ and call it the omega-limit set of $u_{0}$.

The following theorem shows that the functional $E_{i}$ is a Liapunov functional for the dynamical system $S$ in the following sense (cf. [22, Chapter VII, Definition 4.1]): (i) $E_{\omega}: \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right) \rightarrow \mathbf{R}$ is continuous, and (ii) if $u_{0} \in \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$ is such that $E_{\omega}\left[S(t) u_{0}\right]=E_{\omega}\left[u_{0}\right]$ for some $t>0$, then $u_{n}$ is a stationary point for $\lrcorner$.

Theorem 3 The dynamical system $S$ defined in Eq. (4.5) has the following properties:
(i) $E_{i}$ is a Liapunou functional for $S$.
(ii) The orbit of each $u_{0} \in \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$ has compact closure in $W^{1+\alpha, 2}$.
(iii) The omega-limit set of each $u_{0} \in \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$ is a nonempty compact connected set of divergence-free equilibria.

The proof of Theorem 3 is given in Section 5.3.
Property (iii) of Theorem 3 says, in effect, that every element of any omegalimit set is a solution of the stationary GL equations (2.29)-(2.31) in the London gauge, $\nabla \cdot \boldsymbol{A}=0$ in $\Omega$, that satisfies the condition $\boldsymbol{n} \cdot \boldsymbol{A}=0$ on the boundary $\partial \Omega$. It will be seen in the proof of Theorem 3 (iii) that, for the degenerate case $\omega=0$ (hence, $\phi=0$ ), we cannot conclude that the equilibrium solutions are divergence free. For the TDGL equations, the gauge " $\phi=-\omega(\nabla$. $\boldsymbol{A})^{n}$ and the condition $\nabla \cdot \boldsymbol{A}=0$ together imply the identity $\phi=0$.

An attractor of the dynamical system $S$ is the omega-limit set of one of its open neighbors. An attractor is called a global attractor if it attracts all its open bounded neighbors.

An immediate consequence of Theorem 3 is given in the following corollary; see [22, Chapter VII, Theorem 4.1].

Corollary 3 Let $\mathcal{A}$ be the global attractor and $\mathcal{E}$ the set of all stationary points of $S$. If $\mathcal{E}$ is discrete, then $\mathcal{A}$ is the union of $\mathcal{E}$ and the heteroclinic orbits between points of $\mathcal{E}$.

## 5 Proofs

In this section we give the proofs of the theorems presented in the previous section. We begin by recalling some general properties of the fractional powers
of the operator $\mathcal{A}$ defined in Eq. (3.13): cf. [17].
The fractional powers $\mathcal{A}^{\theta}$ of the second-order elliptic differential operator $\mathcal{A}$ defined in Eqs. (3.16) and (3.17) are well defined for all real $\theta$. They are unbounded for $\theta>0$. The dumain $\bar{L}\left(\mathcal{A l}^{i} ;\right.$ is a ciosed linear subspace oi $W^{2 \theta, 2}$ for $0<\theta<1$; hence, $C^{\mathcal{\beta}}\left(0, T ; \mathcal{D}\left(\mathcal{A}^{\theta}\right)\right)$ is a closed linear subspace of $C^{\theta}\left(0 . T: W^{2 \theta .2}\right)$ for this range of values of $\theta$. Furthermore. for $\frac{3}{2}<\theta \leq 2$ (and $n=2$ or 3), the traces of $\nabla v$. A. and $\nabla<A$ belong to the spaces $\left[W^{-\theta-3 / 2.2}(\partial \Omega)\right]^{2 n} \cdot\left[W^{-\theta-1 / 2.2}(\partial \Omega)\right]^{n}$. and $\left[W^{-\theta-3 / 2.2}(\partial \Omega)\right]^{n}$. respectively, and satisfy the boundary conditions specified in Eqs. (3.16) and (3.17). Similarly. the applied vector potential $\boldsymbol{A}_{\mathrm{H}}$ and its curl $\nabla \times \boldsymbol{A}_{\mathrm{H}}$ satisfy the boundary conditions (3.4) if $H \in\left[W^{\theta-1,2}(\Omega)\right]^{n}$.

### 5.1 Proof of Theorem 1

Proof. (i) Local existence and uniqueness. The proof is based on the contraction mapping principle applied to Eq. (3.15) in the space $C\left(0, T ; W^{1+\alpha, 2}\right)$ for $T$ sufficiently small positive. The choice of the target space $W^{1+\alpha, 2}$ is justified because $W^{1+\alpha, 2}$ is continuously imbedded in $W^{1,2} \cap L^{\infty}$ for $\frac{1}{2}<\alpha<1$.

It suffices to prove that $\mathcal{F}(s, \cdot)$ is locally Lipschitz for each $s \in(0, T)$, where $T$ may depend on the Lipschitz constant. Each term in $\mathcal{F}$ is estimated separately. For example, for any two elements $u_{1}=\left(\dot{w}_{1}, \boldsymbol{A}_{1}^{\prime}\right)$ and $u_{2}=\left(\psi_{2}, \boldsymbol{A}_{2}^{\prime}\right)$ of $W^{1+\alpha, 2}$, we have

$$
\begin{gathered}
\left\|w_{1}^{*} \nabla \dot{\psi}_{1}-\psi_{2}^{*} \nabla \dot{\psi}_{2}\right\|_{L^{2}} \leq\left\|w_{1}\right\|_{L^{\infty}} \mid \psi_{1}-\dot{\psi}_{2}\left\|_{W^{1,2}}+\right\| \psi_{2}\left\|_{W^{1,2}}\right\| \psi_{1}-\dot{\psi}_{2} \|_{L^{\infty}} \\
\leq C\left\|u_{1}-u_{2}\right\|_{W^{1+\alpha, 2}}
\end{gathered}
$$

where $C$ is a positive constant that depends only on the norms of $u_{1}$ and $u_{1}$ in $W^{1+\alpha, 2}$. Similar estimates hold for the other terms in $\mathcal{F}$.

Let $B_{R}$ be the ball of radius $R$ centered at the origin in $W^{1+\alpha \cdot 2}$. Then, for any pair $u_{1}, u_{2} \in B_{R}$,

$$
\begin{equation*}
\left\|\mathcal{F}\left(s, u_{1}\right)-\mathcal{F}\left(s, u_{2}\right)\right\|_{L^{2}} \leq C\left\|u_{1}-u_{2}\right\|_{W^{1+\alpha, 2}}, \quad s \in(0, T) \tag{5.1}
\end{equation*}
$$

where the Lipschitz constant $C$ depends on $R$, but not on $s$. The remainder of the proof is standard; see [17, Theorem 3.3.3].
(ii) Global existence. The maximum modulus principle (4.1) is a consequence of the maximum principle applied to Eq. (2.5). (Note that every constant $M$ with $M \geq 1$ is a supersolution of Eq. (2.5).)

The functional $E_{\omega}(t)$ defined in Eq. (2.20) is bounded on every interval ( $0, T$ ), by Lemma l, so

$$
\psi \in L^{\infty}\left(0, T ;\left[W^{1,2}(\Omega)\right]^{2}\right) \quad \text { and } \quad A \in L^{\infty}\left(0 . T:\left[W^{-1,2}(\Omega)\right]^{n}\right):
$$

cf. [21, Chapter I. Eq. (5.45)]. Also, $\boldsymbol{A}_{\mathrm{H}} \in L^{\infty}\left(0, T:\left[W^{1,2}(\Omega)\right]^{n}\right)$, because of the hypothesis (H4). Hence, $u=\left(i, A^{\prime}\right) \in L^{\infty}\left(0 . T: W^{1,2}\right)$.

It follows from Lemma 2 . inequality $(-27)$, that $(G) \equiv W^{-1.2}\left(0 . T: L^{2}\right)$ and $\nabla \cdot A \in L^{2}\left(0 . T:\left[W^{1.2}(\Omega)\right]^{n}\right)$. We also have $A_{H} \in W^{-1.2}\left(0 . T:\left[L^{2}(\Omega)\right]^{n}\right.$. again because of the hypothesis (H4). Therefore. $u \equiv W^{-1.2}\left(0 . T: L^{2}\right)$.

We improve this regularity result by taking advantage of the smoothing action of the semigroup $e^{-\mathcal{A t}}$. This smoothing action has already been demonstrated on the term $\partial_{t} \boldsymbol{A}_{\mathrm{H}}$ in Section 3.3. We first treat $\boldsymbol{A}^{\prime}$ and use the result to improve the regularity of $\psi$. Each term in $\boldsymbol{F}$ needs to be estimated separately. For example,

$$
\left\|\psi^{*} \nabla \dot{\psi}\right\|_{L^{2}} \leq\|\psi\|_{L^{\infty}}\|\psi\|_{W^{1,2}} \leq C\|u\|_{W^{1,2}} .
$$

Here, $C=\max \left\{1,\left\|\psi_{0}\right\|_{L^{\infty}}\right\}$, which is independent of $\psi$. Similar estimates hold for the other terms in $F$, so $F \in L^{\infty}\left(0 . T ;\left[L^{2}(\Omega)\right]^{n}\right)$. Therefore,

$$
\left(t \mapsto \int_{0}^{t} e^{-A(t-s)} F(s) \mathrm{d} s\right) \in C\left(0, T ;\left[W^{1+\alpha, 2}(\Omega)\right]^{n}\right)
$$

so $A^{\prime} \in C\left(0, T ;\left[W^{1+\alpha, 2}(\Omega)\right]^{n}\right)$.
Next, we improve the regularity of $w$. Again, each term in $\varphi$ needs to be estimated separately. For example,

$$
\left\|(\nabla \dot{\psi}) \cdot\left(A_{H}+A^{\prime}\right)\right\|_{L^{2}} \leq\left\|(\nabla \psi) \cdot \boldsymbol{A}_{\mathrm{H}}\right\|_{L^{2}}+\left\|(\nabla \psi) \cdot \boldsymbol{A}^{\prime}\right\|_{L^{2}},
$$

where

$$
\left\|(\nabla \psi) \cdot \boldsymbol{A}_{\mathrm{H}}\right\|_{L^{2}} \leq\|\nabla \psi\|_{L^{2}}\left\|\boldsymbol{A}_{\mathbf{H}}\right\|_{L^{\infty}} \leq C\|u\|_{W^{1,2}}\left\|\boldsymbol{A}_{\mathrm{H}}\right\|_{W^{1+\infty, 2}}
$$

and

$$
\left\|(\nabla \psi) \cdot \boldsymbol{A}^{\prime}\right\|_{L^{2}} \leq\|\nabla \psi\|_{L^{2}}\left\|\boldsymbol{A}^{\prime}\right\|_{L^{\infty}} \leq C\|u\|_{W^{1,2}}\left\|\boldsymbol{A}^{\prime}\right\|_{W^{1+\alpha, 2}}
$$

(To obtain the last estimate, we used the Sobolev imbedding theorem.) Similar estimates hold for the other terms in $\varphi$, so $\varphi \in L^{\infty}\left(0, T ;\left[L^{2}(\Omega)\right]^{2}\right)$ and, therefore, $\psi \in C\left(0, T ;\left[W^{1+\alpha, 2}(\Omega)\right]^{2}\right)$. It follows that $u \in C\left(0, T ; W^{1+\alpha, 2}\right)$, as claimed.

### 5.2 Proof of Theorem 2

Proof. We use Eq. (3.15) to prove the regularity of the solution $u$ of the initial-value problem (3.14).

Let $B_{R}$ be the ball of radius $R$ centered at the origin in $W^{1+\alpha .2}$. Let $u_{1}$ and $u_{2}$ satisfy Eq. (3.15) with initial data $u_{10}$ and $u_{20}$. respectively, in $B_{R}$. Define $v=u_{1}-u_{2}$ and $c_{0}=u_{10}-u_{20}$. Combining the inequality (.5.1) with Eq. (3.1.5), we obtain

$$
\begin{gather*}
\|v(t)\|_{W^{i+\alpha, 2}} \leq\left\|e^{-\mathcal{A} t}\right\|_{W^{i+\alpha, 2}\left\|v_{0}\right\|_{W^{1+\alpha 2}}}^{+C \int_{0}^{t}\left\|\mathcal{A}^{(1+\alpha) / 2} e^{-\mathcal{A}(t-3)}\right\|_{W^{1+\alpha, 2} \| v(s)} \|_{W^{1+\alpha, 2} \mathrm{~d} s}}
\end{gather*}
$$

Applying Gronwall's inequality, we find

$$
\begin{equation*}
\|v(t)\|_{W^{1+\alpha, 2}} \leq C\left\|v_{0}\right\|_{W^{1+\alpha, 2}}, \quad 0<t<T \tag{5.3}
\end{equation*}
$$

so the mapping $S_{0}$ defined in Eq. (4.2) is Lipschitz continuous on $B_{R}$.
Set $f(s)=\mathcal{F}\left(s, u_{1}(s)\right)-\mathcal{F}\left(s, u_{2}(s)\right)$. Then, for $0<t<t^{\prime}<T$,

$$
\begin{aligned}
v\left(t^{\prime}\right)-v(t) & =\left(e^{-\mathcal{A}\left(t^{\prime}-t\right)}-I\right) e^{-\mathcal{A} t} v_{0}+\int_{0}^{t^{\prime}-t} e^{-\mathcal{A} s} f\left(t^{\prime}-s\right) \mathrm{d} s \\
& +\left(e^{-\mathcal{A}\left(t^{\prime}-t\right)}-I\right) \int_{0}^{t} e^{-\mathcal{A}(t-s)} f(s) \mathrm{d} s
\end{aligned}
$$

Taking $\alpha, \alpha^{\prime}, 3$, and $\beta^{\prime}$ subject to the conditions of the theorem, we obtain

$$
\begin{gathered}
\mathcal{A}^{(1+\alpha) / 2}\left(v\left(t^{\prime}\right)-v(t)\right) \\
=\left(e^{-\mathcal{A}\left(t^{\prime}-t\right)}-I\right) \mathcal{A}^{-\mathcal{B}} \mathcal{A}^{\mathcal{\beta}^{\prime}} e^{-\mathcal{A} t} \mathcal{A}^{\left(1+\alpha^{\prime}\right) / 2} v_{0}+\int_{0}^{t^{\prime}-t} \mathcal{A}^{(1+\alpha) / 2} e^{-\mathcal{A} s} f\left(t^{\prime}-s\right) \mathrm{d} s \\
+\left(e^{-\mathcal{A}\left(t^{\prime}-t\right)}-I\right) \mathcal{A}^{-\mathcal{B}} \int_{0}^{t} \mathcal{A}^{\beta+(1+\alpha) / 2} e^{-\mathcal{A}(t-s)} f(s) \mathrm{d} s
\end{gathered}
$$

Using the inequalities (3.19) and (3.20), we deduce the estimates

$$
\begin{gathered}
\left\|\mathcal{A}^{(1+\alpha) / 2}\left(v\left(t^{\prime}\right)-v(t)\right)\right\|_{L^{2}} \leq C_{1}\left(t^{\prime}-t\right)^{\beta} t^{-\beta^{\prime}}\left\|\mathcal{A}^{\left(1+\alpha^{\prime}\right) / 2} v_{0}\right\|_{L^{2}} \\
+C_{2}\left(\left(t^{\prime}-t\right)^{(1-\alpha) / 2}+\left(t^{\prime}-t\right)^{\beta} t^{(1-\alpha) / 2-\beta}\right) \text { ess } \sup \left\{\|f(s)\|_{L^{2}}: 0<s<T\right\} \\
\leq C\left(t^{\prime}-t\right)^{\beta} t^{-\beta^{\prime}}\left(\left\|v_{0}\right\|_{W^{1+\alpha^{\prime}, 2}}+C \sup \left\{\|v(s)\|_{W^{1,2}}: 0<s<T\right\}\right)
\end{gathered}
$$

But, as we have seen, the solution map $S_{0}$ defined in (4.2) is Lipschitz continuous, so $\sup \left\{\|v(s)\| W_{W^{1,2}}: 0<s<T\right\} \leq C\left\|v_{0}\right\|_{W^{1,2}}$. Therefore, the mapping (4.3) is Lipschitz continous, as claimed.

### 5.3 Proof of Theorem 3

Proof. (i) The continuity of the functional $E_{\omega}$ follows from the continuous imbedding of $W^{1+\alpha, 2}$ inton $W^{1,2} \cap L^{\infty}$.

Let $u_{0}=\left(\dot{\psi}, A-A_{H}\right) \in \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$ be such that $E_{\omega}\left[S(t) u_{0}\right]=E_{\omega}\left[u_{0}\right]$ for some $t>0$. From the inequality (4.4), we obtain immediately the identities $\partial_{t} \boldsymbol{A}=0$ and $\omega \nabla(\nabla \cdot \boldsymbol{A})=0$ in $\Omega \times(0, t)$. The first identity implies that $\partial_{t}(\nabla \cdot A)=0$ in $\Omega \times(0, t)$. From this and the second identity we deduce that $\omega \nabla \cdot A=c$ in $\Omega \times(0 . t)$, where $c$ is a real constant. We conclude from $E q \cdot 1 \cdot .1 T$ that $J_{s}=0$. Also, the inequality (t.t) implies $\partial_{t} w=i \kappa c w$ in $\left[L^{2}(\Omega \times(0 . t))\right]^{2}$. so Eq. (2.18) reduces to $c|w|^{2}=0$ in $\Omega \times(0, t)$. We claim that $c=0$.

Suppose $c \neq 0$. Then it must be the case that $\dot{w}=0$ in $\Omega \times(0, t)$. Equations (2.15)-(2.16) reduce to the boundary-value problem (3.3)-(3.4) for $\boldsymbol{A}_{\mathbf{H}}$. Therefore, $\boldsymbol{A}=\boldsymbol{A}_{\mathrm{H}}$ and $\boldsymbol{A}^{\prime}=0$ in $\Omega \times(0, t)$, so $c=\omega \nabla \cdot \boldsymbol{A}_{\mathrm{H}}=0$, and we have a contradiction.

The identity $\partial_{t} \psi=0$ in $\Omega \times(0, t)$, together with the identity $\partial_{t}(\nabla \cdot \boldsymbol{A})=0$ established above, implies that $S\left(t^{\prime}\right) u_{0}=u_{0}$ for all $t^{\prime} \in(0, t)$.

We have proved that $E_{\omega}$ is a Liapunov functional for $S$.
(ii) An immediate consequence of Corollary 2.
(iii) It follows from (ii) that the omega-limit set of each $u_{0} \in \mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$ is nonempty and compact. We prove by contradiction that $\omega\left(u_{0}\right)$ is connected. Suppose $\omega\left(u_{0}\right)$ is not connected. Then $\omega\left(u_{0}\right)=K_{1} \cup K_{2}$, where $K_{1}$ and $K_{2}$ are compact and disjoint. Hence, there exist two disjoint open neighborhoods $N_{1}$ and $N_{2}$ of $K_{1}$ and $K_{2}$, respectively, in $\mathcal{D}\left(\mathcal{A}^{(1+\alpha) / 2}\right)$ and $t_{0} \geq 0$, such that $S(t) u_{0} \in N_{1} \cup N_{2}$ for all $t \geq t_{0}$. But $\left\{S(t) u_{0}: t \geq t_{0}\right\}$, being the image of the interval $\left[t_{0}, \infty\right)$, is connected, so we have a contradiction.

The proof that the omega-limit set of $u_{0}$ consists of equilibrium points only is standard; cf. [22, Chapter VII, Proof of Theorem 4.1].

If $w=\left(\psi, \boldsymbol{A}-\boldsymbol{A}_{\mathbf{H}}\right) \in \omega\left(u_{0}\right)$, then $E_{\omega}[S(t) w]=E_{\omega}[w]$ for all $t>0$, and the same argument as in (i) above leads to the conclusion that $\omega(\nabla \cdot \boldsymbol{A})=0$ in $\Omega$.

## DISCLAIMER

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