Nonlinear stability of source defects in the complex Ginzburg–Landau equation

Margaret Beck\(^1,2\) , Toan T Nguyen\(^3\) , Björn Sandstede\(^4\) and Kevin Zumbrun\(^5\)

\(^1\) Department of Mathematics, Boston University, Boston, MA 02215, USA
\(^2\) Heriot-Watt University, Edinburgh, EH14 4AS, UK
\(^3\) Department of Mathematics, Pennsylvania State University, State College, PA 16803, USA
\(^4\) Division of Applied Mathematics, Brown University, Providence, RI 02912, USA
\(^5\) Department of Mathematics, Indiana University, Bloomington, IN 47405, USA

E-mail: mabeck@math.bu.edu, nguyen@math.psu.edu, Bjorn.Sandstede@Brown.edu and kzumbrun@indiana.edu

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Abstract

In an appropriate moving coordinate frame, source defects are time-periodic solutions to reaction–diffusion equations that are spatially asymptotic to spatially periodic wave trains whose group velocities point away from the core of the defect. In this paper, we rigorously establish nonlinear stability of spectrally stable source defects in the complex Ginzburg–Landau equation. Due to the outward transport at the far field, localized perturbations may lead to a highly non-localized response even on the linear level. To overcome this, we first investigate in detail the dynamics of the solution to the linearized equation. This allows us to determine an approximate solution that satisfies the full equation up to and including quadratic terms in the nonlinearity. This approximation utilizes the fact that the non-localized phase response, resulting from the embedded zero eigenvalues, can be captured, to leading order, by the nonlinear Burgers equation. The analysis is completed by obtaining detailed estimates for the resolvent kernel and pointwise estimates for Green’s function, which allow one to close a nonlinear iteration scheme.

Keywords: Ginzburg–Landau equation, nonlinear stability, sources, defects, patterns, Green’s function, Burgers equation

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(Some figures may appear in colour only in the online journal)
1. Introduction

In this paper, we study stability of source defect solutions of the complex cubic-quintic Ginzburg–Landau (qCGL) equation

\[ A_t = (1 + i\alpha)A_{xx} + A - (1 + i\beta)|A|^2 + (\gamma_1 + i\gamma_2)|A|^4. \]  

(1.1)

Here \( A = A(x, t) \) is a complex-valued function, \( x \in \mathbb{R}, t \geq 0 \) and \( \alpha, \beta, \gamma_1 \) and \( \gamma_2 \) are all real constants with \( \gamma = \gamma_1 + i\gamma_2 \) being small but non-zero. It is shown, for instance in [BN85, PSAK95, Doe96, KR00, Leg01, SS04a], that the qCGL equation exhibits a family of defect solutions known as sources (see equation (1.2)). We are interested here in establishing nonlinear stability of these solutions, under suitable spectral stability assumptions.

In general, a defect is a solution \( u_d(x, t) \) of a reaction–diffusion equation

\[ u_t = Du_{xx} + f(u), \quad u : \mathbb{R} \times \mathbb{R}^+ \rightarrow \mathbb{R}^n \]

that is time-periodic in an appropriate moving frame \( \xi = x - c_d t \), where \( c_d \) is the speed of the defect, and spatially asymptotic to wave trains, which have the form \( u_{at}(kx - \omega t; k) \) for some profile \( u_{at}(\theta; k) \) that is \( 2\pi \)-periodic in \( \theta \). Thus, \( k \) and \( \omega \) represent the spatial wave number and the temporal frequency, respectively, of the wave train. Wave trains typically exist as one-parameter families, where the frequency \( \omega = \omega_{nl}(k) \) is a function of the wave number \( k \). The function \( \omega_{nl}(k) \) is referred to as the nonlinear dispersion relation, and its domain is typically an open interval. The group velocity \( c_g(k_0) \) of the wave train with wave number \( k_0 \) is defined as

\[ c_g(k_0) = \frac{d\omega_{nl}}{dk}(k_0). \]

The group velocity is important as it is the speed with which small localized perturbations of the wave train propagate as functions of time, and we refer to [DSSS09] for a rigorous justification of this.

Defects have been observed in a wide variety of experiments and reaction–diffusion models and can be classified into several distinct types that have different existence and stability properties [vSH92, vH98, SS04a]. This classification involves the group velocities \( c_\pm_g := c_g(k_\pm) \) of the asymptotic wave trains, whose wavenumbers are denoted by \( k_\pm \). Sources are defects for which \( c_-g < c_d < c_+g \), so that perturbations are transported away from the defect core towards infinity. Generically, sources exist for discrete values of the asymptotic wave numbers \( k_\pm \), and in this sense they actively select the wave numbers of their asymptotic wave trains. Thus, sources can be thought of as organizing the dynamics in the entire spatial domain; their dynamics are inherently not localized.

For equation (1.1), the properties of the sources can be determined in some detail. We will focus on standing sources, for which \( c_d = 0 \). They have the form

\[ A_{source}(x, t) = r(x)e^{i\psi(x)}e^{-i\omega_0 t}, \]

(1.2)

where

\[ \lim_{x \to \pm\infty} \psi(x) = \pm k_0, \quad \lim_{x \to \pm\infty} r(x) = \pm r_0(k_0), \quad \omega_0 = \omega_0(k_0), \]

and where the details of the functions \( r, \varphi, r_0 \) and \( \omega_0 \) are described in lemma 2.1. In order for such solutions to be nonlinearly stable, they must first be spectrally stable, meaning roughly that the linearization about the source must not contain any spectrum in the closed right half-plane, except at the origin—see hypothesis 2.3. Our goal is to prove that, under this hypothesis, the sources are nonlinearly stable.

To determine spectral stability one must locate both the point and the essential spectrum. The essential spectrum is determined by the asymptotic wave trains. As we will see in...
section 2.2, there are two parabolic curves of essential spectrum. One is strictly in the left half-plane and the other is given by the linear dispersion relation
\[ \lambda_{\text{lin}}(\kappa) = -ic_\kappa \kappa - d\kappa^2 + O(\kappa^3) \]
for small \( \kappa \in \mathbb{R} \), where \( c_\kappa = 2k_0(\alpha - \beta) \) denotes the group velocity and
\[
d := (1 + \alpha \beta) - \frac{2k_0^2(1 + \beta^2)}{r_0^2(1 - 2\gamma_1 r_0^2)}, \quad \beta_s := \frac{\beta - 2\gamma_2 r_0^2}{1 - 2\gamma_1 r_0^2}.
\]
Thus, this second curve touches the imaginary axis at the origin and, if \( d > 0 \), then it otherwise lies in the left half-plane. In this case, the asymptotic plane waves, and therefore also the essential spectrum, are stable, at least with respect to small wave numbers \( k_0 \). Otherwise, they are unstable. Throughout the paper, we assume that \( d > 0 \). We also note that there has been previous work on stability of wave trains in equation (1.1). See, for example [Kap94].

Determining the location of the point spectrum is more difficult. For all parameter values there are two zero eigenvalues, associated with the eigenfunctions \( \partial_x A_{\text{source}} \) and \( \partial_t A_{\text{source}} \), which correspond to space and time translations, respectively. When \( \gamma_1 = \gamma_2 = 0 \), one obtains the cubic Ginzburg–Landau equation (cCGL). In this case, the sources are referred to as Nozaki–Bekki holes, and they are a degenerate family, meaning that they exist for values of the asymptotic wave number in an open interval (if one chooses the wavespeed appropriately), rather than for discrete values of \( k_0 \). Therefore, in this case there is a third zero eigenvalue associated with this degeneracy. Moreover, in the limit where \( \alpha = \beta = \gamma_1 = \gamma_2 = 0 \), which is the real Ginzburg–Landau (rGL) equation, the sources are unstable. This can be shown roughly using a Sturm–Liouville type argument: in this case, the amplitude is \( r(x) = \tanh(x) \) and so \( r'(x) \), which corresponds to a zero eigenvalue, has a single zero, which implies the existence of a positive eigenvalue.

The addition of the quintic term breaks the underlying symmetry to remove the degeneracy [Doe96] and therefore also one of the zero eigenvalues. To find a spectrally stable source, one needs to find parameter values for which both the unstable eigenvalue (from the rGL limit) and the perturbed zero eigenvalue (from the cCGL limit) become stable. This has been investigated in a variety of previous studies, including [Leg01, PSAK95, CM92, KR00, SS05, LF97]. Partial analytical results can be found in [KR00, SS05]. Numerical and asymptotic evidence in [CM92, PSAK95] suggests that the sources are stable in an open region of parameter space near the NLS limit of (1.1), which corresponds to the limit \( |\alpha|, |\beta| \to \infty \) and \( \gamma_1, \gamma_2 \to 0 \). In the present work, we will assume the parameter values have been chosen so that the sources are spectrally stable.

The main issue regarding nonlinear stability will be to deal with the effects of the embedded zero eigenvalues. This has been successfully analysed in a variety of other contexts, most notably viscous conservation laws [ZH98, HZ06, BSZ10]. Typically, the effect of these neutral modes is studied using an appropriate ansatz for the form of the solution that involves an initially arbitrary function. That function can subsequently be chosen to cancel any non-decaying components of the resulting perturbation, allowing one to close a nonlinear stability argument. The key difference here is that the effect of these eigenvalues is to cause a non-localized response, even if the initial perturbation is exponentially localized. This makes determining the appropriate ansatz considerably more difficult, as it effectively needs to be based not just on the linearized operator but also on the leading order nonlinear terms.

The remaining generic defect types are sinks (both group velocities point towards the core), transmission defects (one group velocity points towards the core, the other one away from the core), and contact defects (both group velocities coincide with the defect speed). Spectral stability implies nonlinear stability of sinks [SS04a, theorem 6.1] and transmission
defects [GSU04] in appropriately weighted spaces; the proofs rely heavily on the direction of transport and do not generalize to the case of sources. We also mention work on the nonlinear stability of sinks in the specific case of CGL in [Kap91, Kap96]. We are not aware of nonlinear stability results for contact defects, though their spectral stability was investigated in [SS04b].

We will now state our main result in more detail, in section 1.1. Subsequently, we will explain in section 1.2 the importance of the result and its relationship to the existing literature. The proof will be contained in sections 2–6.

1.1. Main result: nonlinear stability

Let \( A_{\text{source}}(x, t) \) be a source solution of the form (1.2) and let \( A(x, t) \) be the solution of (1.1) with smooth initial data \( A_{in}(x) \). In accordance with (1.2), we assume that the initial data \( A_{in}(x) \) are of the form \( R_{in}(x) e^{\phi(x,t)} \) and close to the source solution in the sense that the norm

\[
\| A_{in}(\cdot) - A_{\text{source}}(\cdot, 0) \|_{\infty} := \| e^{\frac{x^2}{2M_0}} (R_{in} - r(\cdot)) \|_{C^3[\mathbb{R}]} + \| e^{\frac{x^2}{2M_0}} (\phi_{in} - \varphi)(\cdot) \|_{C^3[\mathbb{R}]},
\]

(1.4)

where \( M_0 \) is a fixed positive constant and \( \| \cdot \|_{C^3} \) is the usual \( C^3 \)-sup norm, is sufficiently small. The assumption that the initial perturbations be bounded by Gaussians in the above sense is largely for convenience as it simplifies our analysis considerably: in particular, this assumption allows us to close our nonlinear iteration scheme more easily (and, technically, we use (1.4) in estimate (6.29) to obtain (6.25)). This localization separates, to some extent, the issue of stability of sources versus stability of the asymptotic wave trains. The stability of wave trains under non-localized perturbations was addressed only recently in [JNRZ13, SSSU12]. We believe that the source stability result remains true for initial perturbations that decay like \( |x|^{-\alpha} \) for some sufficiently large \( \alpha \) (and believe that \( \alpha = 3/2 \) is sufficient).

The solution \( A(x, t) \) will be constructed in the form

\[
A(x + p(x, t), t) = (r(x) + R(x, t)) e^{i(\varphi(x) + \phi(x,t))} e^{-i\omega t},
\]

where the function \( p(x, t) \) will be chosen so as to remove the non-decaying terms from the perturbation. The initial values of \( p(x, 0), R(x, 0), \phi(x, 0) \) can be calculated in terms of the initial data \( A_{in}(x) \).

Below we will compute the linearization of (1.1) about the source (1.2) and use this information to choose \( p(x, t) \) in a useful way. Furthermore, the linearization and the leading order nonlinear terms will imply that \( \phi(x, t) \to \phi^e(x, t) \) as \( t \to \infty \), where \( \phi^e \) represents the phase modulation caused by the zero eigenvalues. The notation is intended to indicate that \( \phi^e \) is an approximate solution to the equation that governs the dynamics of the perturbation \( \phi \). In particular, \( \phi^e \) is a solution to an appropriate Burgers-type equation that captures the leading order dynamics of \( \phi \). (See equation (1.12).) The below analysis will imply that the leading order dynamics of the perturbed source are given by the modulated source

\[
A_{\text{mod}}(x, t) := A_{\text{source}}(x, t) e^{i\phi^e(x,t)} = r(x) e^{i(\varphi(x) + \phi^e(x,t))} e^{-i\omega t}.
\]

The functions \( p(x, t)A \) and \( \phi^e(x, t) \) in \( A_{\text{mod}} \) will remove from the dynamics any non-decaying or slowly-decaying terms, resulting from the zero eigenvalues and the quadratic terms in the nonlinearity, thus allowing a nonlinear iteration scheme to be closed. To describe these functions in more detail, we define

\[
e(x,t) := \text{erf} \left( \frac{x + \zeta}{\sqrt{4dt}} \right) - \text{erf} \left( \frac{x - \zeta}{\sqrt{4dt}} \right), \quad \text{erf} \left( z \right) := \frac{1}{2\pi} \int_{-\infty}^{z} e^{-x^2} \, dx \quad (1.5)
\]

and the Gaussian-like term

\[
\theta(x,t) := \frac{1}{(1+|t|)^{1/2}} \left( e^{\frac{i\zeta g(x)}{M_{0}^{1/2}}} + e^{-\frac{i\zeta g(x)}{M_{0}^{1/2}}} \right), \quad (1.6)
\]
where \( M_0 \) is a fixed positive constant. Now define
\[
\phi^a(x, t) := -\frac{d}{2q} \log \left( \frac{1 + \delta^+(t)}{1 + \delta^-(t)} \right), \quad p(x, t) := \frac{d}{2qk_0} \log \left( \frac{1 + \delta^+(t)}{1 + \delta^-(t)} \right),
\]
where the constant \( q \) is defined in (5.10) and \( \delta^\pm = \delta^\pm(t) \) are smooth functions that will be specified later. Our main result asserts that the shifted solution \( A(x + p(x, t), t) \) converges to the modulated source with the decay rate of a Gaussian.

**Theorem 1.1.** Assume that the initial data are of the form \( A_{in}(x) = R_{in}(x)e^{ip_{in}(x)} \) with \( R_{in}, \phi_{in} \in C^3(\mathbb{R}) \). In addition, assume that \( A_{source} \) is spectrally stable in the sense of hypothesis 2.3 and that the assumptions of lemma 2.1 are satisfied. There exists a positive constant \( \epsilon_0 \) such that, if
\[
\epsilon := \| A_{in}(\cdot) - A_{source}(\cdot, 0) \|_{\infty} \leq \epsilon_0,
\]
then the solution \( A(x, t) \) to the qCGL equation (1.1) exists globally in time. In addition, there are constants \( \eta_0, C_0, M_0 > 0, \delta^\pm_\infty \in \mathbb{R} \) with \( |\delta^\pm_\infty| \leq C_0 \), and smooth functions \( \delta^\pm(t) \) so that
\[
|\delta^\pm(t) - \delta^\pm_\infty| \leq \epsilon C_0 e^{-\eta_0 t}, \quad \forall t \geq 0 \]
and
\[
\left| \frac{\partial^\ell}{\partial x^\ell} \left[ A(x + p(x, t), t) - A_{mod}(x, t) \right] \right| \leq \epsilon C_0 (1 + t)^{\epsilon} [1 + t]^{-\ell/2} + e^{\eta x|\ell|} \theta(x, t),
\]
\[
\forall x \in \mathbb{R}, \quad \forall t \geq 0,
\]
for \( \ell = 0, 1, 2 \) and for each fixed \( \kappa \in (0, \frac{1}{2}) \). In particular, \( \| A(\cdot + p(\cdot, t), t) - A_{mod}(\cdot, t) \|_{W^{1, r}} \rightarrow 0 \) as \( t \rightarrow \infty \) for each fixed \( r > \frac{1}{1-2\kappa} \).

Not only does theorem 1.1 rigorously establish the nonlinear stability of the source solutions of (1.1), but it also provides a rather detailed description of the dynamics of small perturbations. The amplitude of the shifted solution \( A(x + p(x, t), t) \) converges to the amplitude of the source \( A_{source}(x, t) \) with the decay rate of a Gaussian: \( R(x, t) \sim \theta(x, t) \). In addition, the phase dynamics can be understood as follows. If we define
\[
\delta_p(t) := -\frac{d}{2q} \log \left( \frac{1 + \delta^+(t)}{1 + \delta^-(t)} \right), \quad \delta_p(t) := \frac{d}{2qk_0} \log \left[ \frac{1 + \delta^+(t)}{1 + \delta^-(t)} \right],
\]
it then follows from (1.7) that
\[
\left| \phi^a(x, t) - \delta_p(t)e(x, t + 1) \right| + \left| p(x, t) - \delta_p(t)e(x, t + 1) \right| \leq \epsilon C_0 (1 + t)^{1/2} \theta(x, t).
\]
The function \( e(x, t) \) resembles an expanding plateau of height approximately equal to one that spreads outwards with speed \( \pm c_r \), while the associated interfaces widen like \( \sqrt{t} \); see figure 1. Hence, the phase \( \psi(x) + \phi(x, t) \) tends to \( \psi(x) + \phi^a(x, t) \), where \( \phi^a(x, t) \) looks like an expanding plateau as time increases.

As a direct consequence of theorem 1.1, we obtain the following corollary.

**Corollary 1.2.** Let \( \eta \) be an arbitrary positive constant and let \( V \) be the space–time cone defined by the constraint: \( -(c_r - \eta)t \leq x \leq (c_r + \eta)t \). Under the same assumptions as in theorem 1.1, there are positive constants \( \eta_1, C_1 \) so that the solution \( A(x, t) \) to the qCGL equation (1.1) satisfies
\[
\left| \phi^a(x, t) - \delta_p(\infty), t - \delta_p(\infty)/\omega_0 \right| \leq \epsilon C_1 e^{-\eta_1 t}
\]
for all \( (x, t) \in V \), in which \( \delta_p \) and \( \delta_p \) are defined in (1.10).
Proof. Indeed, within the cone \( V \), we have
\[
|e(x, t + 1) - 1| + \theta(x, t) \leq C_1 e^{-\eta_1 t}
\]
for some constants \( \eta_1, C_1 > 0 \). Estimate (1.11) shows that \( p(x, t) \) and \( \phi^a(x, t) \) are constants up to an error of order \( e^{-\eta_1 t} \). The main theorem thus yields the corollary at once. □

As will be seen in the proof of theorem 1.1, the functions \( \delta^\pm(t) \) will be constructed via integral formulae that are introduced to precisely capture the non-decaying part of Green’s function of the linearized operator. The choices of \( p(x, t) \) and \( \phi^a(x, t) \) are made based on the fact that the asymptotic dynamics of the translation and phase variables is governed (to leading order) by a nonlinear Burgers-type equation:
\[
\left( \partial_t - \frac{c_g}{k_0} \partial_x \right) (\phi^a \pm k_0 p) = q \left( \partial_x \phi^a \pm k_0 \partial_x p \right)^2,
\]
where \( q \) is defined in (5.10). See section 5.4. Formulae (1.7) are related to an application of the Cole–Hopf transformation to the above equation.

1.2. Difficulties and framework

In the proof, we will have to overcome two difficulties. The first, the presence of the embedded zero eigenvalues, can be dealt with using the now standard, but non-trivial, technique first developed in [ZH98]. Roughly speaking, this technique involves the introduction of an initially arbitrary function into the perturbation ansatz, which is later chosen to cancel with the non-decaying parts of Green’s function that result from the zero eigenvalues. The second difficulty is dealing with the quadratic order nonlinearity.

To illustrate this second difficulty, for the moment ignore the issue of the zero eigenvalues. Suppose we were to linearize equation (1.1) in the standard way and set
\[
A(x, t) = A_{\text{source}}(x, t) + \tilde{A}(x, t),
\]
with the hope of proving that the perturbation, \( \tilde{A}(x, t) \), decays. The function \( \tilde{A}(x, t) \) would then satisfy an equation of the form
\[
(\partial_t - \mathcal{L}) \tilde{A} = Q(\tilde{A}),
\]
where \( \mathcal{L} \) denotes the linearized operator, with the highest order derivatives being given by \( (1 + i\alpha) \partial_x^2 \), and \( Q(\tilde{A}) = O(|\tilde{A}|^2) \) denotes the nonlinearity, which contains quadratic terms. Since the temporal Green’s function (also known as the fundamental solution) for the heat operator is the Gaussian \( t^{-1/2} e^{-|x-y|^2/4t} \) centred at \( x = y \), Green’s function of \( \partial_t - \mathcal{L} \) at best behaves like a Gaussian centred at \( x = y \pm c_g t \). (In fact it is much worse, once we take into account the effects of the embedded zero eigenvalues.) Quadratic terms can have a non-trivial and subtle effect on the dynamics of such an equation: consider, for example, \( u_t = u_{xx} - u^2 \).
The zero solution is stable with respect to positive initial data, but is in general unstable. For such situations, standard techniques for studying stability are often not effective. In particular, the nonlinear iteration procedure that is typically used in conjunction with pointwise Green’s function estimates does not work when quadratic terms are present (unless they have a special conservative structure). This is because the convolution of a Gaussian (Green’s function) against a quadratic function of another Gaussian, $Q(\tilde{A})$, would not necessarily yield Gaussian behaviour. Therefore, if we were to use this standard ansatz, it would not be possible to perform the standard nonlinear iteration scheme and show that $\tilde{A}$ also decays like a Gaussian. To overcome this, we must use an ansatz that removes the quadratic terms from the equation.

Returning to the first difficulty, as mentioned above (see also lemma 2.4), the essential spectrum of $L$ touches the imaginary axis at the origin and $L$ has a zero eigenvalue of multiplicity two. The associated eigenfunctions are $\partial_t A_{\text{source}}$ and $\partial_t A_{\text{source}}$, which correspond to time and space translations, respectively. Neither of these eigenfunctions are localized in space (nor are they localized with respect to the $(R, \phi)$ coordinates—see (2.7)). This is due to the fact that the group velocities are pointing outward, away from the core of the defect, and so (localized) perturbations will create a non-local response of the phase. More precisely, the perturbed phase $\phi(x, t)$ will resemble an outwardly expanding plateau. This behaviour will need to be incorporated in the analysis if we are to close a nonlinear iteration scheme.

In the proof, we write the solution $A(x, t)$ in the form

$$A(x + p(x, t), t) = [r(x) + R(x, t)]e^{i\phi(x, t)}e^{-k_0 x},$$

and work with perturbation variables $(R(x, t), \phi(x, t))$. The advantages when working with these polar coordinates are (i) they are consistent with the phase invariance (or gauge invariance) associated with (1.1); (ii) the quadratic nonlinearity is a function of $R, \phi_\times$, and their higher derivatives, without any zero-order term involving $\phi$; (iii) based upon the leading order terms in the equation (see section 5), we expect that the time decay in the amplitude $\tilde{R}$ is faster than that of the phase $\phi$. Roughly speaking, these coordinates effectively replace the equation $u_t = u_{x x} - u^2$, which is essentially what we would have for $\tilde{A}$, with an equation like $u_t = u_{x x} - uu_x$, which is essentially what we obtain in the $(R, \phi)$ variables (but without the conservation law structure). In other words, with respect to $\tilde{A}$, the nonlinearity is relevant, but with respect to $(R, \phi)$, it is marginal [BK94].

In the case of a marginal nonlinearity, if there is an additional conservation law structure, as in for example the Burgers equation $(uu_x = (u^2)_x)/2$, then one can often exploit this structure to close the nonlinear stability argument. Here, however, that structure is absent, and so we must find another way to deal with the marginal terms. The calculations of section 5.4 show that, to leading order, the dynamics of $(R, \phi)$ are essentially governed by

$$\partial_t \begin{pmatrix} R \\ \phi \end{pmatrix} \approx \begin{pmatrix} 1 - \alpha \beta^* & \alpha r \\ -\alpha (1 + \beta^*) & r (1 + \beta^*) \end{pmatrix} \partial_x^2 \begin{pmatrix} R \\ \phi \end{pmatrix} + \begin{pmatrix} -2(\alpha + \beta^*)\phi_x \\ -2\phi_x(1 + (\beta^*)^2) \end{pmatrix} + O(R^2, \phi_x, R\phi_x),$$

where $q$ is defined in (5.10). The presence of the zero-order term $-2r^2(1 - 2\gamma_{l_1}^2)R$ in the $R$ equation implies that it will decay faster than $\phi$. In fact, the above equation implies that to leading order $R \sim \phi_x$. Moreover, if we choose an approximate solution so that $R \sim k_0\phi_{x l_0}/(r_0(1 - 2\gamma_{l_1}^2))$, then we see that $\phi$ satisfies exactly the Burgers equation given in (1.12) (up to terms that are exponentially localized). In order to close the nonlinear iteration, we will then need to incorporate these Burgers-type dynamics for $\phi$ into the ansatz, which is done exactly through the approximate solution $\phi^\circ$. This is similar to the analysis of the toy model in [BNSZ12].
When working with the polar coordinates, however, there is an apparent singularity when \( r(x) \) vanishes. Such a point is inevitable since \( r(x) \rightarrow \pm r_0 \) with \( r_0 \neq 0 \) as \( x \rightarrow \pm \infty \). We overcome this issue by writing the perturbation system as

\[
(\partial_t - \mathcal{L})U = \mathcal{N}(R, \phi, p),
\]

for \( U = (R, r\phi) \), instead of \((R, \phi)\). Here \( \mathcal{L} \) again denotes the linearized operator and \( \mathcal{N}(R, \phi, p) \) collects the remainder; see lemmas 2.2 and 5.1 for details. Note that we do not write the remainder in terms of \( U \), but leave it in terms of \( R \) and \( \phi \). Later on, once all necessary estimates for \( U(x, t) \) and its derivatives are obtained, we recover the estimates for \((R(x, t), \phi(x, t))\) from those of \( U(x, t) \), together with the observation that \( \phi(x, t) \) should contain no singularity near the origin if \( r(x)\phi(x, t) \) and its derivatives are regular; see section 6.9.

To make the above discussion rigorous, there will be four main steps. After stating some preliminary facts about sources and their linearized stability in section 2, the first step in section 3 will be to construct the resolvent kernel by studying a system of ODEs that corresponds to the eigenvalue problem. In the second step, in section 4, we derive pointwise estimates for the temporal Green’s function associated with the linearized operator. These first two steps, although non-trivial, are by now routine following the seminal approach introduced by Zumbrun and Howard [ZH98]. The third step, in section 5, is to construct the approximate ansatz for the solution of qCGL, and the final step, in section 6, is to introduce a nonlinear iteration scheme to prove stability. These last two steps are the novel and most technical ones in our analysis.

It is perhaps worth noting that recent work on the nonlinear stability of spatially periodic patterns in dissipative systems, in particular [JZ11], is closely related to this work. Both rely on two key components: (1) a suitable perturbation ansatz and (2) detailed pointwise estimates for the linearized Green’s function, with the residual error in the former sufficiently small that, when paired with the latter, the result is small enough to close a nonlinear iteration scheme. In both [JZ11] and this work, a potential issue is the appearance of quadratic terms in the nonlinearity. For the periodic wave trains treated in [JZ11], the linearized estimates decay sufficiently rapidly that such terms are asymptotically irrelevant and may be absorbed in the nonlinear iteration; for the source defects treated here, by contrast, the linearized estimates decay more slowly by a factor of \( t^{1/2} \), and so quadratic terms do not absorb. It is for this reason that we find it necessary to build a more complicated nonlinear ansatz taking into account quadratic (i.e. Burgers) order coupling, so as to remove these terms from the residual error.

In some sense, the main difference is the nature and dynamics of perturbations. For wave trains, there is enough space so that even non-localized (but bounded) phase perturbations can effectively decay like rarefaction waves. For sources, this is not possible, because the source has a fixed phase: thus we get phase fronts, which do not decay; this is then promulgated through the system. In some sense, the analogue of the latter in the wave-train case would be non-localized wavenumber perturbations, which are not understood.

To the best of the authors’ knowledge, this work is the first nonlinear stability result for a defect of source type, extending the theoretical framework to include this case. An interesting open problem at a practical level is to verify the spectral stability assumptions made here in some asymptotic regime; this is under current investigation. An important extension in the theoretical direction would be to treat the case of source defects of general reaction–diffusion equations not possessing a gauge invariance naturally identifying the phase. This would involve constructing a suitable approximate phase, sufficiently accurate to carry out a similar nonlinear analysis, a step that appears to involve substantial additional technical difficulty. We hope to address this in future work.
Universal notation. Throughout the paper, we write $g = O(f)$ to mean that there exists a universal constant $C$ so that $|g| \leq C|f|$.

2. Preliminaries

2.1. Existence of a family of sources for $q$CGL

In this subsection, we prove the following lemma concerning the existence and some qualitative properties of the source solutions defined in (1.2). This lemma relies on the results of, for example, [Doe96] and [SS04a], where the existence of sources was studied. Such results are valid near the real cubic Ginzburg–Landau limit ($\alpha = \beta = \gamma = 0$) and near the NLS limit ($|\alpha| = |\beta| = \infty$, $\gamma = 0$). We note that, in this limit, the only sources that exist are the standing ones ($c_\delta = 0$).

Lemma 2.1. Assume that $|\alpha|$ and $|\beta|$ are either both sufficiently small or both sufficiently large; in addition, once $\alpha$ and $\beta$ are fixed, we assume that $|\gamma|$ is sufficiently small. There exists a $k_0 \in \mathbb{R}$ with $|k_0| < 1$ such that a source solution $A_{source}(x, t)\) of (1.1) of the form (1.2) exists and satisfies the following properties.

1. The functions $r(x)$ and $\phi(x)$ are $C^\infty$. Let $x_0$ be a point at which $r(x_0) = 0$. Necessarily, $r'(x_0) \neq 0$ and $r(x_0) = \phi(x_0) = 0$.

2. The functions $r$ and $\phi$ satisfy $r(x) \to \pm r_0(k_0)$ and $\phi(x) \to \pm k_0$ as $x \to \pm \infty$, respectively, where $r_0$ is defined in (2.1). Furthermore, \[ \left| \frac{d^\ell}{dx^\ell} \left( r(x) \mp r_0(k_0) \right) \right| + \left| \frac{d^{\ell+1}}{dx^{\ell+1}} \left( \phi(x) \mp k_0 \right) \right| \leq C_0 e^{-\eta_0 |x|}, \]
for integers $\ell \geq 0$ and for some positive constants $C_0$ and $\eta_0$.

3. As $x \to \pm \infty$, $A_{source}(x, t)$ converges to the wave trains $A_{wave}(x, t; \pm k_0) = \pm r_0(k_0) e^{i(\pm k_0 x - \omega_0(k_0)t)}$, respectively, with \[ r_0^2 = 1 - k_0^2 + \gamma_1 r_0^4, \quad \omega_0(k_0) = \beta + (\alpha - \beta) k_0^2 + (\beta \gamma_1 - \gamma_3) r_0^2. \] (2.1) Necessarily, $\omega_0 = \omega_0(k_0)$.

4. If $k_0 \neq 0$, the asymptotic group velocities $c_\pm := \frac{\partial \omega_0}{\partial k}|_{k=\pm k_0}$ have opposite sign at $\pm \infty$ and satisfy \[ c_\pm = \mp c_g, \quad c_g := 2k_0(\alpha - \beta_\pm), \] (2.2) where $\beta_\pm$ is defined (1.3). Without loss of generality, we assume that $c_g > 0$.

Before proving this, we recall that, for the cubic GLE, equation (1.1), an explicit formula for the (traveling) source, which is known in this case as a Nozaki–Bekki hole, is given in [BN85,Leg01,PSAK95]. These Nozaki–Bekki holes are degenerate solutions of CGL in the sense that they exist in a non-transverse intersection of stable and unstable manifolds. More precisely, when $\gamma_1 = \gamma_2 = 0$, the formula for the standing Nozaki–Bekki holes is given by \[ A_{source}(x, t) = r_0 \tanh(\kappa x)e^{-i\kappa_0 t}e^{-i k_0(x - \delta \log(2\cosh(\kappa x))}, \]
where $r_0 = \sqrt{1 - k_0^2}$, $k_0 = -\delta \kappa$, and $\delta$ and $\kappa$ are defined by \[ \kappa^2 = \frac{\alpha - \beta}{\alpha - \beta} \delta^2 - 3\delta(1 + \alpha^2), \quad \delta^2 + \frac{3(1 + \alpha \beta)}{(\beta - \alpha)} \delta - 2 = 0, \]
with $\delta$ chosen to be the root of the above equation such that $\delta(\alpha - \beta) < 0$. See [Leg01] for details. The asymptotic phases and group velocities are \[ \omega_{\pm}(k_0) = \beta + (\alpha - \beta) k_0^2, \quad c_{\pm} = \frac{\partial \omega_{\pm}(k)}{\partial k}|_{k=\pm k_0} = \pm 2k_0(\alpha - \beta). \] (2.3)
which are the identities (2.1) and (2.2) with $\gamma_1 = \gamma_2 = 0$. We note that since $c_g^+ + c_g^- = 0$, the asymptotic group velocities must have opposite signs. To see that the solution is really a source, one can check that

$$c_g^+ - c_g^- = -4\delta\kappa(\alpha - \beta) = 4\text{sgn}(\kappa)\kappa\delta(\alpha - \beta).$$

If $\kappa > 0$, this difference is positive, and $c_g^+$, which is then the group velocity at $+\infty$ is positive, and $c_g^-$, which is then the group velocity at $-\infty$, is negative. Thus, the solution is indeed a source. If $\kappa < 0$, the signs of $c_g^\pm$ are reversed, but so are the ends to which they correspond. Thus, the solution is indeed a source in all cases.

The (qCGL) equation (1.1) is a small perturbation of the cubic CGL, and it has been shown that the above solutions persist as standing sources for an open set of parameter values [Doe96, SS04a]. Furthermore, they are constructed via a transverse intersection of the two-dimensional centre-stable manifold of the asymptotic wave train at infinity and the two-dimensional centre-unstable manifold of the wave train at minus infinity that is unfolded with respect to the wavenumbers of these wave trains; in particular, the standing sources connect wave trains with a selected wavenumber. These facts are essential for the proof of lemma 2.1.

**Proof of lemma 2.1.** As mentioned above, standing sources have been proven to exist in [Doe96] and [SS04a]. Let $A_{\text{source}}(x, t)$ be that source, which we can write in the form

$$A_{\text{source}}(x, t) = r(x)e^{i(\gamma(x) - \omega t)}.$$

Plugging this into (1.1), we find that $(r, \psi)$ solves

$$0 = r_{xx} + r - r\psi_x^2 - 2ar_x\psi_x - ar\psi_{xx} - r^3 + \gamma_1 r^5$$

$$0 = r\psi_{xx} + 2r_x\psi_x + ar_{xx} + \omega_0 r - ar\psi_x^2 - \beta r^3 + \gamma_2 r^5.

It is shown in [Doe96] that there exists a locally unique wavenumber $k_0$ and a smooth solution $(r, \psi)$ of (2.3) so that $A_{\text{source}}(x, t)$ converges to the wave trains of the form

$$A_{\text{wt}}(x, t; \pm k_0) = \pm r_0(k_0)e^{i(\pm k_0x - \omega_0(t))},$$

respectively as $x \to \pm \infty$. Putting this asymptotic ansatz into (2.3) then yields

$$\gamma_1 r_0^3 - r_0^3 + 1 - k_0^2 = 0,$$

$$\omega_0(k_0) = \alpha k_0^2 + \beta r_0^2 - \gamma_2 r_0^3.$$

Rearranging terms then gives (2.1), and hence item 3 in the lemma. Differentiating the above identities with respect to $k$ and solving for $\frac{\partial A_{\text{source}}(k)}{\partial k}$, we obtain item 4 as claimed.

For the first item, since $r(x) \to \pm r_0$ with $r_0 \approx \sqrt{1 - k_0^2} \neq 0$ (because $\gamma_1 \approx 0$), there must be a point at which $r(x)$ vanishes. Without loss of generality, we assume that $r(0) = 0$. Since $r(x)$ and $\psi(x)$ are smooth, evaluating system (2.3) at $x = 0$ gives

$$r_{xx}(0) - 2ar_x(0)\psi_x(0) = 0,$$

$$2r_x(0)\psi_x(0) + ar_{xx}(0) = 0.$$

These equations imply that $r_{xx}(0) = r_x(0)\psi_x(0) = 0$. We now argue that $r_x(0)$ must be non-zero. Otherwise, since $A_{\text{source}}(0) = r(0)e^{i(\psi(0) - \omega t)}$, $\partial_r A_{\text{source}}(0) = [r_x(0) + i\psi_x(0)]e^{i(\psi(0) - \omega t)}$, and $\partial_{xx} A_{\text{source}}(0) = [r_{xx}(0) + 2ir_x(0)\psi_x(0) + i(r_x(0)\psi_{xx}(0) - r(0)\psi_x^2(0))e^{i(\psi(0) - \omega t)}]$, we would have $A_{\text{source}}(0) = \partial_r A_{\text{source}}(0) = \partial_{xx} A_{\text{source}}(0) = 0$ and so $\partial_r^2 A_{\text{source}}(0) = 0$ for all $k \geq 0$ by repeatedly using equation (1.1). This would imply at once that $A_{\text{source}}(x) \equiv 0$. Item 1 is thus proved. The exponential decay stated in item 2 is a direct consequence of the fact that the solutions in [Doe96, SS04a] are shown to lie in the intersection of stable and unstable manifolds of saddle equilibria.
2.2. Linearization

In order to linearize (1.1) around $A_{\text{source}}(x, t)$, we introduce the perturbation variables $(R, \phi)$ via

$$A(x, t) = [r(x) + R(x, t)]e^{i(\psi(x) + \phi(x, t))}e^{-i\omega_0 t}, \quad (2.4)$$

where $r$ and $\psi$ are the amplitude and phase of $A_{\text{source}}(x, t)$. Throughout the paper, we shall work with the vector variable

$$U := \begin{pmatrix} R \\ r\phi \end{pmatrix}$$

for the perturbation $(R, \phi)$. We have the following lemma.

**Lemma 2.2.** If the $A(x, t)$, defined in (2.4), solves (1.1), then the linearized dynamics of $U$ are

$$(\partial_t - L)U = 0,$$

where

$$L := D_2\partial_x^2 - 2\psi_x \partial_x + D_0^\psi(x) + D_0(x)$$

with

$$D_1 := \begin{pmatrix} \alpha & 1 \\ -1 & \beta \end{pmatrix}, \quad D_2 := \begin{pmatrix} 1 & -\alpha \\ \alpha & 1 \end{pmatrix},$$

$$D_0^\psi(x) := \begin{pmatrix} 0 & \frac{1}{r}(2\psi_x r_x + \alpha r_{xx}) \\ 0 & \frac{1}{r}(2\psi r_x - r_{xx}) \end{pmatrix}, \quad D_0(x) := \begin{pmatrix} 1 - 3r^2 - \alpha \psi_{xx} - \psi_x^2 + 5\gamma_1 r^4 & 0 \\ 0 & 3r^2 - \alpha \psi_x^2 + \psi_{xx}^2 + 5\gamma_2 r^4 \end{pmatrix}. \quad (2.5)$$

**Proof.** Plugging (2.4) into (1.1) and using equation (2.3), we obtain the linearized system

$$R_t = r_{xx} - 2\alpha \psi_x R_x + (1 - 3r^2 - \alpha \psi_{xx} - \psi_x^2 + 5\gamma_1 r^4)R - 3(\alpha r_x + r \psi_x)\phi_x,$$

$$r \phi_t = r \phi_{xx} - (2\alpha r \psi_x - 2r_x)\phi_x + \alpha \psi_{xx}^2 + 2\psi_x R_x + (\omega_0 + \psi_{xx} - \alpha \psi_x^2 - 3\beta r^2 + 5\gamma_2 r^4)R.$$

A rearrangement of terms yields the lemma. \qed

Next, note that $L$ is a bounded operator from $H^2(\mathbb{R}; \mathbb{C}^2) \cap L^\infty(\mathbb{R}; \mathbb{C}^2)$ to $L^2(\mathbb{R}; \mathbb{C}^2)$, where the function space $H^2(\mathbb{R}; \mathbb{C}^2)$ consists of functions $U = (u_1, u_2)$ so that $u_1 \in H^2(\mathbb{R}; \mathbb{C})$ and $\partial_t u_1, \partial_x^2 u_2 \in L^2(\mathbb{R}; \mathbb{C})$. (Note that we do not require that $u_2 \in L^2$, since $r \phi$ need not be localized.) We make the following assumption about the spectral stability of $L$.

**Hypothesis 2.3.** The spectrum of the operator $L$ in $L^2(\mathbb{R}, \mathbb{C}^2)$ satisfies the following two conditions:

- The spectrum does not intersect the closed right half-plane, except at the origin.
- In the weighted space $L^2_\eta(\mathbb{R}, \mathbb{C}^2)$, defined by $\|u\|^2_\eta = \int e^{-\eta |x|} |u(x)|^2 dx$ with $\eta > 0$ sufficiently small, there are exactly two eigenvalues in the closed right half-plane, and they are both at the origin.

**Lemma 2.4.** The essential spectrum of the linearized operator $L$ in $L^2(\mathbb{R}, \mathbb{C}^2)$ lies entirely in the left half-plane $\{\Re \lambda \leq 0\}$ and touches the imaginary axis only at the origin as a parabolic curve; see figure 2. In addition, the nullspace of $L$ is spanned by functions $V_1(x)$ and $V_2(x)$, defined by

$$V_1(x) := \begin{pmatrix} r_x \\ r \psi_x \end{pmatrix} \quad \text{and} \quad V_2(x) := \begin{pmatrix} 0 \\ r \end{pmatrix}, \quad (2.7)$$
The essential spectrum of $L$ is contained in the shaded region determined by the algebraic curves $\lambda_{\pm}^1(\kappa), \lambda_{\pm}^2(\kappa)$. The set $\Omega_\vartheta$ and its boundary $\Gamma = \partial \Omega_\vartheta$, defined in (3.2), are also illustrated. Note that they both lie outside of the essential spectrum.

and the corresponding adjoint eigenfunctions, $\psi_{1,2}(x)$, are both exponentially localized: $|\psi_{1,2}(x)| \leq C e^{-\eta_0|x|}$ for some $C, \eta_0 > 0$.

**Proof.** It follows directly from the translation and gauge invariance of qCGL that $\partial_x A_{\text{source}}(x, t)$ and $\partial_t A_{\text{source}}(x, t)$ are exact solutions of the linearized equation about $A_{\text{source}}(x, t)$. Consequently, $V_1(x), V_2(x)$ belong to the kernel of $L$. By hypothesis 2.3 these are the only elements of the kernel.

Next, by standard spectral theory (see, for instance, Henry [Hen81]), the essential spectrum of $L$ in $L^2(\mathbb{R}, \mathbb{C}^2)$ is the same as that of the limiting, constant-coefficient operator $L_\pm$ defined by

$$L_\pm := D_2 \mp 2k_0 D_1 \partial_x + D_0^\infty,$$

with $D_0^\infty := -2i \gamma r_0^2 \beta - 2\gamma_2 r_0^2$. To find the spectrum of $L_\pm$, we denote by $\Lambda_\pm^\kappa$ for each fixed $\kappa \in \mathbb{R}$ the constant matrix

$$\Lambda_\pm^\kappa := \kappa^2 D_2 \mp 2i \kappa D_1 - D_0^\infty,$$

and let $\lambda_{\pm}^1(\kappa), \lambda_{\pm}^2(\kappa)$ be the two eigenvalues of $\Lambda_\pm^\kappa$ with $\text{Re} \lambda_{\pm}^1(0) = 0$. Clearly, $\text{Re} \lambda_{\pm}^2(\kappa) < 0$ for all $\kappa \neq 0$, and $\lambda_{\pm}^1(\kappa)$ touches the origin as a parabolic curve. Precisely, we have

$$\lambda_{\pm}^1(\kappa) = -i \kappa^2 \kappa - \kappa^2 + O(\kappa^3)$$

for sufficiently small $\kappa$, whereas $\lambda_{\pm}^2(\kappa)$ is bounded away from the imaginary axis. It then follows that the spectrum of $L_\pm$ is confined to the shaded region to the left of the curves $\lambda_{\pm}^1(\kappa)$ for $\kappa \in \mathbb{R}$, as depicted in figure 2.

Finally, one can see from studying the asymptotic limits of (3.10) that the adjoint eigenfunctions are exponentially localized; alternatively, this property was shown more generally in [SS04a, corollary 4.6].
3. Construction of the resolvent kernel

We now construct the resolvent kernel and derive resolvent estimates for the linearized operator $L$, defined in (2.5). These resolvent estimates will be used below, in section 4, to obtain pointwise estimates for Green’s function.

Let $G(x, y, \lambda)$ denote the resolvent kernel associated with the operator $L$, which is defined to be the distributional solution of the system

$$
(\lambda - L)G(\cdot, y, \lambda) = \delta_y(\cdot)
$$

for each $y \in \mathbb{R}$, where $\delta_y(x)$ denotes the Dirac delta function centred at $x = y$. The main result in this section provides pointwise bounds on $G(x, y, \lambda)$. We divide it into three different regions: low-frequency ($\lambda \to 0$), mid-frequency ($\theta \leq |\lambda| \leq M$) and high-frequency ($\lambda \to \infty$). The reason for this separation has to do with the behaviour of the spatial eigenvalues of the four-dimensional, first-order ODE associated with (3.1), which is given in (3.4). These spatial eigenvalues determine the key features of the resolvent kernel $G$ and depend on the spectral parameter $\lambda$ in such a way so that the asymptotics of $G$ can be characterized differently in these three regions.

Define

$$
\Omega_\theta := \left\{ \lambda \in \mathbb{C} : \text{Re} \lambda \geq -\theta_1 - \theta_2 |\text{Im} \lambda|, |\lambda| \geq \theta_3 \right\},
$$

(3.2)

where the real constants $\theta_1, \theta_2, \theta_3$ are chosen so that $\Omega_\theta$ does not intersect the spectrum of $L$; see figure 2 and lemma 2.4. The following propositions are the main results of this section, and their proofs will be given below.

**Proposition 3.1 (High-frequency bound).** There exist positive constants $\theta_1, \theta_2, M, C$ and $\eta$ so that

$$
|\partial^k_x G(x, y, \lambda)| \leq C|\lambda|^{k-1/2} e^{-\eta|\lambda|^{1/2}|x-y|}
$$

for all $\lambda \in \Omega_\theta \cap \{||\lambda| \geq M\}$ and $k = 0, 1, 2$.

**Proposition 3.2 (Mid-frequency bound).** For any positive constants $\theta_3$ and $M$, there exists $C = C(M, \theta_3)$ sufficiently large so that

$$
|\partial^k_x G(x, y, \lambda)| \leq C(M, \theta_3),
$$

for all $\lambda \in \Omega_\theta \cap \{||\lambda| \leq M\}$ and $k = 0, 1, 2$.

**Proposition 3.3 (Low-frequency bound).** There exists an $\eta_3 > 0$ sufficiently small such that, for all $\lambda$ with $|\lambda| < \eta_3$, we have the expansion

$$
G(x, y, \lambda) = \frac{1}{\lambda} \sum_{j=1}^2 e^{v_j(\lambda)|x|} V_j(x) \left( \psi_j(y), \cdot \right) e^{\eta_3(|\lambda|^{1/2}|x-y|)} + O(e^{\eta_3(|\lambda|^{1/2}|x-y|)}) + O(e^{-\eta|x-y|}),
$$

where

$$
v_j(\lambda) = -\frac{\lambda}{c_g} + \frac{d\lambda^2}{c_g^3} + O(\lambda^3).
$$

Here $c_g > 0$ is the group velocity defined in (2.2), $V_1, V_2$ are the eigenfunctions defined in (2.7), and the adjoint eigenfunctions satisfy $\psi_j(y) = O(e^{-\eta|y|})$, for some fixed $\eta > 0$. 

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3.1. Spatial eigenvalues

We first consider the linearized eigenvalue problem

\[(\lambda - \mathcal{L})U = 0\]  \hspace{1cm} (3.3)

and derive necessary estimates on the behaviour of solutions as \(x \to \pm \infty\). We write the eigenvalue problem (3.3) as a four-dimensional first-order ODE system. For simplicity, we denote

\[
B_0(x, \lambda) := D^{-1}[\lambda I - D_0(x) - D\varphi_0(x)], \quad C_0 := D^{-1}D_1 = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}
\]

where \(I\) denotes the 2 \(\times\) 2 identity matrix and \(D_{0,1,2}\) are defined in (2.6). Let \(W = (U, U_x)\) be the new variable. By (2.5), the eigenvalue problem (3.3) then becomes

\[
W_x = \mathcal{A}(x, \lambda)W, \quad \mathcal{A}(x, \lambda) := \begin{pmatrix} 0 & I \\ B_0(x, \lambda) & 2\varphi_x C_0 \end{pmatrix}.
\]  \hspace{1cm} (3.4)

Let \(A_\pm(\lambda)\) be the asymptotic limits of \(\mathcal{A}(x, \lambda)\) at \(x = \pm \infty\), and let

\[
B_0(\lambda) := D^{-1}[\lambda I - D_0^\infty]
\]

so that, by lemma 2.1, \(B_0(x, \lambda) \to B_0(\lambda)\) as \(x \to \pm \infty\). Thus, we have

\[
A_\pm(\lambda) = \begin{pmatrix} 0 & I \\ B_0(\lambda) & \pm 2k_0 C_0 \end{pmatrix}.
\]

The solutions of the limiting ODE system \(W_x = A_\pm(\lambda)W\) are of the form \(W_\infty(\lambda)e^{\nu(\lambda)x}\), where \(W_\infty = (w, v_\pm(\lambda)w)\), \(w \in C^2\), and \(v_\pm(\lambda)\) are the eigenvalues of \(A_\pm(\lambda)\). These eigenvalues are often referred to as spatial eigenvalues, to distinguish them from the temporal eigenvalue parameter \(\lambda\), and they satisfy

\[
\text{det}(B_0(\lambda) \pm 2k_0 v_\pm(\lambda)C_0 - \nu^2(\lambda)I) = 0.
\]  \hspace{1cm} (3.5)

The behaviour of these spatial eigenvalues as functions of \(\lambda\), which determines the key properties of the resolvent kernel, can be understood by considering the limiting cases \(\lambda \to 0\) and \(|\lambda| \to \infty\), along with the intermediate regime between these limits. The easiest cases are the mid- and high-frequency regimes.

3.2. Mid- and high-frequency resolvent bounds

**Proof of proposition 3.1.** The spatial eigenvalues and resolvent kernel can be analysed in this regime using the following scaling argument. Define

\[
\tilde{x} = |\lambda|^{1/2}x, \quad \tilde{\lambda} = |\lambda|^{-1/2}, \quad \tilde{W}(\tilde{x}) = W(|\lambda|^{-1/2}x).
\]

In these scaled variables, (3.4) becomes

\[
\tilde{W}_x = \tilde{A}(\tilde{\lambda})\tilde{W} + O(|\lambda|^{-1/2}\tilde{W}) \quad \text{with} \quad \tilde{A}(\tilde{\lambda}) = \begin{pmatrix} 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \\ \frac{\tilde{\lambda}}{1 + \alpha^2} & \frac{\alpha \tilde{\lambda}}{1 + \alpha^2} & 0 & 0 \\ \frac{\alpha \tilde{\lambda}}{1 + \alpha^2} & \frac{\tilde{\lambda}}{1 + \alpha^2} & 0 & 0 \end{pmatrix}
\]  \hspace{1cm} (3.6)

The eigenvalues of \(\tilde{A}(\tilde{\lambda})\) are

\[
\pm \sqrt{\frac{\tilde{\lambda}(1 \pm i\alpha)}{1 + \alpha^2}}.
\]

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Figure 3. Illustration of the spatial eigenvalues (depicted on the real axis for convenience). For \(|\lambda| \to \infty\) they are well-separated from the imaginary axis. When \(|\lambda| \to 0\) (see section 3.3), one spatial eigenvalue, determined by the spatial limit \(x \to \pm \infty\), approaches the origin.

For \(\lambda \in \Omega_\rho, \tilde{\lambda}\) is on the unit circle and bounded away from the negative real axis. Therefore, there exists some fixed \(\eta > 0\) so that the real part of \(|\lambda|(1 \pm i\omega)\) is greater than \(\eta\) for all \(\lambda \in \Omega_\rho\) as \(|\lambda| \to \infty\). See figure 3(a).

Furthermore, we can use this asymptotic behaviour to determine information about the resolvent kernel. Let \(\tilde{P}^{s,u}(\tilde{\lambda})\) be the eigenprojections associated with \(\tilde{A}(\tilde{\lambda})\) onto the stable and unstable subspaces, which are smooth with respect to \(\lambda\). We define \(\tilde{W}^{s,u}(\tilde{x}) = \tilde{P}^{s,u}(\tilde{\lambda}) \tilde{W}(\tilde{x})\), then

\[
\begin{pmatrix}
\tilde{W}^s \\
\tilde{W}^u
\end{pmatrix}
= \begin{pmatrix}
\tilde{A}^s(\tilde{\lambda}) & 0 \\
0 & \tilde{A}^u(\tilde{\lambda})
\end{pmatrix}
\begin{pmatrix}
\tilde{W}^s \\
\tilde{W}^u
\end{pmatrix}
+ \mathcal{O}(|\lambda|^{-1/2}|\tilde{W}|)
\]

(3.7)

where the matrices \(\tilde{A}^s,u(\tilde{\lambda})\) satisfy

\[\text{Re}(\tilde{A}^s(\tilde{\lambda})W, W)_{\mathbb{C}^4} \leq -\eta|W|^2,\quad \text{Re}(\tilde{A}^u(\tilde{\lambda})W, W)_{\mathbb{C}^4} \geq \eta|W|^2,\]

for all \(W \in \mathbb{C}^4\). Here, \((\cdot, \cdot)_{\mathbb{C}^4}\) denotes the usual inner product in \(\mathbb{C}^4\). Taking the inner product of equation (3.7) with \((\tilde{W}^s, \tilde{W}^u)\), we obtain

\[
\frac{1}{2} \frac{d}{dx} |\tilde{W}^s|^2 \leq -\eta|\tilde{W}^s|^2 + \mathcal{O}(|\lambda|^{-1/2}|\tilde{W}|^2),
\]

\[
\frac{1}{2} \frac{d}{dx} |\tilde{W}^u|^2 \geq \eta|\tilde{W}^u|^2 + \mathcal{O}(|\lambda|^{-1/2}|\tilde{W}|^2).
\]

This shows that for \(|\lambda|\) is sufficiently large, the stable solutions decay faster than \(e^{-\eta|x|}\) and the unstable solutions grow at least as fast as \(e^{\eta|x|}\). Thus, the resolvent kernel \(\tilde{G}_W(\tilde{x}, \tilde{y}, \tilde{\lambda})\) associated with the scaled equation (3.6) satisfies the uniform bound

\[|\tilde{G}_W(\tilde{x}, \tilde{y}, \tilde{\lambda})| \leq C e^{-\eta|x-y|}.
\]

Going back to the original variables, the resolvent kernel associated with (3.4) satisfies

\[|G_W(x, y, \lambda)| \leq C e^{-\eta|x-y|}.
\]

The resolvent kernel \(G(x, y, \lambda)\) is by definition is the two-by-two upper-left block of the matrix \(G_W(x, y, \lambda)\). This proves the proposition.

**Proof of proposition 3.2.** The proof is immediate by the analyticity of \(G(x, y, \lambda)\) in \(\lambda\).
3.3. Low frequency resolvent bounds via exponential dichotomies

As $\lambda$ approaches the origin, it approaches the boundary of the essential spectrum, due to lemma 2.4. Here the asymptotic matrices $A_{\pm}(\lambda)$ lose hyperbolicity, and there is one spatial eigenvalue from each $+\infty$ and $-\infty$ that approaches zero. Equation (3.5) implies that this eigenvalue satisfies $\nu_{\pm}(\lambda) = O(\lambda)$. More precisely, we have

$$
\nu^+(\lambda) = -\lambda/c^+_g + \frac{d\lambda^2}{c^+_g} + O(\lambda^3),
$$

where $c^\pm_g$ is defined as in (2.2). Since $c^+_g > 0$ and $c^-_g < 0$, we see that, at $+\infty$, there are three spatial eigenvalues that are bounded away from the imaginary axis as $\lambda \to 0$, two of which have positive real parts and one of which has a negative real part, whereas at $-\infty$ there are three spatial eigenvalues that are bounded away from the imaginary axis as $\lambda \to 0$, two of which have negative real parts and one of which has a positive real part. See figures 3(b) and (c).

Consider the positive half line $x \geq 0$. Let $\nu^+(\lambda)$ denote the spatial eigenvalue, coming from $+\infty$, that has a negative real part bounded away from zero for all $\lambda > 0$ and hence corresponds to a strong-stable direction. Similarly, let $\nu^-(\lambda)$ denote the spatial eigenvalue, coming from $+\infty$, that has a negative real part that is $O(\lambda)$ as $\lambda \to 0$, and hence corresponds to a centre-stable direction. We then have $\nu^+(\lambda) = -\eta + O(\lambda)$ for some $\eta > 0$ and $\nu^-(\lambda) = -\lambda/c^+_g + O(\lambda^2)$ as $\lambda \to 0$. A direct application of the conjugation lemma (see, for example, [ZH98]) shows that a basis of bounded solutions to equation (3.4) on the positive half line $\mathbb{R}^+$ consists of

$$
e^{\nu^+(\lambda)x} W^e_1(x, \lambda), \quad e^{\nu^-(\lambda)x} W^s_1(x, \lambda)
$$

where $W^e_1(x, \lambda)$ approaches $W^e_\infty(\lambda)$ exponentially fast as $x \to +\infty$ for each $j = e, s$. The eigenfunctions $V_{1,2}$, defined in (2.7), correspond to solutions $W_{1,2}$ of (3.4) for $\lambda = 0$. Since,

$$
\lim_{x \to +\infty} W_1(x) = \begin{pmatrix} 0 \\ r_0 k_0 \\ 0 \\ 0 \end{pmatrix}, \quad \lim_{x \to -\infty} W_2(x) = \begin{pmatrix} 0 \\ \pm r_0 \\ 0 \\ 0 \end{pmatrix}
$$

we see that

$$
W_{1,2}(x) = a^{1,2}_1 e^{\nu^+(\lambda)x} W^e_1(x, 0) + a^{1,2}_2 e^{\nu^-(\lambda)x} W^s_1(x, 0), \quad x \geq 0,
$$

where $a^{1,2}_1 \neq 0$. There are also corresponding solutions $e^{-\nu_j(\lambda)y} \Psi^j_e(y, \lambda)$ of the adjoint equation associated with (3.4),

$$
W^*_e = -A(x, \lambda)^T W^e.
$$

For $\lambda = 0$, the adjoint eigenfunctions are denoted by $\Psi_{1,2}(x)$, and they are related to the adjoint eigenfunctions $\psi_{1,2}$ via $(\Psi_{1,2}) = D^T T \psi_{1,2}$, where $(\Psi_{1,2})$ denotes the second component of $\Psi_{1,2}$. We can choose $\Psi_{1,2}$ so that $\{W_{1,2}(x), \Psi_{1,2}(x)\}$ forms a basis of $\mathbb{C}^4$ for each fixed $x$.

The four-dimensional ODE associated with equation (3.1) is just system (3.4) with an additional non-autonomous term corresponding to the Dirac delta function. Let $\Phi^{(\lambda)}_0$ denote the associated exponential dichotomy (see, for instance, [BHSZ10, definition 3.2]). In other words, $\Phi^{(\lambda)}_0(x, y, \lambda)$ decays to zero exponentially fast as $x, y \to \infty$, for $x \geq y \geq 0$, and $\Phi^{(\lambda)}_0(x, y, \lambda)$ decays to zero exponentially fast as $x, y \to \infty$, for $y \geq x \geq 0$. The resolvent kernel, which solves (3.1), on the positive half line corresponds to the upper-left two-by-two block of $\Phi^{(\lambda)}_0$. By definition and equation (3.9), we can take $\Phi^{(\lambda)}_0(x, y, \lambda)$ to be of the form

$$
\Phi^{(\lambda)}_0(x, y, \lambda) = e^{\nu^+(\lambda)(x-y)} W^e_1(x, \lambda) \langle \Psi^e_1(y, \lambda), \cdot \rangle + e^{\nu^-(\lambda)(x-y)} W^s_1(x, \lambda) \langle \Psi^s_1(y, \lambda), \cdot \rangle
$$

(3.11)
for $0 \leq y \leq x$, where all terms are analytic in $\lambda$. Here $\langle \cdot, \cdot \rangle$ denotes the usual inner product in $\mathbb{C}^4$. Similarly, since the two spatial eigenvalues that correspond to the unstable direction have positive real parts bounded away from zero, we have $\Phi^u_{\pm}(x, y, \lambda) = O(e^{-|\lambda|^2/M})$ for $0 \leq x \leq y$. A similar construction can be obtained for the negative half line using the unstable spatial eigenvalues $\nu_{u, \pm}(\lambda)$.

We need to extend the exponential dichotomy so that it is valid not just on the half line, but on the entire real line. This extension will not be analytic, as there is an eigenvalue at the origin, which will correspond to a pole in the resolvent kernel. However, we can construct the extension so that it is meromorphic. The strategy will be similar to that of [BHSZ10, section 4.4] and [BSZ10, section 4.2].

Write
\[ E^j(\lambda) := \text{Rg} \Phi_j(0, 0, \lambda), \quad E^j(\lambda) := \text{Rg} \Phi_j(0, 0, \lambda), \quad j = s, u. \]

And note that $\text{span}[W_1, W_2] = E_s^+(0)$. Similarly, $\text{span}[W_1, W_2] = E_u^+(0)$. Next, set
\[ E^0_s := \text{span}[W_1(0), W_2(0)], \quad E^0_u := \text{span}[\Psi_1(0), \Psi_2(0)], \]

so that $E^0_s \oplus E^0_u = \mathbb{C}^4$. The following lemma is analogous to [BHSZ10, lemma 4.10] and [BSZ10, lemma 6], and more details of the proof can be found in those papers.

**Lemma 3.4.** There exists an $\epsilon > 0$ sufficiently small such that, for each $\lambda \in B_\epsilon(0) \setminus \{0\}$, there exists a unique map $h^+(\lambda) : E^+_s(\lambda) \to E^+_s(\lambda)$ such that $E^+_s(\lambda) = \text{graph } h^+(\lambda) := \{W + h^+(\lambda)(W) : W \in E^+_s(\lambda)\}$. Furthermore, the mapping $h^+_0(\lambda) := h^+(\lambda)|_{E_0^s}$ is of the form
\[ h^+_0(\lambda) : E^0_s \to E^0_s, \quad h^+_0(\lambda) = h^+_p(\lambda) + h^+_m(\lambda), \quad (3.12) \]

where $h^+_s(\lambda)$ is analytic for $\lambda \in B_\epsilon(0) \setminus \{0\}$ and
\[ h^+_p(\lambda)W^\psi = \frac{1}{\lambda}MW^\psi \quad (3.13) \]

for $W^\psi \in E^0_s$, where $M = (M^\psi(0))^{-1}$ and $M^\psi(0)$ is a $4 \times 4$ matrix as defined in (3.16). Similarly, for each $\lambda \in B_\epsilon(0) \setminus \{0\}$, there is a unique map $h^-_s(\lambda) : E^+_s(\lambda) \to E^+_s(\lambda)$ so that $E^+_s(\lambda) = \text{graph } h^-_s(\lambda)$. This map has a meromorphic representation analogous to the one given above for $h^+_s(\lambda)$, with $h^-_m(\lambda) = -h^-_m(\lambda)$.

**Proof.** In this proof we will use the coordinates $(W^s, W^\psi) \in E^0_s \oplus E^0_s$. Because the exponential dichotomies $\Phi^s_{\pm}(x, y, \lambda)$ and $\Phi^u_{\pm}(x, y, \lambda)$ are analytic for all $x \geq y \geq 0$ and $y \geq x \geq 0$, respectively, there exist functions $h^s(\lambda), g^s(\lambda)$ that are analytic for all $\lambda$ near zero and such that
\[ E^s_s(\lambda) : \quad W = W^s + \lambda h^s(\lambda) W^\psi, \quad \tilde{W}^s \in E^0_s, \]
\[ E^s_u(\lambda) : \quad W = W^s + \lambda g^s(\lambda) W^\psi, \quad \tilde{W}^u \in E^0_s, \]
\[ E^u_u(\lambda) : \quad W = W^\psi + g^s(\lambda) W^s, \quad W^\psi \in E^0_s. \]

The superscripts indicate the range of the associated function. For example, $g^s(\lambda)W^s \in E^0_s$ for all $W^s \in E^0_s$. We want to write $E^s_s(\lambda)$ as the graph of a function $h_s(\lambda) : E^s_s(\lambda) \to E^s_s(\lambda)$. This requires
\[ \tilde{W}^s + \lambda h^s(\lambda) W^\psi = \left\{ W^\psi + g^s(\lambda) W^s \right\} + \left\{ W^s + \lambda g^s(\lambda) W^\psi \right\}. \]
where for each $\tilde{W}^p$ we need to write $W^p$ in terms of $W^\psi$ so that the above equation holds. In components,

$$\tilde{W}^p = W^p + g^p(\lambda) W^\psi, \quad \lambda h^\psi(\lambda) \tilde{W}^p = W^\psi + \lambda g^\psi(\lambda) W^p,$$

which implies

$$\lambda h^\psi(\lambda) \left( W^p + g^p(\lambda) W^\psi \right) = W^\psi + \lambda g^\psi(\lambda) W^p.$$

Rearranging the terms in this equation, we find

$$\lambda \left( h^\psi(\lambda) - g^\psi(\lambda) \right) W^p = \left( 1 - \lambda h^\psi(\lambda) g^p(\lambda) \right) W^\psi. \quad (3.14)$$

We need to consider $M^\psi(\lambda)$:

$$M^\psi(\lambda) := h^\psi(\lambda) - g^\psi(\lambda) : E^p_0 \to E^\psi_0.$$ 

For the moment, assume that $M^\psi(0)$ is invertible, with

$$M^p(0) := M^\psi(0)^{-1} : E^\psi_0 \to E^p_0, \quad (3.15)$$

where $(M^\psi(0))_{ij}$ is given in (3.16). We then know that $M^\psi(\lambda)$ is invertible for $\lambda$ sufficiently near zero, so we can write $(M^\psi(\lambda))^{-1} = M^p(0) + \lambda \tilde{M}^p(\lambda)$, where $M^p(\lambda)$ is analytic in $\lambda$. Hence, (3.14) can be solved via

$$W^p = \frac{1}{\lambda} \left( M^p(0) + \lambda \tilde{M}^p(\lambda) \right) \left( 1 - \lambda h^\psi(\lambda) g^p(\lambda) \right) W^\psi$$

$$= \frac{1}{\lambda} M^p(0) W^\psi + \tilde{M}^p(\lambda) (1 - \lambda h^\psi(\lambda) g^p(\lambda)) - M^p(0) h^\psi(\lambda) g^p(\lambda)] W^\psi$$

$$= \frac{1}{\lambda} M^p(0) + h_0^\psi(\lambda)) W^\psi,$$

where $h_0^\psi(\lambda)$ is analytic in $\lambda$. This defines the map $h_0^\psi(\lambda)$ from $E^\psi_0$ to $E^p_0$ as stated in the lemma. The lifted map $h^\psi(\lambda)$ from $E^\psi_0(\lambda)$ to $E^p_0(\lambda)$ is then defined by $h^\psi(\lambda) = (I + \lambda g^\psi(\lambda)) h_0^\psi(\lambda) (I + g^p(\lambda))^{-1}$, which is well defined for sufficiently small $\lambda$. The property of $h^\psi(\lambda)$ in the lemma follows directly from the above construction.

It only remains to justify (3.15). We have that $M^\psi(0) = h^\psi(0) - g^\psi(0)$, where $h^\psi$ and $g^\psi$ represent the graphs of $E^\psi_0$ and $E^p_0$, respectively. Following the same argument as in the proof of [BSZ10, lemma 6], we find that

$$\left \langle \Psi_j(0), (h^\psi(0) - g^\psi(0)) W_j(0) \right \rangle = \int_\mathbb{R} \left \langle \psi_j(y), \begin{pmatrix} 0 & 0 \\ 0 & D^{-1}_2 \end{pmatrix} W_j(y) \right \rangle dy.$$ 

Recall that the first component of $W_j(y)$ corresponds to the eigenfunctions $V_j$ defined in (2.7), and $(D^{-1}_2)^T \Psi_j$ corresponds to $\psi_j$. Thus, we find that

$$(M^\psi(0))_{i,j} = \int_\mathbb{R} \left \langle \psi_i(y), V_j(y) \right \rangle dy. \quad (3.16)$$

It now follows from [SS04a, corollary 4.6] in conjunction with hypothesis 2.3 that $M^\psi(0)$ is invertible: as shown there, invertibility of $M^\psi(0)$ encodes precisely the property that the linearization about the source in an appropriately weighted space has zero as an eigenvalue with multiplicity two.

Following [BHSZ10, BSZ10], the meromorphic extension of the exponential dichotomy for $x > y$ is then given by

$$\Phi(x, y, \lambda) := \begin{cases} 
\Phi^\psi_+ (x, y, \lambda) & 0 \leq y < x, \\
\Phi^\psi_+ (x, 0, \lambda) \tilde{\Phi}^p (0, y, \lambda) & y < 0 \leq x, \\
\Phi^\psi_+ (x, y, \lambda) & y < x < 0.
\end{cases} \quad (3.17)$$
with

\[ \Phi^+_s(x, y, \lambda) := \Phi^+_s(x, y, \lambda) - \Phi^+_s(x, 0, \lambda) h^+(\lambda) \Phi^+_s(0, y, \lambda) \quad 0 \leq y \leq x, \]

\[ \Phi^1(x, y, \lambda) := \Phi^1(x, y, \lambda) - \Phi^1(x, 0, \lambda) h^- (\lambda) \Phi^1(0, y, \lambda) \quad y \leq 0, \]

and

\[ \Phi^+_s(x, y, \lambda) = e^{\nu_j(\lambda)(x-y)} W^+_s(x, \lambda) (\Psi^+_j(y, \lambda), \cdot) + O(e^{-\eta|x-y|}) \quad 0 \leq y \leq x, \]

\[ \Phi^1(x, y, \lambda) = O(e^{-\eta|x-y|}) \quad y \leq x \leq 0, \]  

(3.18)

We are now ready to derive pointwise bounds on the resolvent kernel for \( \lambda \in B_0(0) \).

**Proof of proposition 3.3.** We give the details only for \( y \leq x \); the analysis for \( x \leq y \) is similar. Recall that the resolvent kernel \( G \) is just the upper-left two-by-two block of \( \Phi \).

**Case I:** \( 0 \leq y \leq x \). Due to (3.12) and (3.17), we have

\[ \Phi(x, y, \lambda) = -\Phi^+_s(x, 0, \lambda) h^+_s(\lambda) \Phi^+_s(0, y, \lambda) + e^{\nu_j(\lambda)(x-y)} W^+_s(x, \lambda) (\Psi^+_j(y, \lambda), \cdot) + O(e^{-\eta|x-y|}). \]

By definition (3.13), we can expand the first term as

\[ \Phi^+_s(x, 0, \lambda) h^+_s(\lambda) \Phi^+_s(0, y, \lambda) = \Phi^+_s(x, 0, \lambda) \sum_{j=1}^2 \frac{1}{\lambda} W_j(0) (\Psi_j(y), \Phi^+_s(0, y, \lambda)) \]

\[ = \frac{1}{\lambda} \sum_{j=1}^2 \Phi^+_s(x, 0, \lambda) W_j(0) (\Phi^+_s(0, y, \lambda)^* \Psi_j(y), \cdot) \]

By definition (3.11), we have

\[ \Phi^+_s(x, 0, \lambda) W_s(0) = e^{\nu_j(\lambda)x} W_s(x)(1 + O(\lambda)) + O(\lambda e^{\nu_j(\lambda)x}) = e^{\nu_j(\lambda)x} W_s(x) + O(\lambda e^{\nu_j(\lambda)x}) \]

for any \( W_s \) that is a linear combination of \( W_1 \) and \( W_2 \), where we have used the fact that \( (\Phi^+_s(0, 0), W_j(0)) = 0 \). Similarly, (3.18) implies that we also have \( \Phi^+_s(0, y, \lambda)^* \Psi_j(y) = \Psi_j(y) + O(\lambda e^{-\eta|y|}) \). Thus, we have

\[ \Phi^+_s(x, 0, \lambda) h^+_s(\lambda) \Phi^+_s(0, y, \lambda) = \frac{1}{\lambda} \sum_{j=1}^2 e^{\nu_j(\lambda)x} W_j(x) (\Psi_j(y), \cdot) + O(e^{\nu_j(\lambda)x} e^{-\eta|y|}). \]

We can also write

\[ e^{\nu_j(\lambda)(x-y)} W^+_s(x, \lambda) (\Psi^+_j(y, \lambda), \cdot) = e^{\nu_j(\lambda)(x-y)} W_s(x) (\Psi_j(y), \cdot) + O(\lambda e^{\nu_j(\lambda)(x-y)}), \]

where \( W_s(x) = W^+_s(x, 0) \) and \( \Psi_j(y) = \Psi^+_j(y, 0) \).

Thus, we have obtained

\[ G(x, y, \lambda) = \frac{1}{\lambda} \sum_{j=1}^2 e^{\nu_j(\lambda)x} V_j(x) (\Psi_j(y), \cdot) + e^{\nu_j(\lambda)(x-y)} V_s(x) (\Psi_j(y), \cdot) \]

\[ + O((\lambda + e^{-\eta|y|}) e^{\nu_j(\lambda)(x-y)}) + O(e^{-\eta|x-y|}), \]

(3.19)

where \( V_s(x) \) is a linear combination of \( V_1(x) \) and \( V_2(x) \). This proves proposition 3.3 in this case.

**Case II:** \( y \leq x \leq 0 \). A a result of (3.17), we obtain

\[ \Phi(x, y, \lambda) = \Phi^u(x, 0, \lambda) h^- (\lambda) \Phi^1(0, y, \lambda) + O(e^{\nu_j(\lambda)x} e^{-\eta|y|}) + O(e^{-\eta|x-y|}) \]
Remark 3.5.\\n
In this section, we derive pointwise bounds on the temporal Green’s function defined by

\[ G(x, y, \lambda) = \frac{1}{\lambda} \sum_{j=1}^{2} e^{\nu c(\lambda)x} W_j(x) (\Psi_j(y), \cdot) + O(e^{\nu c(\lambda)x} e^{-|y|} + e^{-|x-y|}), \]

which yields the proposition for \( y \leq x \leq 0. \)

Case III: \( y \leq 0 \leq x. \) Again by (3.17), we write

\[ \Phi(x, y, \lambda) = \Phi^a(x, 0, \lambda) \Phi^b(0, 0, \lambda) \Phi^c(0, y, \lambda) + O(e^{\nu c(\lambda)x} e^{-|y|}) \]

\[ = \frac{1}{\lambda} \sum_{j=1}^{2} \Phi^a(x, 0, \lambda) \Phi^b(0, 0, \lambda) \Phi^c(0, y, \lambda) W_j(y) (\Psi_j(y), \cdot) + O(e^{\nu c(\lambda)x} e^{-|y|}) \]

\[ = \frac{1}{\lambda} \sum_{j=1}^{2} e^{\nu c(\lambda)x} W_j(x) (\Psi_j(y), \cdot) + O(e^{\nu c(\lambda)x} e^{-|y|}). \]

Combining all the three cases yields proposition 3.3. \( \square \)

Remark 3.5. If we denote by \( G_R(x, y, \lambda) \) the first row in the matrix \( G(x, y, \lambda) \), then \( G_R(x, y, \lambda) \) has a better bound than that of \( G(x, y, \lambda) \) by \( O(e^{-|y|}) \). This is due to the structure of \( V_1(x) = (r_x, \phi_x) \) and \( V_2(x) = (0, r_x) \), since \( r_x = O(e^{-|y|}) \). More precisely, we have

\[ G_R(x, y, \lambda) = -\frac{r_x}{\lambda} e^{\nu c(\lambda)x} (\psi_1(y), \cdot) + O((\lambda + e^{-|y|}) e^{\nu c(\lambda)(x-y)}) + O(e^{-|x-y|}). \]

This last estimate means that the \( R \)-component of Green’s function will decay faster than the \( \phi \)-component by a factor of \( \lambda^{-1/2} \), after subtracting the terms resulting from the eigenfunctions.

4. Temporal Green’s function

In this section, we derive pointwise bounds on the temporal Green’s function defined by

\[ G(x, y, t) := \frac{1}{2\pi i} \int_{\gamma} e^{i\lambda t} G(x, y, \lambda) d\lambda, \]  

(4.1)

where \( \Gamma \) is a contour outside of the essential spectrum and \( G(x, y, \lambda) \) is the resolvent kernel constructed in section 3. Since \( G(x, y, \lambda) \) is analytic in \( \lambda \) outside of the essential spectrum, the contour \( \Gamma \) can be taken to be the boundary of the set \( \Omega_0 \) which is defined as in (3.2); see figure 2.

Lemma 4.1. If \( G(x, y, t) \) is defined in (4.1), then \( \tilde{G}(x, y, t) \) is the temporal Green’s function of \( \partial_t - \mathcal{L} \), where \( \mathcal{L} \) is the linearized operator defined as in (2.5). In particular, the solution to the inhomogeneous system

\[ (\partial_t - \mathcal{L}) U(x, t) = f(x, t) \]

is given by the standard Duhamel formula

\[ U(x, t) = \int_{\mathbb{R}} G(x, y, t) U(y, 0) dy + \int_{0}^{t} \int_{\mathbb{R}} \tilde{G}(x, y, t-s) f(y, s) ds dx, \]

as long as the integrals on the right-hand side are well defined.
Proof. By propositions 3.1 and 3.2, \(G(x, y, \lambda)\) is uniformly bounded outside of the essential spectrum, and therefore \(e^{itG(x, y, \lambda)}\) is integrable by moving the contour \(\Gamma\) so that \(\Gamma = \partial \Omega_\delta\) for \(\lambda\) large. Hence, \(\hat{G}(x, y, t)\) is well defined for \(x \neq y\) and for \(t > 0\). A direct calculation then yields
\[
(\partial_t - \mathcal{L})\hat{G}(x, y, t) = \delta(x - y)\delta(t).
\]
This verifies that \(\hat{G}(x, y, t)\) is indeed Green’s function.

The main result of this section is the following Proposition, which contains pointwise bounds on Green’s function. This proposition will be proven in sections 4.1 and 4.2.

Proposition 4.2. Green’s function \(\hat{G}(x, y, t)\) defined as in (4.1) may be decomposed as
\[
\hat{G}(x, y, t) = e(x, t)V_1(x)(\psi_1(y), \cdot) + e(x, t)V_2(x)(\psi_2(y), \cdot) + \hat{G}(x, y, t),
\]
where \(V_{1,2}(x)\) are defined in (2.7), the adjoint functions satisfy \(\psi_{1,2}(y) = O(e^{-\eta|y|})\), \(e(x, t)\) is the sum of error functions defined in (1.5), and \(\hat{G}(x, y, t)\) satisfies
\[
\hat{G}(x, y, t) = \sum_{\pm} O\left(t^{-\frac{1}{2}}e^{-\frac{(x-y)^2}{M^2}}\right) + O\left(e^{-\eta|y|+e^{-\eta|y|+\eta^2t}}\right).
\]
In addition, the first row of \(\hat{G}(x, y, t)\), denoted by \(\hat{G}_R(x, y, t)\), satisfies the better bound
\[
\hat{G}_R(x, y, t) = \sum_{\pm} O\left(t^{-\frac{1}{2}}(t^{-\frac{1}{2}} + e^{-\eta|y|})e^{-\frac{(x-y)^2}{M^2}}\right) + O\left(e^{-\eta|y|+e^{-\eta|y|+\eta^2t}}\right).
\]

This proof is essentially contained in [ZH98]. However, for completeness, we outline the key steps. The main idea is to deform the contour \(\Gamma\) in such a way as to minimize the integral in the definition of \(\hat{G}\). The choice of the minimizing contour is based on the method known as the saddle point method, the method of stationary phase, or the method of steepest descents. As we are only integrating with respect to the spectral parameter \(\lambda\), we are free to choose the contour to depend on \((x, y, t)\).

We first split the contour \(\Gamma\) into \(\Gamma_1 \cup \Gamma_2\), with
\[
\Gamma_1 := \partial B(0, M) \cap \overline{\Omega}_\delta, \quad \Gamma_2 := \partial \Omega_\delta \setminus B(0, M),
\]
for some appropriate \(M\) to be determined below (see figure 4). We consider two cases: when \(|x - y|/t\) is sufficiently large and when it is bounded.

4.1. Large \(|x - y|/t\)

Set \(z := \eta^2|x - y|/(2t)\) and \(M := (z/\eta)^2\). Since \(M\) is sufficiently large, we can apply proposition 3.1 so that
\[
|G(x, y, \lambda)| \leq C|\lambda|^{-1/2}e^{-|\eta^2x - y|}/(|x - y|).
\]
Let \(\lambda_0\) and \(\lambda_0^\ast\) be the two points of intersection of \(\Gamma_1\) and \(\Gamma_2\). We observe that, after possibly making \(\eta\) smaller, for all \(\lambda \in \Gamma_1\) we have
\[
\text{Re} \lambda = M \cos(\text{arg} \lambda) \leq M(1 - \eta \arg(\lambda)^2).
\]
Also, on \(\Gamma_2\) we have \(\text{Re} \lambda = -\vartheta(1 + |\text{Im} \lambda|)\) (choosing \(\vartheta = \min\{\vartheta_1, \vartheta_2\}\) and deforming the contour slightly) and so
\[
\text{Re} \lambda = \text{Re} \lambda_0 - \vartheta(|\text{Im} \lambda| - |\text{Im} \lambda_0|), \quad \lambda \in \Gamma_2.
\]
Thus, we can compute
\[
\left| \int_{\Gamma_1} e^{i\lambda t} G(x, y, \lambda) d\lambda \right| \leq C \int_{\Gamma_1} |\lambda|^{-1/2} e^{Re\lambda t - \eta|\lambda|^{1/2}|x-y|} |d\lambda|
\]
\[
\leq C e^{-\frac{z}{\eta}} \int_{-arg(\lambda_0)}^{arg(\lambda_0)} M^{-1/2} e^{-\eta M t \arg(\lambda)^2} d(\arg(\lambda))
\]
\[
\leq C t^{-1/2} e^{-\frac{z}{\eta}},
\]
and
\[
\left| \int_{\Gamma_2} e^{i\lambda t} G(x, y, \lambda) d\lambda \right| \leq C \int_{\Gamma_2} |\lambda|^{-1/2} e^{Re\lambda t - \eta|\lambda|^{1/2}|x-y|} |d\lambda|
\]
\[
\leq C e^{Re(\lambda_0)t - \eta|\lambda_0|^{1/2}|x-y|} \int_{\Gamma_2} |\lambda|^{-1/2} e^{Re\lambda t} |d\lambda|
\]
\[
\leq C e^{-\frac{z}{\eta}} \int_{\Gamma_2} |\lambda|^{-1/2} e^{-\theta |\Im\lambda| - |\Im\lambda_0| t} |d\lambda|
\]
\[
\leq C t^{-1/2} e^{-\frac{z}{\eta}}.
\]
Combining these last two estimates and noting that \(z = \frac{2|x-y|}{\eta}\) is large and \(\eta\) is small, we have
\[
|\Omega(x, y, t)| \leq Ct^{-1/2} e^{-\frac{z}{\eta}} = Ct^{-1/2} e^{-\frac{z}{\eta}} e^{-\frac{z}{\eta}} \leq Ct^{-1/2} e^{-\frac{z}{\eta}} e^{-\frac{z^2}{2|\lambda_0| t}} \leq C t^{-1/2} e^{-\frac{z}{\eta}} e^{-\frac{z^2}{2|\lambda_0| t}}
\]
for \(\eta > 0\) independent of the amplitude of \(\frac{|x-y|}{t}\), as long as it is sufficiently large.

4.2. Bounded \(|x-y|/t\)

We now turn to the critical case where \(|x-y|/t \leq S\) for some fixed \(S\). We first deform \(\Gamma\) into \(\Gamma_1 \cup \Gamma_2\) as in (4.2) with \(M\) now being sufficiently small.

The integral over \(\Gamma_2\) is relatively straightforward. Again let \(\lambda_0, \lambda_0^*\) be the points where \(\Gamma_1\) meets \(\Gamma_2\). We have \(Re \lambda_0 = -\bar{\vartheta}(1 + \tilde{\vartheta})\) and \(Im \lambda_0 = \tilde{\vartheta}\) for some \(\tilde{\vartheta} > 0\). Moreover, on \(\Gamma_2\) we
have \(|G(x, y, \lambda)| \leq C|\lambda|^{-\frac{1}{2}}\) by the mid- and high-frequency resolvent bounds. Thus, we can estimate
\[
\left| \int_{\Gamma_2} e^{itG(x, y, \lambda)} d\lambda \right| \leq C |\text{Re}\lambda|^{\frac{1}{2}} e^{-\eta |\text{Im}\lambda|^{1/2}} |d\text{Im}\lambda| \leq Ct^{-\frac{1}{4}} e^{-\delta t}.
\]
Noting that \(|x - y| \leq St\), we have
\[
\left| \int_{\Gamma_2} e^{itG(x, y, \lambda)} d\lambda \right| \leq Ct^{-\frac{1}{4}} e^{-\eta |x-y|^\sigma}
\]
for some small \(\eta > 0\).

In order to estimate the integral over \(\Gamma_1\), we take \(M\) and \(\theta\) small enough so that \(\Gamma_2\) remains outside of the essential spectrum and the \(\lambda\)-expansion obtained in section 3.3 for the low-frequency resolvent kernel \(G(x, y, \lambda)\) holds. We consider several cases, depending on the position of \(x\) and \(y\).

**Case I:** \(0 \leq y \leq x\). In this case, we have expansion (3.19), and by equation (3.8) we also have
\[
v_c^e(\lambda) = -\frac{\lambda}{c_g} + \frac{d\lambda^2}{c_g^2} + O(\lambda^3),
\]
where \(d, c_g > 0\) are defined as in (1.3) and (2.2), respectively. We first estimate the contribution of the term \(O(\lambda^2 e^{\nu_c^e(\lambda)(x-y)})\). Set
\[
\text{where } d := \frac{x - y - c_g t}{2t}, \quad z_2 := \frac{d(x - y)}{c_g^2 t}.
\]
Then \(\lambda = z_1/z_2\) minimizes \(v_c^e(\lambda)(x-y)\) when \(\lambda\) is real. Define \(\Gamma_1\) to be the portion contained in \(\Omega_r := \{|\lambda| \leq r\}\) of the hyperbola
\[
-\frac{1}{c_g^2} \text{Re}(\lambda - \lambda^2 d/c_g^2) \equiv -\frac{1}{c_g^2} (\lambda_{min} - \lambda^2 d/c_g^2)
\]
where \(\lambda_{min}\) is defined by \(z_1/z_2\) if \(|z_1/z_2| \leq \epsilon\) and by \(\pm \epsilon\) if \(z_1/z_2 \geq \epsilon\), for \(\epsilon\) small.

With these definitions, we readily obtain that
\[
\text{Re}(\lambda t + v_c^e(\lambda)(x-y)) \leq -\frac{1}{c_g^2} (z_1^2 t/2z_2) - \eta |\text{Im}\lambda|^2 t \leq -z_1^2 t/C_0 - \eta |\text{Im}\lambda|^2 t,
\]
for \(\lambda \in \Gamma_1\); here, we note that \(z_2\) is bounded above, and we have used the crucial fact that \(z_1\) controls \((|x| + |y|)/t\), in bounding the error term \(O(\lambda^3(|x| + |y|)/t)\) arising from expansion. Thus, we obtain for any \(\eta\) that
\[
\left| \int_{\Gamma_1} e^{\text{Re}(\lambda t + v_c^e(\lambda)(x-y))} d\lambda \right| \leq C e^{-z_1^2 t/C_0} \int_{\Gamma_1} (|\lambda_{min}|^2 + |\text{Im}\lambda|^2) e^{-\eta |\text{Im}\lambda|^2 t} d\lambda \\
\leq Ct^{-\frac{1}{4}} e^{-\eta t/C_0},
\]
for suitably large \(C, M_0 > 0\).

Thus, the contribution from \(O(\lambda^2 e^{\nu_c^e(\lambda)(x-y)})\) to the Green function bounds is
\[
O(t^{-\frac{1}{4}} e^{-\frac{(|x-y|^2)}{C_0}} + O(e^{-\eta |(x-y)|^\sigma})).
\]
Clearly, the term \(O(e^{-\eta |(x-y)|^\sigma})\) in \(\Phi^e(x, y, \lambda)\) contributes a time- and space-exponential decay: \(O(e^{-\eta |(x-y)|^\sigma})\). Thus, we are left with the term involving \(\lambda^{-1}\). Precisely, consider the term
\[
-\frac{1}{\lambda} \sum_{j=1}^{2} e^{-(\lambda/c_g^2 + d\lambda^2/c_g^4) t} V_j(x).
\]
Figure 5. The contour $\Gamma_1$ is deformed into three straight lines.

We denote

$$\alpha(x, t) := \frac{1}{2\pi i} \int_{\Gamma_1 \cap \{ |\lambda| \leq r \}} \lambda^{-1} e^{\lambda t} e^{-\lambda/c_g + \frac{1}{c_g^2} x} d\lambda.$$ 

Thus, by Cauchy’s theorem, we can move the contour $\Gamma_1$ (as shown in figure 5) to obtain

$$\alpha(x, t) = \frac{1}{2\pi} \text{P.V.} \int_{-\infty}^{\infty} (i\xi)^{-1} e^{i\xi (t-x/c_g)} e^{-\xi^2 (d/c_g^2)} d\xi$$

$$+ \frac{1}{2\pi i} \left( \int_{-\eta-i\infty}^{-\eta+i\infty} + \int_{i\eta}^{i\infty} \right) \lambda^{-1} e^{\lambda t} e^{(-\lambda/c_g + \frac{1}{c_g^2} x)} d\lambda$$

$$+ \frac{1}{2} \text{Residue}_{\lambda=0} e^{\lambda t} e^{(-\lambda/c_g + \frac{1}{c_g^2} x)}.$$ 

for some $\eta > 0$. Rearranging and evaluating the residue term, one then has

$$\alpha(x, t) = \left( \frac{1}{2\pi} \text{P.V.} \int_{-\infty}^{\infty} (i\xi)^{-1} e^{i\xi (t-x/c_g)} e^{-\xi^2 (d/c_g^2)} d\xi + \frac{1}{2} \right)$$

$$- \frac{1}{2\pi i} \left( \int_{-\eta-i\infty}^{-\eta+i\infty} + \int_{i\eta}^{i\infty} \right) (i\xi)^{-1} e^{i\xi (t-x/c_g)} e^{-\xi^2 (d/c_g^2)} d\xi$$

$$+ \frac{1}{2\pi i} \left( \int_{-\eta-i\infty}^{-\eta+i\infty} + \int_{i\eta}^{i\infty} \right) \lambda^{-1} e^{\lambda t} e^{(-\lambda/c_g + \frac{1}{c_g^2} x)} d\lambda.$$ 

Note that the first term in $\alpha(x, t)$ can be explicitly evaluated, again by Cauchy’s theorem and the standard dominated convergence theorem, as

$$\frac{1}{2} - \exp \left( -\frac{(x-c_g t)^2}{4dx/c_g} \right) \lim_{r \to \infty} \frac{1}{2\pi} \int_0^{\pi} e^{-\frac{1}{2} \frac{(r-c_g \theta)^2}{d/c_g^2}} d\theta = \frac{1}{2} \left( 1 - \exp \left( -\frac{(x-c_g t)^2}{4dx/c_g} \right) \right),$$

which is conveniently simplified to

$$\frac{1}{2} \text{erf} \left( \frac{x + c_g t}{\sqrt{4dx/c_g}} \right),$$

(4.3)
plus a time-exponentially small error. The second and third terms are clearly bounded by $Ce^{-\eta|x|}$ for $\eta$ sufficiently small relative to $t$, and thus time-exponentially small for $t \leq C|\lambda|$. In the case $t \geq C|\lambda|$, $C > 0$ sufficiently large, we can simply move the contour to $[-\eta-ir, -\eta+ir]$ to obtain (complete) residue 1 plus a time-exponentially small error corresponding to the shifted contour integral. In this case, we note that the result again can be expressed as the errfn (4.3) plus a time-exponentially small error. Indeed, for $t \geq C|\lambda|$, large, one can estimate

$$1 - \text{errfn} \left( \frac{-x + cg t}{\sqrt{4d|x/c_g|}} \right) \sim 1 - \text{errfn} \left( \sqrt{t} \right) = \frac{1}{2\pi} \int_{\sqrt{t}}^{+\infty} e^{-x^2} \, dx = O(e^{-nt}).$$

Thus, we have obtained for $0 \leq y \leq x$:

$$G(x, y, t) = \sum_{j=1}^{2} V_j(x)(W_j(y), \cdot)\text{errfn} \left( \frac{-x + cg t}{\sqrt{4d|x/c_g|}} \right) + O(t^{-\frac{1}{2}}e^{-\frac{|x-y|^2}{4t}}) V_j(x)$$

$$+ O(t^{-\frac{1}{2}}(t^{-\frac{1}{2}} + e^{-\eta|y|})e^{-\frac{|x-y|^2}{4t}}) + O(e^{-\eta|x-y|t})) \quad (4.4)$$

in which we recall that $V_j(x)$ belongs to the span of $V_1(x)$ and $V_2(x)$. We note that the errfn in (4.4) may be rewritten as

$$\text{errfn} \left( \frac{-x + cg t}{\sqrt{4d|x/c_g|}} \right)$$

plus error

$$\text{errfn} \left( \frac{-x + cg t}{\sqrt{4d|x/c_g|}} \right) - \text{errfn} \left( \frac{-x + cg t}{\sqrt{4d|x/c_g|}} \right) = O(t^{-1}e^{(x-cg t^2)/M^2}),$$

for $M > 0$ sufficiently large. Note also that errfn $(\frac{-x+cgt}{\sqrt{4dt}})$ is time-exponentially small since $x, t, c_g$ are all positive. Thus, an equivalent expression for $G(x, y, t)$ is

$$G(x, y, t) = \sum_{j=1}^{2} V_j(x)\left[ e_+(x, t)(W_j(y), \cdot) + O(t^{-\frac{1}{2}}e^{-\frac{|x-y|^2}{4t}}) \right]$$

$$+ O(t^{-\frac{1}{2}}(t^{-\frac{1}{2}} + e^{-\eta|y|})e^{-\frac{|x-y|^2}{4t}}) + O(e^{-\eta|x-y|t})),$$

with $e_\pm(x, t) = \left[ \text{errfn} \left( \frac{-x+cgt}{\sqrt{4dt}} \right) - \text{errfn} \left( \frac{-x+cgt}{\sqrt{4dt}} \right) \right]$. Again, we note that at the leading term the first row (or the $R$-component) of the Green function matrix $G(x, y, t)$ enjoys a better bound due to the special structure of $V_j(x)$; see (2.7).

**Case II:** $y \leq x \leq 0$. In this case we recall that

$$G(x, y, \lambda) = O(e^{\nu_+(\lambda)x}e^{-\eta|y|}) + O(e^{-\eta|x-y|}) = \frac{1}{\lambda} \sum_{j=1}^{2} \langle W_j(y), \cdot \rangle e^{\nu_+(\lambda)x} V_j(x),$$

where

$$\nu_+(\lambda) = -\frac{\lambda}{c_g} + \frac{d\lambda^2}{(c_g)^2} + O(\lambda^3),$$

with $c_g^- = -c_g < 0$. Similar computations as in the previous case yield

$$G(x, y, t) = \sum_{j=1}^{2} V_j(x)e_-(x, t)(W_j(y), \cdot) + O(t^{-\frac{1}{2}}e^{-\eta|y|}e^{-\frac{|x-y|^2}{4t}}) + O(e^{-\eta|x-y|t}))$$

with $e_\pm(x, t) = \left[ \text{errfn} \left( \frac{-x+cgt}{\sqrt{4dt}} \right) - \text{errfn} \left( \frac{-x+cgt}{\sqrt{4dt}} \right) \right].$
Case III: \( y \leq 0 \leq x \). In this case we have
\[
G(x, y, \lambda) = O(e^{-\eta x}e^{-|y|}) + O(e^{\eta x}e^{-|y|}) - \frac{1}{2} \sum_{j=1}^{2} (W_j(y), \cdot) e^{\eta x} V_j(x).
\]
Thus, similar computations as done in case I yield
\[
G(x, y, t) = 2 \sum_{j=1}^{2} V_j(x) e^{\nu c} (\cdot) \langle W_j(y), \cdot \rangle + O(t^{-1/2} e^{-\eta|x+y|} e^{-(x-cg t)2Mt}) + O(e^{-\eta(|x+y|t)}).
\]
This establishes proposition 4.2.

5. Asymptotic ansatz

In this section, we shall construct the asymptotic ansatz that will be used in the analysis of (1.1). For convenience, we recall that the qCGL equation (1.1) (after rescaling so that \( \mu = 1 \)) is
\[
A_t = (1+i\alpha)A_{xx} + A - (1+i\beta)|A|^2 + (\gamma_1 + i\gamma_2)|A|^4.
\]

5.1. Setup

Recall that the source solution of (5.1) is given by
\[
A_{\text{source}}(x, t) = r(x) e^{i\phi(x)} e^{-i\omega_0 t}.
\]
Given a solution \( A(x, t) \) of (5.1), define
\[
B(x, t) := A(x+p(x, t), t) = [r(x) + R(x, t)] e^{i(\phi(x) + \phi(x,t))} e^{-i\omega_0 t},
\]
where for the moment \( p(x, t) \) is an arbitrary smooth function. The functions \( R \) and \( \phi \) denote the perturbation in the amplitude and phase variables. We will insert (5.2) into (5.1) and derive a system for \( U = (R, r\phi) \). We obtain the following simple lemma whose proof will be given in section 5.3.

Lemma 5.1. Let \( p(x, t) \) be a given smooth function so that the map \( (x, t) \mapsto (\xi(x, t), \tau(x, t)) \), defined by \( \xi = x + p(x, t) \) and \( \tau = t \), is invertible. If \( A(x, t) \) is a solution of (5.1), then \( U = (R, r\phi) \), where \( R \) and \( \phi \) are defined by (5.2), satisfies
\[
(\partial_t - \mathcal{L}) U - T(p) - Q(R, \phi, p) = \mathcal{N}(R, \phi, p).
\]
In this equation, \( \mathcal{L} \) is the linearized operator, \( T(p) \) denotes the linear residual effects resulting from the function \( p \), \( Q(R, \phi, p) \) denotes the quadratic nonlinear terms, and \( \mathcal{N}(R, \phi, p) \) denotes the higher order terms. They are defined in equations (5.4), (5.6), (5.7), and (5.8), respectively.

We recall from lemma 2.2 that
\[
\mathcal{L} = D_2 \partial_t^2 - 2\varphi_x D_1 \partial_t + D_0^\phi(x) + D_0(x),
\]
with
\[
D_1 := \begin{pmatrix} \alpha & 1 \\ -1 & \alpha \end{pmatrix}, \quad D_2 := \begin{pmatrix} 1 & -\alpha \\ \alpha & 1 \end{pmatrix}
\]
\[
D_0^\phi(x) := \begin{pmatrix} 0 & (2\varphi_x r_x + \alpha r_{xx}) \\ 0 & (2\alpha \varphi_x r_x - r_{xx}) \end{pmatrix}, \quad D_0(x) := \begin{pmatrix} 1 - 3r^2 - \alpha \varphi_{xx} - \varphi_x^2 + 5\gamma_1 r^4 & 0 \\ \omega_0 - 3\beta r^2 - \alpha \varphi_x^2 + \varphi_{xx} + 5\gamma_2 r^4 & 0 \end{pmatrix}.
\]
The function $T(p)$ is defined by
\[
T(p) := p_t \left( \frac{r_x}{r \phi_x} \right) + 2r \psi_x^2 p_x \left( \frac{1}{\alpha} \right) + r \phi_x p_{xx} \left( \frac{\alpha}{-1} \right),
\]
and the quadratic term $Q(R, \phi, p) := (Q_R(R, \phi, p), Q_\phi(R, \phi, p))$ is defined by
\[
Q_R(R, \phi) := -2\psi_x R\phi_x - r\phi_x^2 - 3r R^2 + 10\gamma_1 r^3 R^2 - 3\beta r R^2 + 10\gamma_2 r^3 R^2 - r\phi_x p_x + r\phi_x p_x R - 3\alpha\rho^2 p_x^2 + 4\alpha r\phi_x p_x \phi_x + 2\alpha\rho^2 p_x R.
\]
The remainder $N(R, \phi, p)$, corresponding to higher order nonlinear terms as well as terms that are small due to exponential localization in space, satisfies the uniform bound
\[
N(R, \phi, p) = O \left[ |R\phi_x^2| + |R_x\phi_x| + |R\phi_{xx}| + |R|R_x| + |p_x|^3 + |\rho|^2 \right]
+ |p_x p_{xx}| + |p_t R_x| + |\phi_{xx}||(|R| + |\phi_x|) + |p_x||(|R_x| + |\phi_{xx}| + |R\phi_x| + e^{-\eta_0|x|}).
\]

5.2. Approximate solution

We will construct an approximate solution $U^a$ that solves (5.3) up to some error terms that are harmless for the nonlinear analysis. In the course of constructing the approximate ansatz, the Burgers-type equation (1.12) appears as a governing equation at leading order. Let
\[
L_B := \partial_t + 2(\alpha - \beta_x)\phi_x \partial_x - d\delta_x^2
\]
denote the linear part of the Burgers equation, where we recall that
\[
d := (1 + \alpha \beta_x) - \frac{2k_0^2(1 + \beta_x^2)}{r_0^2(1 - 2\gamma_1 r_0^2)}, \quad \beta_x := \frac{\beta - 2\gamma_2 r_0^2}{1 - 2\gamma_1 r_0^2}.
\]
Note that the operator $L_B$ corresponds precisely with the curve essential spectrum of the linearized operator $\mathcal{L}$ that touches the origin; see equation (2.8).

Recall from the discussion of section 1 that the dynamics of $R$ are higher order, relative to $\phi$, due to the zero-order terms in the $R$-component, which can be seen in the operator $D_0$ defined in (5.5). In order to construct the approximate solution, we will (asymptotically) diagonalize the operator $D_0$ so that these zero-order terms only appear in the equation for $R$. Variables labeled with a hat correspond to variables to which this diagonalizing transformation has been applied. We first state the approximate solution in these variables, and then undo the diagonalizing transformation.

Our approximate solution is constructed as follows. First, take $\mathcal{B}(x, t)$ to be an arbitrary smooth solution to the linear Burgers equation: $L_B \mathcal{B} = 0$. Define
\[
\begin{align*}
\phi^\delta(x, t, \delta^\pm) := & \frac{d}{2q} \log \left( 1 + \delta^+ \mathcal{B}(x, t) \right) + \log \left( 1 + \delta^- \mathcal{B}(x, t) \right), \\
p(x, t, \delta^\pm) := & \frac{d}{2q k_0} \log \left( 1 + \delta^+ \mathcal{B}(x, t) \right) - \log \left( 1 + \delta^- \mathcal{B}(x, t) \right),
\end{align*}
\]
where
\[
q := (\alpha - \beta_x) + \frac{4k_0^2(\gamma_1 \beta - \gamma_2)}{(1 - 2\gamma_1 r_0^2)^2}
\]
and $\delta^\pm \in \mathbb{R}$ are for the moment arbitrary constants.
Remark 5.2. For each fixed $\delta^\pm$, the function $(\hat{\phi}^a \pm k_0 p)(x, t, \delta^\pm)$, defined via (5.9), solves the nonlinear Burgers-type equation

$$L_B(\hat{\phi}^a \pm k_0 p) = q(\hat{\phi}^a \pm k_0 p)^2,$$

(5.11)

where $q$ is defined in (5.10). In fact, $\hat{\phi}^a$ and $p$ were chosen essentially by applying the Cole–Hopf transformation to the function $B$. Therefore, they will exactly cancel the potentially problematic quadratic terms in the nonlinearity. Moreover, (5.11) implies that

$$\hat{\phi}^a_x + 2(\alpha - \beta_s)\psi_x \hat{\phi}^a_x = O\left(|\hat{\phi}^a_x|^2 + |p_x|^2 + |p_t|^2\right),$$

$$p_x + 2(\alpha - \beta_s)\psi_x p_x = O\left(|\hat{\phi}^a_x|^2 + |p_x|^2 + |p_t|^2\right).$$

(5.12)

For the approximate solution of the $R$-component, define

$$\hat{R}^a(x, t, \delta^\pm) := \hat{R}^a_0(x, t, \delta^\pm) + \hat{R}^a_1(x, t, \delta^\pm),$$

(5.13)

where

$$\hat{R}^a_0(x, t, \delta^\pm) := \frac{k_0}{r_0(1 - 2\gamma r_0^2)}(\hat{\phi}^a_x + \psi_x p_x)(x, t, \delta^\pm)$$

and $\hat{R}^a_1$ is a higher order correction defined in (5.29). Finally, define $\hat{U}^a := (\hat{R}^a, r\hat{\phi}^a)$.

In the original variables, the approximate solution $U^a = (R^a, r\phi^a)$ is given by

$$U^a(x, t, \delta^\pm) := S\hat{U}^a(x, t, \delta^\pm),$$

(5.14)

where $S := \begin{pmatrix} 1 & 0 \\ \beta_s & -1 \end{pmatrix}$.

We shall prove that for each fixed $\delta^\pm$, $U^a(x, t, \delta^\pm)$ solves (5.3) up to good error terms of the form

$$\gamma(\hat{\phi}^a_x, p_x) := \sum_{\psi = \phi^a, p} \left[e^{-\eta_0|\psi_x|} + |\psi_x|^2 + |\psi_t|^3\right],$$

(5.15)

for $\eta_0 > 0$ as in lemma 2.1. The main result of this section is the following proposition, whose proof will be given in section 5.4.

Proposition 5.3. Let $\delta^\pm = \delta^\pm(t)$ be arbitrary smooth functions. The function $U^a = (R^a, r\phi^a)$ defined in (5.14) solves

$$(\partial_t - L)U^a - T(p) - Q(R^a, \phi^a, p) = \sum_{\pm} \delta^\pm(t)\left(\Sigma^\pm(x, t, \delta^\pm(t)) + O(|\hat{\phi}^a_x| + |p_x|)\right)$$

+ $O(\gamma(\hat{\phi}^a_x, p_x))$.

(5.16)

Here $L$, $T(p)$ and $Q(R, \phi, p)$ are defined in (5.4), (5.6), and (5.7), respectively, $\gamma(\hat{\phi}^a_x, p_x)$ is defined in (5.15), and $\Sigma^\pm(x, t, \delta^\pm)$ are defined by

$$\Sigma^\pm(x, t, \delta^\pm) := \frac{\partial U^a}{\partial \delta^\pm}(x, t, \delta^\pm) - \frac{\partial p}{\partial \delta^\pm}(x, t, \delta^\pm) \left(r_x \psi_x\right).$$

(5.17)

We now choose the function $B(x, t)$ used in (5.9) to be

$$B(x, t) = e(x, t + 1), \quad e(x, t) = \text{erf}_n\left(\frac{x + c_d t}{\sqrt{4d \tau}}\right) - \text{erf}_n\left(\frac{x - c_d t}{\sqrt{4d \tau}}\right).$$

(5.18)

Note that, although it is not the case that $L_B B = 0$, this equation is satisfied asymptotically in an appropriate sense. In other words, the function $B$ is still sufficient for our ansatz in the sense that proposition 5.4, which will be proven in section 5.5, holds.
Proposition 5.4. Let $\delta^\pm = \delta^\pm(t)$ be arbitrary smooth functions and let $\mathcal{B}(x,t)$ defined as in (5.18). The approximate solution $U^a$ constructed as in (5.14) satisfies
\[
(\partial_t - \mathcal{L})U^a - T(p) - Q(R^a, \phi^a, p) = \frac{d}{2\gamma} \sum_{\pm} \frac{\delta^\pm(t)}{1 + \delta^\pm} \mathcal{E}^\pm(x,t) + \mathcal{R}_1^{\text{app}}(x,t, \delta^\pm),
\]
where $\mathcal{E}^\pm(x,t)$ is defined by
\[
\mathcal{E}^\pm(x,t) = \left( \frac{\phi_x \pm k_0}{r_0(1 - 2\gamma r_0^2)} B_s(x,t) \right) \mp \frac{B(x,t)}{k_0} \left( \frac{r_s}{r \phi_s} \right),
\]
and the remainder $\mathcal{R}_1^{\text{app}}(x,t, \delta^\pm)$ satisfies
\[
\mathcal{R}_1^{\text{app}}(x,t, \delta^\pm) \leq C(|\delta^+| + |\delta^-|)e^{-\eta_0|x|} \theta(x,t) + C(1 + t)^{-1} \theta(x,t)(|\delta^+|^2 + |\delta^-|^2)
\]
\[
+ C(|\delta^+| + |\delta^-|)(|\dot{\delta}^+| + |\dot{\delta}^-|)(1 + t)^{1/2} \theta(x,t),
\]
for $\eta_0 > 0$ as in lemma 2.1. Here $\theta(x,t)$ denotes the Gaussian-like behaviour (as in (1.6)):
\[
\theta(x,t) = \frac{1}{(1 + t)^{1/2}} \left( e^{-\frac{c^2 - (x - c)^2}{2M_0^2(t+1)}} + e^{-\frac{c^2 - (x + c)^2}{2M_0^2(t+1)}} \right),
\]
for some $M_0 > 0$.

The following lemma is relatively straightforward, but crucial to our analysis later on.

Lemma 5.5. Let $\mathcal{E}^\pm(x,t)$ be defined as in (5.19). Then
\[
(\partial_t - \mathcal{L})\mathcal{E}^\pm(x,t) = O((1 + t)^{-1} \theta(x,t)) + O(e^{-\eta_0|x|} \theta(x,t)).
\]

5.3. Proof of lemma 5.1

In this subsection, we shall prove lemma 5.1. First, we obtain the following simple lemma.

Lemma 5.6. Let $p(x,t)$ be a given smooth function so that the map $(x, t) \mapsto (\xi(x, t), \tau(x, t))$, defined by
\[
\xi = x + p(x,t) \quad \text{and} \quad \tau = t,
\]
is invertible. If $A(x,t)$ solves the qCGL equation (5.1), then the function
\[
B(x,t) := A(x + p(x,t), t)
\]
solves the modified qCGL equation
\[
\partial_t B = (1 + i\alpha)\partial_x B + B - (1 + i\beta)B|B|^2 + \gamma B|B|^4 + T(p, B),
\]
where $\gamma = \gamma_1 + i\gamma_2$ and
\[
T(p, B) := \frac{p_x}{1 + p_x} B_x - \frac{p_x}{1 + p_x} \left( (1 + i\alpha)B_x - \frac{p_x(2 + p_x)}{(1 + p_x)^2} (1 + i\alpha)B_{xx} \right).
\]

Proof. Write $B(x,t) = A(\xi(x,t), t)$. Then
\[
B_t = A_t \dot{p} + A_t, \quad B_x = A_x(1 + p_x), \quad B_{xx} = A_{xx}(1 + p_x)^2 + A_x p_{xx},
\]
which implies
\[
A_t = B_t - \frac{\dot{p} B_x}{1 + p_x}, \quad A_x = B_x \frac{1 + p_x}{1 + p_x}, \quad A_{xx} = B_{xx} \frac{1 + p_x}{(1 + p_x)^3} - \frac{p_{xx} B_x}{(1 + p_x)^3}.
\]
Inserting these expressions into the equation
\[
A_t = (1 + i\alpha)A_{xx} + A - (1 + i\beta)A|A|^2 + \gamma A|A|^4
\]
yields the result.
Next, we write the solution to the new qCGL equation (5.21) in the amplitude and phase variables

\[ B(x, t) = [r(x) + R(x, t)]e^{i[\psi(x) + \phi(x, t)]}e^{-i\omega t}, \]

with \((R, \phi)\) denoting the perturbation variables. As in section 2.2, but now keeping all nonlinear terms, we collect the real and imaginary parts of the equations for \(R\) and \(\phi\) to find

\[
\begin{align*}
R_t &= R_{xx} - 2\alpha r \phi R_x + (1 - 3\psi^2 - \alpha)R_{xx} - \psi_x^2 + 5\gamma r^4 R - \alpha r \phi_x - 2(\alpha r_x + r \phi_x)\phi_x + T_p(p) + Q(R, \phi, p) + N(R, \phi, p), \\
\phi_t &= r \phi_{xx} + (-2\alpha r \phi_x + 2\alpha_{x}) \phi_x + \alpha R_{xx} + 2r \phi_x R_x + (\omega_0 + \frac{\psi_x}{\phi_x} - 2\beta r^2 + 5\gamma r^4)R_t + T_p(p) + Q(R, \phi, p) + N(R, \phi, p).
\end{align*}
\]

The function \(T_j(p)\), \(j = R, \phi\), denotes the terms resulting from \(T(p, B)\) that are linear in \(p\). The function \(Q_j(R, \phi, p)\) collects terms that are quadratic in \((R, \phi, p)\), and \(N_j(R, \phi, p)\) denotes the remaining terms. We now calculate these functions in detail.

First, note that (5.22) can be written

\[
T(p, B) = p_1(1 - p_1)B_t - (1 + i\alpha)p_{xx}B_x - (1 + i\alpha)p_1(2 - 3p_1)B_{xx} + \mathcal{O}(p_1^2 + p_1^3p_t + p_1p_{xxx}),
\]

as long as \(|p_1| < 1\). Thus, by collecting the linear terms in the real and imaginary parts of \(T(p, B)\), we obtain (5.6). Note that we do not include any terms that are of the form \(p_1 e^{-\eta_0|x|}\), as these are higher order and appear in \(N(R, \phi, p)\). Similarly, we also obtain (5.7) and the bounds (5.8) for the remainder \(N(R, \phi, p)\). This completes the proof of lemma 5.1.

### 5.4. Proof of proposition 5.3

We next turn to the construction of the approximate solution to (5.3) and prove proposition 5.3. For simplicity, we first consider the case where \(\delta^2\) are constants. Define

\[
\mathcal{F}(U) := (\partial_t - \mathcal{L}^0)U - T(p, \phi, p) - Q(R, \phi, p),
\]

where \(\mathcal{L}^0\) is the approximate linear operator defined by

\[
\mathcal{L}^0 = D_x^2 - 2\phi_x D_x \partial_x + D_0^\phi(x) + D_0^\infty, \quad D_0^\infty := -2r_0^2 \begin{pmatrix} 1 - 2\gamma r_0^2 & 0 \\ \beta - 2\gamma r_0^2 & 0 \end{pmatrix}.
\]

See (5.4). By lemma 2.1, it follows that

\[
|\mathcal{L}(U)| = \mathcal{O}(e^{-\eta_0|x|})\|R\|.
\]

Our goal in this section is to construct \(U^0\) so that \(\mathcal{F}(U^0) = 0\) up to the good error term as claimed in proposition 5.3.

The zero-order term \(D_0^\infty\) is asymptotically zero. However, we keep \(D_0^\phi(x)\) in the definition of \(\mathcal{L}^0\) precisely because it may be viewed as part of the first- and second-order derivative terms when acting on \((R, r \phi)\). To see this, define

\[
\mathcal{L}_1^0 = D_x^2 - 2\phi_x D_x \partial_x + D_0^\phi(x),
\]

and let \(Z\) to be the diagonal matrix \(\text{diag}(1, r)\) so that \((R, r \phi) = Z(R, \phi)\). A direct calculation shows that

\[
\mathcal{L}_1^0 \begin{pmatrix} R \\ r \phi \end{pmatrix} = \begin{pmatrix} D_x Z D_x^2 - 2\phi_x D_x Z \partial_x \\ 0 \end{pmatrix} \begin{pmatrix} R \\ \phi \end{pmatrix} + \begin{pmatrix} 0 \\ -2\phi_x \end{pmatrix} \partial_x \begin{pmatrix} R \\ \phi \end{pmatrix}.
\]

Thus, we see that the right-hand side consists of no zero-order term. This cancellation is crucial to our analysis as it implies that the zero-order term in \(\mathcal{L}^0\) involves \(R\) only, which leads to the faster decay of \(R\), relative to \(\phi\).
To extract the governing (leading order) equations in system (5.23), we first diagonalize the zero-order constant matrix $D_0^\infty$. We can write

$$D_0^\infty = S \begin{pmatrix} -2r_0^2(1-2y_1r_0^2) & 0 \\ 0 & 0 \end{pmatrix} S^{-1} \quad \text{with} \quad S := \begin{pmatrix} 1 & 0 \\ \beta_s & -1 \end{pmatrix},$$

where we recall that $\beta_s = \frac{\beta - 2\gamma_1 r_0^2}{1-2y_1 r_0^2}$. Define the change of variables $\hat{U} = (\hat{R}, \hat{\phi}) = S^{-1}U$, and set

$$\hat{\mathcal{F}}(\hat{U}) := \left( \hat{\partial}_t - S^{-1}L_1^0 S - S^{-1}D_0^\infty S \right) \hat{U} - S^{-1} \left( T(p) + \mathcal{Q}(R, \phi, p) \right),$$

where $(R, \phi) = Z^{-1}S\hat{U}$. It then follows that

$$\hat{\mathcal{F}}(\hat{U}) = S^{-1} \mathcal{F}(U),$$

with $\mathcal{F}(\cdot)$ defined as in (5.23). We note that due to the structure of $S^{-1} D_0^\infty S$, the zero-order terms in $\hat{\mathcal{F}}(\hat{U})$ appear only in the first component. Using this structural property of (5.25), we will choose the ansatz for $\hat{R}$ so that the zero-order term in the $\hat{R}$ equation cancels with the leading order terms involving $\phi$ and $p$, that come from the functions $T$ and $\mathcal{Q}$.

To compute the terms in $\hat{\mathcal{F}}(\hat{U})$, notice that

$$L_1^0 S\hat{U} = L_1^0 S \begin{pmatrix} \hat{R} \\ \hat{\phi} \end{pmatrix}$$

and so

$$S^{-1} L_1^0 S \hat{U} = \begin{pmatrix} \frac{\beta^*}{r} \left( 2\phi_x r_x + \alpha r_{xx} \right) & 0 \\ \left( \beta^* \right) & 0 \end{pmatrix} \begin{pmatrix} \hat{R} \\ \hat{\phi} \end{pmatrix}$$

Note that the zero-order terms in the $\hat{R}$ and $\hat{\phi}$ equation are exponentially localized in space, and so they will not contribute to the leading order dynamics. Similarly, we have

$$S^{-1} T(p) = p_t \begin{pmatrix} r_x \\ \beta_x r_x - r\phi_x \end{pmatrix} + r\phi_x p_{xx} \begin{pmatrix} \alpha \\ 1 + \alpha \beta_s \end{pmatrix} - 2r\phi_x^2 p_x \begin{pmatrix} -1 \\ \alpha - \beta_s \end{pmatrix}$$

and

$$S^{-1} \mathcal{Q}(R, \phi, p)$$

$$= \begin{pmatrix} 2\phi_x \hat{\phi}_x R - r\phi_x^2 & -3r \hat{R}^2 + 10y_1 r^3 \hat{R}^2 - 3r\phi_x^2 p_x^2 - 4r\phi_x p_x \phi_x + 2r\phi_x^2 \hat{R} p_x \\ (\alpha - \beta_s) r(\phi_x + \phi_s p_x)^2 + \frac{4r_0^2(\beta - \gamma_1)}{1-2y_1 r_0^2} \hat{R}^2 \end{pmatrix}.$$
Next, we denote by $\Pi_R$ and $\Pi_\phi$ the projection on the first and second component of a vector in $\mathbb{R}^2$, respectively. First consider the $R$ equation, by taking the projection $\Pi_R$ of (5.25). We wish to choose and approximate solution so that all leading orders terms in (5.25) cancel. 

First consider the terms that are linear in $\hat{R}, \phi, p$:

$$-2(r\psi_+ + a_+ \psi_+ \phi_+ + 2r_0(1 - 2\gamma r_0^2) \hat{R} = 2r_0^2 p_+.$$  

Recall that we are primarily concerned with the behaviour at $\pm \infty$. Since $r \psi_+ \to r_0 k_0$ as $x \to \pm \infty$, we therefore define

$$\hat{R}_0^a := \frac{k_0}{r_0(1 - 2\gamma r_0^2)} \left( \phi_0^a + \psi_+ p_+ \right),$$  

for some $\hat{\phi}^a$ to be determined later. It is then clear that the function $\hat{U}_0^a = (\hat{R}_0^a, r\hat{\phi}^a)$ with $\hat{R}_0^a$ defined as in (5.26) approximately solves the first equation (5.25) in the following sense:

$$\Pi_R \hat{F}(\hat{U}_0^a) = O(e^{-\Theta(|\hat{\phi}^a| + |p_+|)}) + O\left(|\hat{\phi}_x^a| + |\hat{\phi}_x^a|^2 + |p_+| + |p_+|^2 \right).$$  

This is not quite good enough, due to the terms that are quadratic in $\hat{\phi}^a$—recall that we need only the good terms appearing in (5.15). To fix this, we revise the ansatz (5.26) by introducing

$$\hat{R}^a := \frac{k_0}{r_0(1 - 2\gamma r_0^2)} \left( \phi_0^a + \psi_+ p_+ \right) + \hat{R}_1^a := \hat{R}_0^a + \hat{R}_1^a,$$  

with $\hat{R}_1^a$ to be chosen so that it takes care of all quadratic nonlinear or next order decaying terms in the error that appear on the right of (5.27). Precisely, we take

$$2r_0^2 \left(1 - 2\gamma r_0^2\right) \hat{R}_1^a = -(\hat{R}_0^a)_x - 2(\alpha + \beta_+) \psi_+ (\hat{R}_0^a)_x + 2r_+ (\hat{\phi}_x^a + \psi_+ p_+) \hat{R}_0^a - r_0^2 \hat{\phi}_x^a - 3r_0^2 p_+^2 - 4r_0 p_+ \phi_+ - (3r_0 - 10\gamma r_0^3) (\hat{R}_0^a)_x^2.$$  

Here the $t$-derivative term appearing on the right can be replaced by its spatial derivatives thanks to estimate (5.12). Thus, if we define $\hat{U}^a = (\hat{R}^a, r\hat{\phi}^a)$, then $\hat{U}^a$ approximately solves the first equation in (5.25) with a good error as claimed:

$$\Pi_R \hat{F}(\hat{U}^a) = O(\Theta(\hat{\phi}^a, p_+)), \quad (5.30)$$  

where $\Theta$ is defined in (5.15).

It remains to show that

$$\Pi_\phi \hat{F}(\hat{U}^a) = O(\Theta(\hat{\phi}^a, p_+)). \quad (5.31)$$  

Observe that since there is no zero-order term $\hat{R}$ appearing in $\Pi_\phi \hat{F}(\hat{U}^a)$, the contribution of $\hat{R}_1^a$ from (5.28) can be absorbed into the good error term. For this reason, we shall check the contribution of the first term in $\hat{R}_1^a$, or precisely $\hat{R}_0^a$ defined as in (5.26). By plugging the ansatz for $\hat{R}^a$ into $\hat{F}(\hat{U}^a)$ and using (5.12), it follows by direct calculations that

$$\Pi_\phi \hat{F}(\hat{U}^a) = r \left[ L_B \hat{\phi}^a + \psi_+ L_B p - q \left( \hat{\phi}_x^a + \psi_+ p_+ \right)^2 \right] + O(\Theta(\hat{\phi}_x^a, p_+)) \quad (5.32)$$  

where we recall that $L_B$ is defined by

$$L_B = \partial_t + 2(\alpha - \beta_+) \psi_+ \partial_x - \beta_0^2.$$  

By remark 5.2, we have $L_B (\hat{\phi}^a \pm k_0 p_+)$ for each $+/-$ case. This, together with the fact that $\psi_+ \to \pm k_0$ as $|x| \to \infty$, shows that the leading term on the right-hand side of (5.32) vanishes. We thus obtain (5.31) as claimed.

Going back to the original variable $U^n = S \hat{U}^a$, we find that

$$\hat{F}(U^n) = (\partial_t - L^n) \hat{U}^a - T(p) - Q(\hat{R}^a, \hat{\phi}^a, p) = O(\Theta(\hat{\phi}_x^a, p_+)).$$  

This together with (5.24) establishes proposition 5.3 in the case that $\delta^{\pm}$ are constants. The case when $\delta^{\pm}(\tau)$ are not constants follows similarly.
5.5. Proof of proposition 5.4

We now choose $B(x, t)$ to be $e(x, t+1)$, the difference of error functions defined in (1.5). We then have

$$|B(1-B)|(x, t) \leq C \left( e^{-\frac{\left(\text{sign}B\right)^2}{\text{C}(1)}} + e^{-\frac{\left(\text{sign}B\right)^2}{\text{C}(1)}} \right) \leq C(1+t)^{1/2}\theta(x, t),$$

and so

$$\left[ \frac{B}{1+\delta^+B} - \frac{B}{1+\delta^+B} \right](x, t) \leq C|\delta^+| \left( e^{-\frac{\left(\text{sign}B\right)^2}{\text{C}(1)}} + e^{-\frac{\left(\text{sign}B\right)^2}{\text{C}(1)}} \right) \leq C(1+t)^{1/2}\theta(x, t)|\delta^+|.$$

We also have

$$\left| \frac{B}{1+\delta^+B} \right|(x, t) \leq C(1+t)^{-1/2} \left( e^{-\frac{\left(\text{sign}B\right)^2}{\text{C}(1)}} + e^{-\frac{\left(\text{sign}B\right)^2}{\text{C}(1)}} \right) \leq C\theta(x, t).$$

Thus,

$$|\delta^+ \hat{\phi}^\alpha| \leq C(1+t)^{1/2-k/2}\theta(x, t)(|\delta^+| + |\delta^-|), \quad |\delta^+ p| \leq C(1+t)^{1/2-k/2}\theta(x, t)(|\delta^+| + |\delta^-|),$$

for $k \geq 1$. Here these bounds are valid uniformly in $x \in \mathbb{R}$ and $t \geq 0$, provided that $\delta^\pm$ is sufficiently small, with a bound on $\delta^\pm$ that depends only on the maximum of $|B|$. By using these bounds, it follows immediately that

$$\Upsilon(\hat{\phi}^\alpha, p_x) \leq C(|\delta^+| + |\delta^-|)\varepsilon^{-q_k}\theta(x, t) + C(1+t)^{-1}\theta(x, t)(|\delta^+|^2 + |\delta^-|^2).$$

**Proof of proposition 5.4.** The estimate for $R^{\text{app}}_\alpha(x, t, \delta^\pm)$ follows directly from (5.16) and (5.34). It remains to give details of $\Sigma^\pm(x, t, \delta^\pm)$ introduced in proposition 5.3. By (5.9), we obtain

$$\frac{\partial p}{\partial \delta^\pm}(x, t, \delta^\pm(t)) = \frac{d}{d} \frac{B}{1+\delta^\pm B}, \quad \frac{\partial p}{\partial \delta^\pm}(x, t, \delta^\pm(t)) = -\frac{d}{2qk_0} \frac{B}{1+\delta^\pm B}.$$

By definition (5.28), we obtain a similar result for $\frac{\partial}{\partial \delta^\pm} \hat{R}^\alpha(x, t, \delta^\pm(t))$. This together with (5.33) yields

$$\frac{\partial p}{\partial \delta^\pm}(x, t, \delta^\pm(t)) = \frac{d}{2qk_0} \frac{B}{1+\delta^\pm} + O(|\delta^\pm|(1+t)^{1/2}\theta(x, t))$$

and

$$\frac{\partial \phi^\alpha}{\partial \delta^\pm}(x, t, \delta^\pm(t)) = \frac{d}{2q} \frac{B}{1+\delta^\pm} + O(|\delta^\pm|(1+t)^{1/2}\theta(x, t)).$$

Putting these estimates into $\Sigma^\pm(x, t, \delta^\pm)$ (see (5.17)), we obtain

$$\Sigma^\pm(x, t) = \frac{d}{2q(1+\delta^\pm)} \left[ \left( \frac{\phi^\alpha \pm k_0}{r_0(1-2\gamma r_0^2)} \right) B_x \frac{B}{k_0} \left( \frac{r_x}{r_\phi^\alpha} \right) \right] + O(|\delta^\pm|(1+t)^{1/2}\theta(x, t)).$$

Here the error term $O(|\delta^\pm|(1+t)^{1/2}\theta(x, t)$ times $\delta^\pm$ can be absorbed into the last term in the remainder $R^{\text{app}}_\alpha(x, t, \delta^\pm)$. Rearranging terms, we obtain the proposition as claimed. \[\square\]
Proof of lemma 5.5. Note that
\[ E^\pm = BV_2 \pm \frac{R}{k_0} V_1 + \mathcal{O}(\theta), \]
where \( V_{1,2} \) are the eigenfunctions of \( \mathcal{L} \) defined in (2.7). A direct calculation shows that
\[ (\partial_t - \mathcal{L})E^\pm = \mathcal{O}(e^{-\eta|x|}B_x) + \mathcal{O}((\partial_t - \mathcal{L})\theta). \]
The final term on the right-hand side is of the claimed order because the Gaussian terms each satisfy the asymptotic equation at one end \((\pm \infty)\), where the maximum \( x = \pm c_g t \) is located, and hence decay faster at the other end, away from this maximum.

6. Stability analysis

In this section, we show how the ansatz (5.2) and the approximate solution \((\phi^a, R^a, p)\), constructed in section 5.2, can be used to prove theorem 1.1.

6.1. Nonlinear perturbed equations

Recall, that the ansatz (5.2) gives
\[ B(x,t) = A(x + p(x,t), t) = [r(x) + R(x,t)]e^{i(\phi(x) + \phi(x,t))}e^{-\eta|t|}, \]
where \( A(x,t) \) solves the qCGL equation (5.1). The functions \((R, \phi)\) represent the perturbation, and we define \((\tilde{R}, \tilde{\phi})\) via
\[ (R, \phi)(x,t) = (R^a, \phi^a)(x,t, \delta^\pm(t)) + (\tilde{R}, \tilde{\phi})(x,t), \quad (6.1) \]
where \( U^a = (R^a, r\phi^a) \) is the approximate solution constructed in proposition 5.4. Note that the functions \( \delta^\pm(t) \) are still arbitrary smooth functions. They will be chosen, in equation (6.16), so that they cancel non-decaying terms in the nonlinear iteration scheme. The following lemma is then a direct combination of lemma 5.1 and proposition 5.4.

Lemma 6.1. Let \( \delta^\pm = \delta^\pm(t) \) be arbitrary smooth functions. The perturbation function \( \tilde{U} = (\tilde{R}, \tilde{\phi}) \) defined through (6.1) satisfies
\[
(\partial_t - \mathcal{L})\tilde{U} = -\frac{d}{2q} \sum_{\pm} \frac{\delta^\pm(t)}{1 + \delta^\pm(t)} E^\pm(x,t) + N_1(\tilde{R}, \tilde{\phi}, \delta^\pm)(x,t),
\]
with
\[
\begin{align*}
\mathcal{L}(\tilde{R}, \tilde{\phi}) &= \mathcal{O}(|\tilde{R}|^2 + |\tilde{R}_x|^2 + |\tilde{\phi}|^2 + |\tilde{R} \tilde{\phi}|^2 + |\tilde{R} \tilde{R}_x|^2 + |\tilde{\phi}_x|^2), \\
Q(\tilde{R}, \tilde{\phi}) &= \mathcal{O}(|\tilde{R}|^2 + |\tilde{R}_x|^2 + |\tilde{\phi}|^2 + |\tilde{R} \tilde{\phi}|^2 + |\tilde{R} \tilde{R}_x|^2 + |\tilde{\phi}_x|^2).
\end{align*}
\]

Proof. By lemma 5.1, the perturbation \( U = (R, \phi) \) solves
\[ (\partial_t - \mathcal{L})U - \mathcal{T}(p) - Q(R, \phi, p) = \mathcal{N}(R, \phi, p), \]
We recall the Green function decomposition of proposition 4.2:

\[ (\partial_t - \mathcal{L})U^a - T(p) - Q(R^a, \phi^a, p) = \frac{d}{dt} \sum_{\pm} \hat{\delta}^\pm(t) \mathcal{E}^\pm(x, t) + \mathcal{R}^\text{app}_{1\pm}(x, t, \delta^\pm). \]

Here the same notations as in lemma 5.1 and proposition 5.4 are used. It then follows that the difference \( \hat{U} = U - U^a \) solves (6.2) with the remainder \( \mathcal{N}_1(\hat{R}, \hat{\phi}, \delta^\pm)(x, t) \) defined by

\[ \mathcal{N}_1(\hat{R}, \hat{\phi}, \delta^\pm)(x, t) := \mathcal{N}(R, \phi, p) + Q(R, \phi, p) - Q(R^a, \phi^a, p) - \mathcal{R}^\text{app}_{1\pm}(x, t, \delta^\pm). \]

The claimed estimate for \( \mathcal{N}_1(\hat{R}, \hat{\phi}, \delta^\pm) \) now follows directly from (5.8) and (5.20).

### 6.2. Green function decomposition

We recall the Green function decomposition of proposition 4.2:

\[ \mathcal{G}(x, y, t) = e(x, t) V_1(x) \langle \psi_1(y), \cdot \rangle + e(x, t) V_2(x) \langle \psi_2(y), \cdot \rangle + \hat{\mathcal{G}}(x, y, t), \]

where \( e(x, t) \), defined in (1.5), is the sum of error functions and \( \hat{\mathcal{G}}(x, y, t) \) satisfies the Gaussian-like estimate

\[ \hat{\mathcal{G}}(x, y, t) = \sum_{\pm} \mathcal{O}(t^{-\frac{1}{2}} e^{-\frac{\|x-y\|^2}{4t^2}}) + \mathcal{O}(e^{-\gamma(x-y)^2}). \]

Now observe that from the definitions of \( V_j(x) \) in (2.7) and \( E^\pm(x, t) \) in (5.19), we have

\[ E^\pm(x, t) = -B(x, t) V_2 \pm \frac{B(x, t)}{k_0} V_1 + \mathcal{O}(\theta(x, t)), \]

(6.5)

where \( B(x, t) = e(x, t + 1) \). In addition, since \( \psi_1(y) \) and \( \psi_2(y) \) are exponentially localized and \( \hat{\mathcal{G}}(x, y, t) \) behave as Gaussians, we can write

\[ e(x, t + 1) V_1(x) \langle \psi_1(y), \cdot \rangle + e(x, t + 1) V_2(x) \langle \psi_2(y), \cdot \rangle = E^+(x, t) \langle \Psi^+(y), \cdot \rangle + E^-(x, t) \langle \Psi^-(y), \cdot \rangle, \]

(6.6)

where

\[ \Psi^+ = \frac{1}{2} \psi_1 - \frac{k_0}{2} \psi_2, \quad \Psi^- = \frac{1}{2} \psi_1 + \frac{k_0}{2} \psi_2, \]

with an error of \( \mathcal{O}(e^{-\gamma|y|/|B_e|} + |e_\tau|) = \mathcal{O}(e^{-\gamma|\psi|}) \theta(x, t) \) that can easily be absorbed into the bound for \( \hat{\mathcal{G}}(x, y, t) \). In addition, we observe that the difference \( e(x, t) - e(x, t + 1) \) satisfies

\[ |e(x, t) - e(x, t + 1)| = \left| \int_t^{t+1} e_\tau(x, \tau) \, d\tau \right| \leq C t^{-\frac{1}{2}} e^{-\frac{\|x-y\|^2}{4t^2}}, \quad t \geq 1 \]

and

\[ |e(x, t) - e(x, t + 1)| \leq C \int_t^{t+1} e^{-z^2} \, dz \leq C e^{-t^2/2} e^{-ct^2} \quad t \leq 1. \]

These, together with the localization property of \( \psi_j(y) \), show that the term

\[ (e(x, t) - e(x, t + 1)) V_1(x) \langle \Psi^+(y), \cdot \rangle + (e(x, t) - e(x, t + 1)) V_2(x) \langle \Psi^-(y), \cdot \rangle \]

can be absorbed into the bound for \( \hat{\mathcal{G}}(x, y, t) \). Therefore, we decompose the Green function as follows:

\[ \mathcal{G}(x, y, t) = \mathcal{E}^+(x, t) \langle \Psi^+(y), \cdot \rangle + \mathcal{E}^-(x, t) \langle \Psi^-(y), \cdot \rangle + \hat{\mathcal{G}}(x, y, t), \]

(6.7)

where \( \mathcal{G}(x, y, t) \) satisfies the same bound as \( \hat{\mathcal{G}}(x, y, t) \). The motivation for this will become clear below.

We obtain the following simple lemma.
Lemma 6.2. For all $s, t \geq 0$

$$
\int_{\mathbb{R}} G(x, y, t - s) \mathcal{E}^\pm(y, s) \, dy = \mathcal{E}^\pm(x, t) + \int_{s}^{t} \int_{\mathbb{R}} G(x, y, t - \tau) \mathcal{R}^\pm(y, \tau) \, dy \, d\tau,
$$

where

$$
|\mathcal{R}^\pm(y, s)| \leq C((1 + s)^{-1} + e^{-|y|}) \theta(y, s).
$$

(6.8)

**Proof.** We recall from lemma 5.5 that

$$(\partial_t - L)\mathcal{E}^\pm(x, t) = \mathcal{O}((1 + t)^{-1} + e^{-|y|}) \theta(x, t).$$

By applying the Duhamel principle to the solution of this equation with ‘initial’ data at $t = s$, we easily obtain the lemma.

\[ \square \]

6.3. Initial data for the asymptotic ansatz

We denote $U^a_0(x) := U^a(x, 0, \delta^\pm(0))$, $\tilde{U}_0(x) = \tilde{U}(x, 0)$, $p_0(x) := p(x, 0, \delta^\pm(0))$, and

$$
U_{\text{source}}(x) := \begin{pmatrix} r(x) \\ r(x)\phi(x) \end{pmatrix}, \quad U_{\text{source}, r}(x) := \begin{pmatrix} r(x + p_0(x)) \\ r(x)\phi(x + p_0(x)) \end{pmatrix}.
$$

By (1.4), (6.1) and the assumption on the initial data $A_m(x)$, we have

$$
\left\| e^{\mp \gamma \mathcal{L}_0} \left( U^a_0(\cdot) + \tilde{U}_0(\cdot) + U_{\text{source}}(\cdot) - U_{\text{source}, r}(\cdot) \right) \right\|_{C^1(\mathbb{R})} \leq \epsilon, \tag{6.9}
$$

for $\epsilon = \| A_m(\cdot) - A_{\text{source}}(\cdot) \|_{\text{lin}}$. We prove the following.

**Lemma 6.3.** There are small constants $\delta_0^\pm$ so that if we take $\delta^\pm(0) = \delta_0^\pm$ in the definition of $U^a_0(x)$ and $p_0(x)$ then

$$
\int_{\mathbb{R}} (\Psi^\pm(y), \tilde{U}_0(y)) \, dy = 0 \tag{6.10}
$$

for each $+/-$ case, with $\Psi^\pm(\cdot)$ the same as in (6.7). In addition, $\sum_\pm |\delta_0^\pm| \leq C\epsilon$, and

$$
\left\| e^{\mp \gamma \mathcal{L}_0} \tilde{U}_0(\cdot) \right\|_{C^1(\mathbb{R})} \leq C\epsilon, \tag{6.11}
$$

for some positive constant $C$.

**Proof.** We denote

$$
F_0(x) := U^a_0(x) + \tilde{U}_0(x) + U_{\text{source}}(x) - U_{\text{source}, r}(x),
$$

so that $\|e^{\mp \gamma \mathcal{L}_0} F_0(\cdot)\|_{C^1(\mathbb{R})} \leq \epsilon$ by (6.9). Note that $\Psi^\pm(\cdot)$ are linear combinations of $\Psi_1(\cdot)$, $\Psi_2(\cdot)$, which are in the nullspace of the adjoint of $\mathcal{L}$. Therefore, (6.10) will follow if we can find $\delta_0^\pm$ so that

$$
G_j(\delta_0^\pm) := \int_{\mathbb{R}} (\psi_j(y), U^a_0(y) + U_{\text{source}}(y) - U_{\text{source}, r}(y)) \, dy = \int_{\mathbb{R}} (\psi_j(y), F_0(y)) \, dy, \tag{6.12}
$$

for $j = 1, 2$. Notice that $G_j(0) = 0$ since $p_0$ and $U^a_0$ vanish when $\delta_0^\pm = 0$. Thus, we can define a new function

$$
\tilde{G}(F, \delta_0^\pm) = (G_1(\delta_0^\pm), G_2(\delta_0^\pm)) - \left( \int_{\mathbb{R}} (\psi_1(y), F_0(y)) \, dy, \int_{\mathbb{R}} (\psi_2(y), F_0(y)) \, dy \right).
$$
Since \( \tilde{G}(0,0) = 0 \), it suffices to show that the Jacobian determinant \( J_G \) of \((G_1(\delta_0^+), G_2(\delta_0^+))\) with respect to \( \delta_0^+ \) is non-zero at \((\delta_0^+, \delta_0^-) = 0\). For this, we compute
\[
\frac{\partial G_j(\delta_0^+)}{\partial \delta_0^+}|_{(\delta_0^+, \delta_0^-)=0} = \left( \psi_j(\cdot), \frac{\partial U_j^0(\cdot)}{\partial \delta_0^+} - \frac{\partial p_0(\cdot)}{\partial \delta_0^-} \right)_{L^2} = \left( \psi_j(\cdot), \frac{d}{2q} \left( \beta^* B_j - r B \right) \right)_{L^2} = \frac{d}{2q} \left( \beta^* B_j - r B \right)_{L^2}
\]
where we have used the definition of \( \tilde{p} \) and \( U_a^0 \); see (5.9). For our convenience, we denote
\[
\tilde{V}_1 = \left( \frac{r_y}{r \phi_j} \right) B(y,0), \quad \tilde{V}_2 = \left( \frac{B^*_y}{\beta^* B_j - r B} \right) (\cdot,0).
\]
It then follows that
\[
J_G = \frac{d^2}{2q^2 k_0} \left[ (\psi_1, \tilde{V}_1)_L^2 (\psi_2, \tilde{V}_2)_L^2 - (\psi_1, \tilde{V}_2)_L^2 (\psi_2, \tilde{V}_1)_L^2 \right] = \frac{d^2}{2q^2 k_0} \det M^\phi(0),
\]
where the matrix \( M^\phi(0) \) is defined as in (3.16) and is invertible. This proves that \( J_G \neq 0 \). The existence of small \( \delta_0^+ = \delta_0^+(F_0) \) satisfying (6.12) then follows from the standard implicit function theorem.

Finally, estimate (6.11) follows directly from the fact that the difference of the two error functions \( \tilde{B}(x,0) = \phi(x,1) \) decays to zero as fast as \( e^{-\gamma} \) when \( x \to \pm \infty \). \( \square \)

### 6.4. Integral representations

Using lemmas 4.1 and 6.1, we obtain the integral formulation of the solution \( \tilde{U} = (\tilde{R}, r \tilde{\phi}) \) of (6.2):
\[
\tilde{U}(x,t) = \int_{\mathbb{R}} G(x,y,t) \tilde{U}_0(y) dy + \int_t^0 \int_{\mathbb{R}} G(x, y, t-s) \mathcal{N}_1(\tilde{R}, \tilde{\phi}_y, \delta^\pm)(y,s) dy ds - \frac{d}{2q} \int_0^t \int_{\mathbb{R}} G(x, y, t-s) \left[ \frac{\delta^+(s)}{1 + \delta^+(s)} \mathcal{E}^+(y,s) + \frac{\delta^-(s)}{1 + \delta^-(s)} \mathcal{E}^-(y,s) \right] dy ds. \tag{6.13}
\]
For simplicity, we denote
\[
\Theta^{\pm}(t) := \frac{d}{2q} \left[ \log(1 + \delta^{\pm}(t)) - \log(1 + \delta_0^{\pm}) \right] \tag{6.14}
\]
with \( \delta^{\pm}_0 \) defined as in lemma 6.3. Using lemma 6.2, the last integral term in (6.13) is equal to
\[
- \sum_{\pm} \int_0^t \frac{d}{ds} \Theta^{\pm}(s) \left[ \mathcal{E}^{\pm}(x,t,s) + \int_s^t \int_{\mathbb{R}} G(x, y, t-r) \mathcal{R}^{\pm}(y,r) dy dr ds \right] ds
\]
\[
= - \sum_{\pm} \Theta^{\pm}(t) \mathcal{E}^{\pm}(x,t) - \sum_{\pm} \int_0^t \int_{\mathbb{R}} G(x, y, t-s) \Theta^{\pm}(s) \mathcal{R}^{\pm}(y,s) dy ds.
\]

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Here in the last equality integration by parts in $s$ has been used. That is, we can write the integral formulation (6.13) as

$$
\tilde{U}(x, t) = \int_\mathbb{R} G(x, y, t) \tilde{U}_0(y) dy + \int_0^t \int_\mathbb{R} G(x, y, t - s) N_2(\tilde{R}, \tilde{\phi}_y, \delta^\pm)(y, s) dy ds - \sum_\pm \Theta^\pm(t) \mathcal{E}^\pm(x, t),
$$

(6.15)

with

$$
N_2(\tilde{R}, \tilde{\phi}_y, \delta^\pm)(y, s) := N_1(\tilde{R}, \tilde{\phi}_y, \delta^\pm)(y, s) - \sum_\pm \Theta^\pm(s) \mathcal{R}^\pm(y, s),
$$

where $N_1(\tilde{R}, \tilde{\phi}_y, \delta^\pm)$ is defined in (6.3) and $\mathcal{R}^\pm(y, s)$ is defined in (6.8).

We recall the decomposition (6.7) of Green’s function:

$$
G(x, y, t) = E^+(x, t) \langle \Psi^+(y), \cdot \rangle + E^-(x, t) \langle \Psi^-(y), \cdot \rangle + \tilde{G}(x, y, t).
$$

Using this decomposition, we write the integral formulation (6.15) as follows. To avoid repeatedly writing the lengthy integrals, we denote

$$
I^+(t) := \int_\mathbb{R} (\langle \Psi^+(y), \tilde{U}_0(y) \rangle) dy + \int_0^t \int_\mathbb{R} (\langle \Psi^+(y), N_2(\tilde{R}, \tilde{\phi}_y, \delta^\pm)(y, s) \rangle) dy ds
$$

$$
I^-(t) := \int_\mathbb{R} (\langle \Psi^-(y), \tilde{U}_0(y) \rangle) dy + \int_0^t \int_\mathbb{R} (\langle \Psi^-(y), N_2(\tilde{R}, \tilde{\phi}_y, \delta^\pm)(y, s) \rangle) dy ds
$$

$$
I_d(x, t) := \int_\mathbb{R} \tilde{G}(x, y, t) \tilde{U}_0(y) dy + \int_0^t \int_\mathbb{R} \tilde{G}(x, y, t - s) N_2(\tilde{R}, \tilde{\phi}_y, \delta^\pm)(y, s) dy ds
$$

$$
+ \int_0^t \int_\mathbb{R} \left( E^+(x, t - s) - E^+(x, t) \right) \langle \Psi^+(y), N_2(\tilde{R}, \tilde{\phi}_y, \delta^\pm)(y, s) \rangle dy ds
$$

$$
+ \int_0^t \int_\mathbb{R} \left( E^-(x, t - s) - E^-(x, t) \right) \langle \Psi^-(y), N_2(\tilde{R}, \tilde{\phi}_y, \delta^\pm)(y, s) \rangle dy ds,
$$

referring to the contributions accounting for the translation and phase shifts, and the decaying part, respectively. Thus, the integral formulation (6.15) for solutions $(\tilde{R}, \tilde{\phi})$ to (6.2) now becomes

$$
\tilde{U}(x, t) = - \sum_\pm \Theta^\pm(t) \mathcal{E}^\pm(x, t) + \sum_\pm I^\pm(t) \mathcal{E}^\pm(x, t) + I_d(x, t).
$$

Neither of the two terms $I^\pm(t) \mathcal{E}^\pm(x, t)$ decay in time. To capture these non-decaying terms, we are led to choose $\Theta^\pm$ such that

$$
\Theta^\pm(t) = I^\pm(t),
$$

or equivalently,

$$
\log(1 + \delta^\pm(t)) = \log(1 + \delta^\pm_0) + \frac{2q_d}{d} I^\pm(t).
$$

(6.16)

Note that such a choice is possible for the following reason. Lemma 6.3 implies that the above equation is satisfied at $t = 0$ if $\delta^\pm_0$ are appropriately chosen. We can then just define $\delta^\pm(t)$ to be a solution of the corresponding integral equation. Thus, the representation for solutions $(\tilde{R}, \tilde{\phi})$ simply reads

$$
\tilde{U}(x, t) = I_d(x, t).
$$

(6.17)
6.5. Spatio-temporal template functions

In this section, we introduce template functions that are useful for the construction and estimation of the solutions. We let

\[ h(t) := h_1(t) + h_2(t) \]

with

\[ h_1(t) := \sup_{0 \leq s \leq t} \left( |\delta^+(s)| + |\delta^-(s)| \right) e^{\kappa s}, \]

\[ h_2(t) := \sup_{0 \leq s \leq t, y \in \mathbb{R}} (1 + s)^{-\kappa} \left( \frac{|\tilde{\phi}|}{\theta} + |\tilde{\phi}_y| + |\tilde{\phi}_{yy}| + |R| + |\tilde{R}| + |\tilde{R}_{yy}| \right) (y, s) \]

for some \( \kappa \in (0, 1/2) \) and some fixed, small \( \eta > 0 \). Here, \( \theta(x, t) \) denotes the Gaussian-like behavior defined in (1.6). We note that the constant \( M_0 \) in (1.6) is a fixed, large, positive number. At various points in the below estimates, there will be a similar quantity, which we denote by \( M_0 \), that will need to be taken to be sufficiently large. The number \( M_0 \) is then the maximum value of \( M_0 \) at the end of the proof. From standard short time theory, we see that \( h(t) \) is well defined and continuous for \( 0 < t \ll 1 \). In addition, standard parabolic theory implies that \( h(t) \) retains these properties as long as \( h(t) \) stays bounded. The key issue is therefore to show that \( h(t) \) stays bounded for all times \( t > 0 \), and this is what the following proposition asserts.

**Proposition 6.4.** There exists an \( \epsilon_0 > 0 \) sufficiently small such that the following holds. Given any initial data \( A_{in} \) with \( \epsilon := \| A_{in}(\cdot) - A_{source}(\cdot) \|_{in} \leq \epsilon_0 \) and any \( \kappa \in (0, 1/2) \), there exist positive constants \( \kappa, C_0, \) and \( M_0 \) such that

\[ h_1(t) \leq C_0(\epsilon + h^2(t)), \quad h_2(t) \leq C_0(\epsilon + h_1(t) + h(t)^2), \] (6.18)

for all \( t \geq 0 \).

Using this proposition, we can add the inequalities in (6.18) and eliminate \( h_1 \) on the right-hand side to obtain

\[ h(t) \leq C_0(C_0 + 2)(\epsilon + h(t)^2). \]

Using this inequality and the continuity of \( h(t) \), we find that \( h(t) \leq 2C_0(C_0 + 2)\epsilon \) provided \( 0 \leq \epsilon \leq \epsilon_0 \) is sufficiently small. Thus, the main theorem will be proved once we establish proposition 6.4. The following sections will be devoted to proving this proposition.

6.6. Bounds on the nonlinear terms

We recall that the nonlinear remainder \( N_2(\tilde{R}, \tilde{\phi}, \delta^\pm)(x, t) \) is defined by

\[ N_2(\tilde{R}, \tilde{\phi}, \delta^\pm)(x, t) = N_1(\tilde{R}, \tilde{\phi}, \delta^\pm)(x, t) - \sum_{\pm} \Theta^\pm(t) R^\pm(x, t), \]

with \( \Theta^\pm(s) \) defined in (6.14), \( R^\pm(x, t) \) defined in (6.8), and \( N_1(\tilde{R}, \tilde{\phi}, \delta^\pm) \) defined in (6.3). We first note that

\[ |\delta^\pm(t)| \leq |\delta^\pm(0)| + \int_0^t |\delta^\pm(s)| ds \leq |\delta^\pm_0| + \int_0^t e^{-\eta s} h_1(s) ds \leq C(\epsilon + h_1(t)), \]

where we recall that lemma 6.3 implies that \( \delta^\pm_0 \sim \epsilon \). By (5.20) and the definition of \( h_1(t) \)

\[ |R_0^{app}(x, t, \delta^\pm)| \leq C(|\delta^+| + |\delta^-|) e^{-\eta t} |\theta(x, t) + C(1 + t)^{-1} \theta(x, t)(\epsilon + h_1(t)^2). \]
Here we leave the linear term in $\delta^\pm(t)$ in the above estimate for a different treatment, below. Similarly, by definitions (6.14) and (6.8), we obtain
\[
| \sum_{\pm} \Theta_{\pm}(t) R_{\pm}(\cdot, t) | \leq C \left( (1 + t)^{-1} + e^{-\eta|x|} \right) (|\delta^+| + |\delta^-|) \theta(x, t).
\]

Next, from (6.4) and the definition of $h_2(t)$, we have
\[
L(\tilde{R}, \tilde{\phi}) = \mathcal{O}(|\tilde{R}| + |\tilde{\phi}_x|) \leq C (1 + t)^{-\frac{1}{2}} \theta(x, t) h_2(t),
\]
\[
Q(\tilde{R}, \tilde{\phi}) = \mathcal{O}(\tilde{R}^2 + |\tilde{R}\tilde{\phi}_x| + \tilde{\phi}_x^2 + |\tilde{R}\tilde{\phi}_t| + |\tilde{\phi}_t|) \leq C (1 + t)^{-1 + 2\epsilon} \theta(x, t) h_2(t).
\]
Recalling estimate (6.3) for $\mathcal{N}_1$, we find that the nonlinear term $\mathcal{N}_2(\tilde{R}, \tilde{\phi}, \delta^\pm)(x, t)$ satisfies
\[
|\mathcal{N}_2(\tilde{R}, \tilde{\phi}, \delta^\pm)(x, t)| \leq C \left[ e^{-\eta|x|} + (1 + t)^{-1} \right] \theta(x, t) \varphi(x, t) \left( |\delta^+| + |\delta^-| \right) + C (1 + t)^{-1 + 2\epsilon} \theta(x, t)^2 h_2(t).
\]
where we have used $\theta(x, t) \leq 2(1 + t)^{-\kappa}$ and $(1 + t) \varphi \leq C (1 + t)^{-1}$.

**6.7. Estimates for $h_1(t)$**

To establish the estimate for $h_1(t)$, we differentiate expression (6.16) to get
\[
\frac{d}{dt} \sum_{\pm} |\delta^\pm(t)|^2 \leq C e^{-\eta t} \sum_{\pm} |\delta^\pm(t)|^2 + C e^{-\eta t} \left[ \epsilon + h(t)^2 \right]^2,
\]
for each $+/-$ case. We first recall that $|\Psi^\pm(y)| \leq 2e^{-\eta|y|}$ and that
\[
e^{-\eta t} e^{-\frac{\eta}{M} \phi^\pm} \leq C e^{-\frac{\eta}{M} \phi^\pm} \leq C e^{-\frac{\eta}{M} \phi^\pm},
\]
which holds for each $M \geq 8\varepsilon^2 / \eta_0$ and $\eta$ sufficiently small. This, together with the bound (6.19) on the nonlinear term $\mathcal{N}_2(\tilde{R}, \tilde{\phi}, \delta^\pm)$, implies
\[
|\langle \Psi^\pm(y), \mathcal{N}_2(\tilde{R}, \tilde{\phi}, \delta^\pm)(y, t) \rangle | \leq C e^{-\frac{\eta}{M} \phi^\pm} e^{-\eta t} \left[ \epsilon + |\delta^+| + |\delta^-| + h(t)^2 \right].
\]
and so
\[
|\int \langle \Psi^\pm(y), \mathcal{N}_2(\tilde{R}, \tilde{\phi}, \delta^\pm)(y, t) \rangle dy | \leq C e^{-\eta t} \left[ \epsilon + |\delta^+| + |\delta^-| + h(t)^2 \right].
\]
Now, multiplying equation (6.20) by $\delta^\pm$ and using the above estimate on the integral, we obtain
\[
\frac{d}{dt} \sum_{\pm} |\delta^\pm(t)|^2 \leq C e^{-\eta t} \sum_{\pm} |\delta^\pm(t)|^2 + C e^{-\eta t} \left[ \epsilon + h(t)^2 \right]^2,
\]
where we have used Young’s inequality and the fact that, as long as $\delta^\pm$ is bounded, higher powers of $\delta^\pm$ can be bounded by $C|\delta^\pm|^2$. Applying the standard Gronwall’s inequality, we obtain
\[
\sum_{\pm} |\delta^\pm(t)|^2 \leq \sum_{\pm} |\delta^\pm(0)|^2 e^{\int_0^t C e^{-\eta s} ds} + C \int_0^t e^{\int_0^s C e^{-\eta r} dr} e^{-\eta t} \left[ \epsilon + h(s)^2 \right]^2 ds.
\]
Since $\int_0^t e^{-\eta s} ds$ is bounded, $|\delta^\pm(0)| \leq C \epsilon$, and $h(t)$ is an increasing function, the above estimate yields
\[
\sum_{\pm} |\delta^\pm(t)|^2 \leq C \left[ \epsilon + h(t)^2 \right]^2.
\]
Using this into (6.22) and in (6.20), we immediately obtain
\[
|\delta^\pm(t)| \leq C e^{-\eta t} \left[ \epsilon + h(t)^2 \right].
\]
which yields the first inequality in (6.18).

Finally, note that if we combine the bound (6.23) with the nonlinear estimate (6.19), we obtain

\begin{equation}
N_2(\bar{R}, \tilde{\phi}, \delta^\pm)(x, t) \leq C \left[ e^{-\eta |y|} + (1 + t)^{-1/2} \right] \theta(x, t)(\epsilon + h^2(t)).
\end{equation}

(6.24)

6.8. Pointwise estimates for \( \bar{R} \) and \( \tilde{\phi} \)

Let \( \tilde{U} = (\bar{R}, r\tilde{\phi}) \) satisfy the integral formulation (6.17). We shall establish the following pointwise bounds

\begin{equation}
|\bar{R}(x, t)| \leq C \left( \epsilon + h(t)^2 \right)(1 + t)^{-1/2} \theta(x, t),
\end{equation}

(6.25)

and the derivative bounds

\begin{equation}
|\partial_k^\ell \tilde{U}(x, t)| \leq C \left( \epsilon + h(t)^2 \right)(1 + t)^{-1/2} \theta(x, t)
\end{equation}

(6.26)

for \( k + \ell \leq 3 \) with \( \ell = 0, 1 \) and \( k = 1, 2, 3 \). We recall the integral formulation (6.17):

\[ \tilde{U}(x, t) = \int_{\mathbb{R}} \tilde{G}(x, y, t) \tilde{U}_0(y) dy + \int_0^t \int_{\mathbb{R}} \tilde{G}(x, y, t-s) N_2(\bar{R}, \tilde{\phi}, \delta^\pm)(y, s) dy ds \]

\[ + \sum_{\pm} \int_0^t \int_{\mathbb{R}} (\mathcal{E}^\pm(x, t-s) - \mathcal{E}^\pm(x, t)) (\Psi^\pm(y), N_2(\bar{R}, \tilde{\phi}, \delta^\pm)(y, s)) dy ds. \]

(6.27)

We give estimates for each term in this expression. First, we consider the integral term in (6.27) that involves the initial data. We recall that

\begin{equation}
|\tilde{G}(x, y, t)| \leq C t^{-1/2} \left( e^{-\frac{(x-y)^2}{4\pi}} + e^{-\frac{(x-y)^2}{4\pi}} \right).
\end{equation}

(6.28)

Using this bound, together with equation (6.11), we see that

\[ \int_{\mathbb{R}} |\tilde{G}(x, y, t) \tilde{U}_0(y)| dy \leq C \epsilon \int_{\mathbb{R}} t^{-1/2} \left( e^{-\frac{(x-y)^2}{4\pi}} + e^{-\frac{(x-y)^2}{4\pi}} \right) e^{-\frac{y^2}{2\pi}} dy. \]

(6.29)

Using the fact that, for \( t \geq 1 \),

\[ e^{-\frac{(x-y)^2}{4\pi}} e^{-\frac{y^2}{2\pi}} \leq C_1 e^{-\frac{|x-y|^2}{4\pi}} e^{-\frac{y^2}{2\pi}}, \]

and, for \( t \leq 1 \),

\[ e^{-\frac{(x-y)^2}{4\pi}} e^{-\frac{y^2}{2\pi}} \leq 2 e^{-\frac{|x-y|^2}{4\pi}} e^{-\frac{y^2}{2\pi}} \leq C_1 e^{-\frac{y^2}{2\pi}} \]

and

\[ \int_{\mathbb{R}} t^{-1/2} \left( e^{-\frac{(x-y)^2}{4\pi}} + e^{-\frac{(x-y)^2}{4\pi}} \right) dy \leq C_1, \]

we conclude that the integral in (6.29) is again bounded by \( \epsilon C_1 \theta(x, t) \).

Now, we note that if we project Green’s function on the \( R \)-component, say \( \tilde{G}_R(x, y, t) \), we obtain a better bound than that of (6.28); see (4.2). More precisely, we find

\begin{equation}
|\tilde{G}_R(x, y, t)| \leq C t^{-1/2} (t^{-1/2} + e^{-\eta |y|}) \left( e^{-\frac{(x-y)^2}{4\pi}} + e^{-\frac{|x-y|^2}{4\pi}} \right).
\end{equation}

(6.30)

Using this better bound on the \( R \)-component, the above argument shows that

\[ \int_{\mathbb{R}} \tilde{G}(x, y, t) \tilde{R}_0(y) dy \leq C \epsilon (1 + t)^{-1/2} \theta(x, t). \]
Next, for the second term in (6.27), we write
\[
\int_0^t \int_\mathbb{R} \tilde{G}(x, y, t - s)N_2(\tilde{R}, \tilde{\phi}_y, \tilde{\delta}_y)(y, s)dyds
\]
\[
= \int_0^t \left[ \int_{|y| \leq 1} + \int_{|y| > 1} \right] \tilde{G}(x, y, t - s)N_2(\tilde{R}, \tilde{\phi}_y, \tilde{\delta}_y)(y, s)dyds
\]
\[
= I_1 + I_2.
\]
By the nonlinear estimates in (6.24), we have
\[
|I_1| \leq C\left(\epsilon + h(t)^2\right) \int_0^t \int_{|y| \leq 1} \tilde{G}(x, y, t - s)e^{-\frac{c}{2}t^2} dyds
\]
and
\[
|I_2| \leq C\left(\epsilon + h(t)^2\right) \int_0^t \int_{|y| \leq 1} \tilde{G}(x, y, t - s)\left[e^{-\eta|s|} + (1 + s)^{-1+\epsilon}\right] \theta(y, s) dyds.
\]
The following lemma is precisely to give the convolution estimates on the right-hand sides of the above inequalities, and so the desired bound for the second term in (6.27) is obtained. We shall prove the lemma in the appendix.

**Lemma 6.5.** For some C and M sufficiently large,
\[
\left| \int_0^t \int_{|y| \leq 1} \partial_y^k \tilde{G}(x, y, t - s)e^{-\frac{c}{2}t^2} dyds \right| \leq C(1 + t)^{-\frac{k}{2}} \theta(x, t),
\]
\[
\left| \int_0^t \int_{|y| \leq 1} \partial_y^k \tilde{G}(x, y, t - s)\left[e^{-\eta|s|} + (1 + s)^{-1+\epsilon}\right] \theta(y, s) dyds \right| \leq C(1 + t)^{-\frac{k}{2}} \theta(x, t),
\]
for k = 0, 1, 2, 3. In addition, similar estimates hold for \(\Pi_R \tilde{G}(x, y, t)\), with a gain of an extra factor \((1 + t)^{-1/2}\).

Next, we consider the last integral term
\[
\sum_\pm \int_0^t \int_\mathbb{R} (E^k(x, t - s) - E^k(x, t)) \langle \Psi^k(y), N_2(\tilde{R}, \tilde{\phi}_y, \tilde{\delta}_y)(y, s)\rangle dyds
\]
in (6.27). Due to the bounds (6.21) and (6.23), we obtain
\[
|\langle \Psi^k(y), N_2(\tilde{R}, \tilde{\phi}_y, \tilde{\delta}_y)(y, s)\rangle| \leq C e^{-\frac{c}{2}|s|^2} e^{-\eta|s|} [\epsilon + h(s)^2].
\]
Also, due to (6.5) and the fact that \(B(x, t) = e(x, t + 1)\) and \(B_1(x, t) = \mathcal{O}(\theta(x, t))\), we have
\[
\left| \sum_\pm \int_0^t \int_\mathbb{R} (E^k(x, t - s) - E^k(x, t)) \langle \Psi^k(y), N_2(\tilde{R}, \tilde{\phi}_y, \tilde{\delta}_y)(y, s)\rangle dyds \right|
\]
\[
\leq C(\epsilon + h(t)^2) \int_0^t |e(x, t - s + 1) - e(x, t + 1)| e^{-\eta|s|} ds + C(\epsilon + h(t)^2) \theta(x, t).
\]

(6.32)

Again note that the \(R\)-component of \(E^k(x, t)\) is bounded by
\[
C|B(x, t)| + Ce^{-\eta|s|}|B(x, t)|.
\]
Estimate (6.32) is thus improved by either \((1 + t)^{-1/2}\) or \(e^{-\eta|s|}\) when projected on the \(R\)-component. Here recall again that \(e^{-\eta|s|}\theta(x, t)\) is in fact decaying exponentially in time and space. Thus it suffices to show that the integral on the right-hand side of (6.32) is bounded by \(C\theta(x, t)\). We shall prove the following lemma in the appendix.
Lemma 6.6. For each sufficiently large $M$, there is a constant $C$ so that

$$\int_0^1 |\partial_x^k\left[ e(x, t - s + 1) - e(x, t + 1) \right]| e^{-\eta t} \, dt \leq C(1 + t)^{-k/2} \theta(x, t),$$

for $k = 0, 1, 2, 3$.

In summary, collecting all these estimates into (6.27), we have obtained the desired estimate (6.25) for $k = 0$. The estimates for the derivatives follow exactly the same way as done above with a gain of time decay; we omit the proof.

6.9. Estimates for $h_2(t)$

In this section we prove the claimed estimate for $h_2(t)$, stated in proposition 6.4. Estimates (6.25) almost prove the claimed inequality, except the estimate near the core $x = 0$ at which $r(0) = 0$. By lemma 2.1, we can assume without loss of generality that there exist positive constants $a, b_1, b_2$ so that

$$|r(x)| \geq a, \quad \forall |x| \geq 1,$$

and

$$b_1 |x| \leq |r(x)| \leq b_2 |x|, \quad |r_x(x)| \geq b_1, \quad |r_{xx}| \leq b_2 |r(x)|. \quad (6.33)$$

Away from the core $|x| \geq 1$. In this case, the second estimate in (6.25), (6.26), and the fact that $|r_x| + |r_{xx}| = O(e^{-\eta|x|})$ imply

$$|\hat{\phi}(x, t)| \leq C a^{-1} \left( \epsilon + h(t) \right) \theta(x, t)$$

and

$$|\hat{\phi}_x(x, t)| + |\hat{\phi}_{xx}(x, t)| \leq C a^{-1} \left( \epsilon + h(t)^2 \right) \left( e^{-\eta|x|} + (1 + t)^{-1/2} \right) \theta(x, t).$$

Note that $e^{-\eta|x|} \theta(x, t)$ can be bounded by $Ce^{-\eta|x+\epsilon t|}$, which may be neglected. In addition, from (6.2) we can write $r \hat{\phi}_t$ in terms of $(\tilde{R}, \hat{\phi})$ and their spatial derivatives. Thus, $r \hat{\phi}_t$ is bounded by $C(\epsilon + h(t)^2)(1 + t)^{-1/2} \theta(x, t)$, where the extra $(1 + t)^{-1/2}$ is due precisely to the fact that the right-hand side of (6.2) does not contain $\hat{\phi}$ (the term with the slowest decay in the equation). Therefore, the claimed estimate on $\hat{\phi}$ follows when $|x| \geq 1$.

Near the core $|x| \leq 1$. The second estimate in (6.26) with $k = 2$ gives

$$|(r \hat{\phi}_{xx})_x| \leq C (\epsilon + h(t)^2) \theta(x, t).$$

Thus, if we write

$$(r^2 \hat{\phi}_x)_x = r(2r_x \hat{\phi}_x + r \hat{\phi}_{xx}) = r((r \hat{\phi})_{xx} - r_{xx} \hat{\phi}),$$

by integration together with (6.33) we have

$$|r^2 \hat{\phi}_x| = \left| \int_0^x r(y) (g(y, t) - r_{yy} \hat{\phi}) \, dy \right| \leq C \int_0^x |y| (|(r \hat{\phi})_{yy}| + |r \hat{\phi}|) \, dy$$

$$\leq C x^2 \left( \epsilon + h(t)^2 \right) e^{-\eta t}.$$

Here we note that $r \hat{\phi}(x, t)$ is finite, and so $r^2 \hat{\phi}$ vanishes at $x = 0$ since $r(0) = 0$. Again by the estimate $|r(x)| \geq b_1 |x|$ from (6.33), the above estimate yields

$$|\hat{\phi}_x| \leq C b_1^{-1} \left( \epsilon + h(t)^2 \right) e^{-\eta t}.$$ (6.34)
In addition, if we write \( r \tilde{\phi} = (r \tilde{\phi})_x - r \tilde{\phi}_x \) and use the above estimate together with (6.25), we obtain the claimed estimate for \( \tilde{\phi} \) at once thanks to the assumption that \(|r_x| \geq b_1 > 0\).

Similarly, we check the claimed estimate for \( \tilde{\phi}_t \). As above, we write

\[
(r^2 \tilde{\phi}_{xx})_x = r \left( (r \tilde{\phi})_{xx} - \frac{r_{xx}}{r} (r \tilde{\phi})_x \right)
\]

and note that, by (6.26), \(|r \tilde{\phi})_{xx} \) and \(|r \tilde{\phi})_x \) are already bounded by \( C_e + h(t)^2 \) \( e^{-\eta t} \). It thus follows similarly to (6.34) that

\[
|\tilde{\phi}_{xx}| \leq C b_1^{-2} \left( \epsilon + h(t)^2 \right) e^{-\eta t}.
\]

This yields the desired estimate for \( \tilde{\phi}_t \) near the core by writing \( r \tilde{\phi}_t = (r \tilde{\phi})_{xx} - r \tilde{\phi}_x \) and using (6.26).

Finally, we turn to the claimed estimate for \( \tilde{\phi}_{xx} \). First notice that we also have an estimate for \( r \tilde{\phi}_{xx} \) for all \( x \) by writing

\[
(r^3 \tilde{\phi}_{xx})_x = r^2 \left( (r \tilde{\phi})_{xxx} - 3r_{xx} \tilde{\phi}_x - r_{xxx} \tilde{\phi} \right)
\]

and integrate the identity from 0 to \( x \). By (6.35), \( r \tilde{\phi}_{xx} \) is finite at \( x = 0 \) and so \( r^3 \tilde{\phi}_{xx} \) vanishes at \( x = 0 \). Using estimate (6.26) with \( k = 3 \) and the estimates on \( \tilde{\phi}_x \) and on \( \phi \), we thus obtain

\[
|r^3 \tilde{\phi}_{xx}| \leq C \int_0^x |y|^2 \left( (r \tilde{\phi})_{yyy} - 3ry \tilde{\phi}_y - r_{yyy} \tilde{\phi} \right) dy \leq C |x|^3 \left( \epsilon + h(t)^2 \right) e^{-\eta t}.
\]

Again, since \(|r(x)| \geq b_1 |x| \), we then obtain

\[
|\tilde{\phi}_{xx}| \leq C b_1^{-3} \left( \epsilon + h(t)^2 \right) e^{-\eta t},
\]

for \(|x| \leq 1\).

This completes the proof of the claimed estimate \( h_2(t) \). The key proposition (proposition 6.4) is therefore proved, and so is the main theorem.

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Appendix A. Convolution estimates

In this section, we prove the convolution estimates that we used in the previous sections. These estimates can also be found in [BNSZ12].

Proof of lemma 6.5. We recall that

\[
|\tilde{\mathcal{G}}(x, y, t)| \leq Ct^{-1/2} \left( e^{-\frac{\|x-y\|^2}{2t}} + e^{-\frac{\|x+y\|^2}{2t}} \right),
\]

and

\[
|\tilde{\mathcal{G}}_R(x, y, t)| \leq Ct^{-1/2} (t^{-1/2} + e^{-\|y\|^2}) \left( e^{-\frac{\|x+y\|^2}{2t}} + e^{-\frac{\|x-y\|^2}{2t}} \right).
\] (A.1)
We will show that
\[
\begin{align*}
\int_0^t \int_{\mathbb{R}} \tilde{g}(x, y, t-s)(1+s)^{-1+\varepsilon}\theta(y, s) \, dy \, ds & \leq C(1+t)^{\varepsilon}\theta(x, t) \\
\int_0^t \int_{\mathbb{R}} \tilde{g}_R(x, y, t-s)(1+s)^{-1+\varepsilon}\theta(y, s) \, dy \, ds & \leq C(1+t)^{-1/2+\varepsilon}\theta(x, t).
\end{align*}
\tag{A.2}
\]

We start with a proof of the first estimate in (A.2). We first note that there are constants \(C_1, \tilde{C}_1 > 0\) such that
\[
\tilde{C}_1 e^{-y^2/M} \leq |\theta(y, s)| \leq C_1 e^{-y^2/M}
\]
for all \(0 \leq s \leq 1\). Thus, for some constant \(C_1\) that may change from line to line, we have
\[
\int_0^t \int_{\mathbb{R}} |\tilde{g}(x, y, t-s)(1+s)^{-1+\varepsilon}\theta(y, s)(y, s)dy \, ds \leq C_1 \int_0^t \int_{\mathbb{R}} (t-s)^{-1/2}e^{-\frac{(t-s)^2}{M}} e^{-y^2/M} dy \, ds.
\]
Using that the maximum of the quadratic polynomial \(\alpha x^2 + \beta x + \gamma\) is \(-\beta^2/(4\alpha) + \gamma\), it is easy to see that the sum of the first two terms in (A.4), which involve only
\( x \) and not \( y \), is less than or equal to zero. Omitting this term, we therefore obtain the estimate

\[
\exp \left( \frac{(x \pm c_\theta t)^2}{M(1+t)} \right) \leq \exp \left( -\frac{4(t-s) + Ms}{4\kappa M(1+s)(t-s)} \right) \left( \frac{y - \delta_1 c_\theta (1 + s)}{4(t-s) + M(1+s)} \right) \]

(A.5)

for \( \delta_1 = \pm 1 \). Using this result, we can now estimate the integral (A.2). Indeed, we have

\[
\theta(x, t)^{-1} \int_0^t \int_{\mathbb{R}} |\tilde{G}(x, y, t-s)| (1 + s)^{-1+\epsilon} \theta(y, s) dy ds
\]

\[
\leq C_t(1 + t)^{1/2} \int_0^t \frac{1}{\sqrt{t-s(1+s)^{1/2}}} \left( \left| y - \frac{y}{(1+\kappa_1 t)(1+s)} \right| \right)^2 dy ds
\]

\[
\leq C_1(1 + t)^{1/2} \int_0^t \frac{1}{\sqrt{t-s(1+s)^{1/2}}} \frac{4\kappa_1 M(1+s)(t-s)}{4(t-s) + M(1+s)} ds
\]

\[
\leq C_1(1 + t)^{1/2} \int_0^t \frac{1}{(1+s)^{1/2}} \frac{1}{(1+t)^{1/2}} ds + C_1(1 + t)^{1/2} \int_1^t \frac{1}{(1+t)^{3/2}} ds
\]

\[
\leq C_1(1 + t)^{-1/2} + C_1,
\]

which is bounded since \( \kappa > 0 \). This proves the first estimate in (A.2). The second estimate is entirely the same, using the refined estimate \( A.1 \) for \( \tilde{G}_R \). Also, derivative estimates follow very similarly. We omit these further details.

Finally, it remains to show that

\[
\int_0^t \int_{|y| \leq 1} |\tilde{G}(x, y, t-s)| e^{-\frac{y^2}{2\delta_2^2}} dy ds \leq C \theta(x, t),
\]

(A.6)

where the Green function bounds read

\[
|\tilde{G}(x, y, t)| \leq C t^{-1/2} \left( e^{-\frac{(x-y)^2}{2\delta_2^2}} + e^{-\frac{|y|^2}{2\delta_2^2}} \right),
\]

for \( |y| \leq 1 \). Estimate (A.6) is clear when \( 0 \leq t \leq 1 \). We consider the case \( t \geq 1 \). The proof of this estimate uses the following bound:

\[
e^{\frac{(x-y)^2}{2\delta_2^2}} e^{-\frac{|y|^2}{2\delta_2^2}} \leq Ce^{-\frac{(t-1)^2}{2\delta_2^2}},
\]

for fixed constant \( B \) and for large \( M \). This is a simpler version of (A.5); see also (A.7). We thus have

\[
\theta(x, t)^{-1} \int_0^t (t-s)^{-1/2} e^{-\frac{(x-y)^2}{2\delta_2^2}} e^{-\frac{|y|^2}{2\delta_2^2}} ds
\]

\[
\leq C(1 + t)^{1/2} \int_0^t (t-s)^{-1/2} e^{\frac{t}{2\delta_2^2}} e^{-\frac{|y|^2}{2\delta_2^2}} ds
\]

\[
\leq C(1 + t)^{1/2} \left[ 2^{1/2} \int_0^t e^{-|y|^2} ds + e^{\frac{t}{2\delta_2^2}} \int_1^t (t-s)^{-1/2} ds \right]
\]

\[
\leq C(1 + t)^{1/2} t^{-1/2} + C(1 + t) e^{\frac{t^2}{2\delta_2^2}},
\]

which is bounded for \( t \geq 1 \). This proves estimate (A.6), and completes the proof of lemma 6.5. \( \square \)
Proof of lemma 6.6. We need to show that
\[
\int_0^t \left[ |e(x, t - s + 1) - e(x, t + 1)| \right] e^{-\frac{\tau^2}{2}} ds \leq C_1 \theta(x, t).
\]
Intuitively, this integral should be small for the following reason. The difference \(e(x, t - s) - e(x, t + 1)\) converges to zero as long as \(s\) is not too large, say on the interval \([0, t/2]\). For \(s \in [t/2, t]\), on the other hand, we use the exponential decay in \(s\). Indeed, we have
\[
|e(x, t - s + 1) - e(x, t + 1)|
= \left| \int_{t-s+1}^{t+1} e_x(x, \tau) d\tau \right|
\leq C \int_{t-s+1}^{t+1} \left( \frac{e^{-\frac{(\xi-x)^2}{4\tau}} + e^{-\frac{(\xi-x)^2}{4\tau}}}{\tau} \right) \left( \frac{e^{-\frac{(\xi-x)^2}{4\tau}} - e^{-\frac{(\xi-x)^2}{4\tau}}}{\tau} \right) d\tau
\leq C \left| \int_{t-s+1}^{t+1} \left( \frac{1}{\sqrt{\tau}} + \frac{1}{\tau} \right) \left( e^{-\frac{(\xi-x)^2}{4\tau}} + e^{-\frac{(\xi-x)^2}{4\tau}} \right) d\tau \right|
\]
where the last estimate follows by the fact that \(ze^{-z^2}\) is bounded for all \(z\).

We shall give an estimate for \(\theta^{-1}(x, t)(e(x, t - s + 1) - e(x, t + 1))\). For instance, we consider the single exponential term
\[
e^{\frac{(\xi-x)^2}{4\tau}} e^{-\frac{(\xi-x)^2}{4\tau}}.
\]
By combining these and completing the square in \(x\), the terms in the exponential become
\[
\frac{M(t - \tau + 1) + (M - B)\tau}{MB(t+1)} \left[ x + \frac{c_\theta(B - M)\tau(t + 1)}{M(t - \tau + 1) + (M - B)\tau} \right]^2 + \frac{c_\theta^2(t - \tau + 1)^2}{M(t - \tau + 1) + (M - B)\tau}.
\]
Since \(\tau \leq t\), if \(B\) is some fixed constant and \(M\) is sufficiently large we can neglect the exponential in \(x\). That is, we have
\[
e^{\frac{(\xi-x)^2}{4\tau}} e^{-\frac{(\xi-x)^2}{4\tau}} \leq Ce^{\frac{\xi^2}{4\tau}}.
\]
We therefore obtain
\[
\theta^{-1}(x, t)|e(x, t - s + 1) - e(x, t + 1)| \leq C(1 + t)^{1/2} \int_{t-s+1}^{t+1} \left( \frac{1}{\sqrt{\tau}} + \frac{1}{\tau} \right) e^{\frac{\xi^2}{4\tau}} d\tau
\leq C(1 + t)^{1/2}(1 + t - s)^{-1/2} e^{\frac{\xi^2}{4\tau}}.
\]
Using this and taking \(M\) large and \(\eta = c_\theta^2/M\), we obtain
\[
\theta(x, t)^{-1} \int_0^t |e(x, t - s + 1) - e(x, t + 1)| e^{-\frac{\eta^2}{2}} ds
\leq C(1 + t)^{1/2} \int_0^t (1 + t - s)^{-1/2} e^{\frac{\xi^2}{4\tau}} e^{-2\eta^2} ds
\leq C(1 + t)^{1/2} \left[ (1 + t)^{-1/2} \int_0^t e^{-\eta^2} ds + e^{-\eta^2} \int_{t/2}^t (1 + t - s)^{-1/2} ds \right]
\leq C.
\]
This proves the lemma for the case \(k = 0\). The derivative estimates follow easily from the above proof. □
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