Can large inhomogeneities generate target patterns?

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Abstract. We study the existence of target patterns in oscillatory media with weak local coupling and in the presence of an impurity, or defect. We model these systems using a viscous eikonal equation posed on the plane and represent the defect as a perturbation. In contrast to previous results, we consider large defects, which we describe using a function with slow algebraic decay, i.e., $g \sim O(1/|x|^m)$ for $m \in (1, 2]$. We prove that these defects are able to generate target patterns and that, just as in the case of strongly localized impurities, their frequency is small beyond all orders of the small parameter describing their strength. Our analysis consists of finding two approximations to target pattern solutions, one which is valid at intermediate scales and a second one which is valid in the far field. This is done using weighted Sobolev spaces, which allow us to recover Fredholm properties of the relevant linear operators, as well as the implicit function theorem, which is then used to prove existence. By matching the intermediate and far-field approximations, we then determine the frequency of the pattern that is selected by the system.

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1. Introduction

Target patterns are coherent structures that emerge in excitable and in oscillatory media. They are characterized by concentric waves that expand away from a center, or core region, creating a ‘bull’s-eye’ pattern. Although often associated with the Belousov–Zhabotinsky reaction [27], they also appear in colonies of slime mold [3, 5], in the oxidation of carbon monoxide on platinum [26], and in brain tissue [24].

In this paper, we will focus on target patterns that arise in oscillatory media, where three key mechanisms, or processes, contribute to their formation. The first mechanism is associated with the intrinsic dynamics of the system, which must support a limit cycle that results in uniform time oscillations. The second is a transport process that allows for different spatial regions to interact, such as diffusion in chemical reactions, or coupling between neurons in brain tissue. While these two processes are enough to generate traveling and spiral waves, to obtain target patterns one needs a third ingredient, a defect. Indeed, it is believed that the role of defects, or impurities, is to alter the dynamics of the system in a localized area resulting in a change in the frequency of the time oscillations. As a consequence, these defects act as pacemakers entraining the rest of the medium and forming target patterns.

While experiments and previous analytical results confirm that small localized defects give rise to these patterns, [7, 8, 11, 13, 15, 17, 18, 22, 25, 26], in this paper we want to determine the exact level of localization that is needed to generate them. In particular, assuming the inhomogeneity is modeled as a function with algebraic decay of order $O(1/|x|^m)$, we want to determine how small we can take $m > 0$ and still obtain a well-defined target pattern.

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To simplify the analysis, we concentrate only on systems which involve weak local coupling. Because it is well known that under this assumption the amplitude of oscillations is tied, or enslaved, to the dynamics of the phase, this allows us to focus our analysis on this last variable. Indeed, the results presented in [4] show that coherent structures in these systems are well described by the following viscous eikonal equation

\[ \phi_t = \Delta \phi - b|\nabla \phi|^2 - \varepsilon g(x), \quad x \in \mathbb{R}^2, \]  

where the perturbation, \( \varepsilon g \), represents the defect. The above expression is derived using a multiple-scale analysis, and it therefore models phase changes that occur over long spatial and time scales. In this context, target patterns then correspond to solutions of the form \( \phi(x,t) = \phi(x) - \Omega t \), satisfying the boundary condition \( \nabla \phi \to k \) as \( |x| \to \infty \), where the constant \( k \) then represents the pattern’s wavenumber.

Our motivation for considering large inhomogeneities is threefold. First, in all previous work the level of localization imposed on the inhomogeneity was tied to the tools used to prove the existence of these patterns. Yet, numerical simulation like the ones presented here in Sect. 6 shows that these assumptions can be relaxed. For example, in [23] defects are modeled as functions with compact support and target pattern solutions are found using separation of variables. In contrast, in [15] the authors use spatial dynamics to prove the existence of these patterns. This then allows them to model the impurities as radially symmetric functions with exponential decay. In [13], thanks to the use of weighted Sobolev spaces, this assumption is relaxed and general (non-radially symmetric) inhomogeneities with decay of order \( O(1/|x|^m) \), \( m > 2 \), are considered.

Although using different approaches, the references mentioned above show that target patterns can only be generated by inhomogeneities with a positive and finite mass \( M = \int_{\mathbb{R}^2} g \). This obviously restricts the level of decay of \( g \) to be of order \( o(1/|x|^2) \). However, our numerical simulations show that one can obtain target patterns even in the case when the defect is assumed to decay only at order \( O(1/|x|^m) \), for \( m \in (1, 2] \). We are therefore interested in proving the existence of target patterns for these ‘large’ inhomogeneities of infinite mass.

Our second reason for considering this problem is tied to the existence of spiral waves in oscillatory media with nonlocal coupling. In [10], it was shown that the dynamics of these patterns are well described by the following amplitude equation

\[ 0 = K * w + \lambda w + \alpha |w|^2 w + \mathcal{O}(\varepsilon), \quad w \in \mathbb{C}, \quad 0 < \varepsilon << 1 \]

where \( w \) is a radially symmetric complex-valued function and \( K \) is a symmetric convolution kernel of diffusive type. Additional assumptions on \( K \) imply that formally one can write this operator as \( (1 - \varepsilon \Delta_1)^{-1} \sigma \Delta_1 \), and suggest preconditioning the above equation with \( (1 - \varepsilon \Delta_1) \), where \( \Delta_1 = \partial_r r + \frac{1}{r} \partial_r - \frac{1}{r^2} \). This then results in the following expression, which perhaps not surprisingly resembles the complex Ginzburg–Landau equation,

\[ 0 = \beta \Delta_1 w + \lambda w + \alpha |w|^2 w + \mathcal{O}(\varepsilon), \quad \beta = (\sigma - \varepsilon \lambda), \quad r \in [0, \infty). \]

From there, a similar multiple-scale analysis as the one carried out in [4] and that we also summarize in Appendix, gives a hierarchy of equations at different powers of a small parameter \( \delta = \delta(\varepsilon) \). In particular, at order \( \delta^2 \) one finds the steady-state viscous eikonal equation,

\[ -\Omega = \Delta_1 \phi - b|\nabla \phi|^2 - \delta^2 g, \]  

as a description of the phase dynamics of spiral waves. However, in contrast to the case of target patterns, here the inhomogeneity does not represent a defect, but is instead related to the small variations, \( \rho \), of the amplitude of the pattern. More precisely, \( g \sim (1 - \rho^2) \). Although not immediately obvious, one can check using numerical simulations that the perturbation \( (1 - \rho^2) \) decays at infinity at order \( \mathcal{O}(1/|x|^2) \) (see Appendix A). Therefore, the particular viscous eikonal equation that is connected to the phase dynamics of spiral waves in these systems is the same equation that we are trying to solve.
Finally, our third motivation comes from the following change of variables, \( \phi = -\frac{1}{\gamma} \log(\Psi) \), which transforms the steady-state viscous eikonal equation (2), into a Schrödinger eigenvalue problem with potential \( \varepsilon g \).

\[
\Omega \Psi = \Delta \Psi + \varepsilon g(x) \Psi.
\]

The transformation also shows that our target pattern solutions correspond to bound states of this operator. The only result solving the above eigenvalue problem that we are aware of is that of Simon [21], who proved that in the two-dimensional case and under the assumption of localized potentials, i.e.,

\[
\int g(x)(1 + x^2) \, dx < \infty,
\]

bound states exists if and only if the mass \( \int g(x) \, dx > 0 \).

Notice that in the context of the Schrödinger operator, our problem corresponds to the ‘supercritical’ case, in the sense that the potential, \( g \), no longer corresponds to a bounded perturbation of the Laplacian. To see this, fix \( g(x) = 1/(1 + |x|^m) \) with \( 1 < m \leq 2 \), and consider the rescaling \( y = \gamma x \). The Schrödinger operator then reads \( \gamma^2 \Delta_y \Psi + \varepsilon (\gamma^m g)^\gamma \Psi \), and it is then clear that if we choose \( \gamma \) small, the potential is actually ‘large’ in the far field. Consequently, the results from [21] no longer apply for the case considered here.

In this paper, we show that target pattern solutions to the viscous eikonal equation, or equivalently, bound states to the above Schrödinger operator, exists even for these large inhomogeneities. As with small defects, we prove that target pattern solutions have frequencies, \( \Omega \), that are small beyond all orders of the parameter \( \varepsilon \). Consequently, one cannot use a regular perturbation expansion to justify existence. To resolve this issue, we first find two approximations to target patterns, one which is valid at intermediate scales and second one that accounts for the far-field behavior of the solution. By matching these two approximations, we are then able to determine the unique value of the frequency selected by the system.

It is in the course of this analysis that one sees that the slow decay rate of the inhomogeneity plays a major role in shaping the solution at intermediate scales. This is the main difference between the analysis presented here and that of [13], where inhomogeneities of finite mass are considered. It is also why we will split defects into a core region and a far-field region, reflecting the fact that the defects we work with are still too small to alter the shape of the pattern at large scales, but do contribute to the form of the equation at intermediate scales. In particular, we write the impurity as the sum two functions, defined as

\[
g_c = (1 - \chi^D)g \quad g_f = \chi^D g,
\]

where \( \chi^D \) is a \( C^\infty \) radial cut-off function, with \( \chi^D(|x|) = 0 \) for \( |x| < D \) and \( \chi^D(|x|) = 1 \) for \( |x| > 2D \). To prove the existence of target patterns, the value of the parameter \( D \) can remain arbitrary, so long as it is a finite number. This follows because even though in the above definition the function \( g_c \) has compact support, our results hold for more general ‘core’ functions. The only requirement being that this core defect has finite mass. We therefore make the following assumption.

**Hypothesis 1.1.** The inhomogeneity, \( g \), lives in \( H^k_s(\mathbb{R}^2) \), with \( k \geq 2 \) and \( \sigma \in (0, 1) \), is radially symmetric, and is positive. In addition, the defect can be split into the sum of two positive functions, \( g_c, g_f \), satisfying

- The function \( g_f \) is in \( H^k_s(\mathbb{R}^2) \) for \( 0 < \sigma < 1 \). In particular, \( g_f \sim O(1/r^m) \) as \( r \to \infty \), with \( 1 < m \leq 2 \), while near the origin \( g_f(|x|) = 0 \) for \( |x| < 1 \).
- The function \( g_c \) is in \( H^k_s(\mathbb{R}^2) \) for \( \gamma > 1 \). In particular, \( g_c \sim O(1/r^d) \) with \( d > 2 \) as \( r \to \infty \).

**Remark 1.2.** The spaces \( H^k_s(\mathbb{R}^2) \), with \( \sigma \in \mathbb{R} \), are weighted Sobolev spaces with norm

\[
\|u\|_{H^k_s(\mathbb{R}^2)} = \sum_{|\alpha| \leq k} \|(1 + |x|^2)^{\sigma/2}D^\alpha u(x)\|_{L^2(\mathbb{R}^2)}.
\]

Notice that for positive values of \( \sigma \), they impose a level of decay on functions. For a precise definition of these spaces, see Sect. 2.
Theorem 1. Let $k \geq 2$ and $\sigma \in (0, 1)$ and consider a function $g \in H_{r,\sigma}^k(\mathbb{R}^2)$ satisfying Hypothesis 1.1. Then, there exist a constant $\varepsilon_0 > 0$ and a $C^1([0, \varepsilon_0])$ family of eigenfunctions $\phi = \phi(r; \varepsilon)$ and eigenvalues $\Omega = \Omega(\varepsilon)$ that bifurcate from zero and solve the equation

$$\Delta_0 \phi - b(\partial_r \phi)^2 - \varepsilon g(r) + \Omega = 0 \quad r = |x| \in [0, \infty).$$

Moreover, this family has the form

$$\phi(r; \varepsilon) = -\frac{1}{b} \chi_1(\Lambda r) \log(K_0(\Lambda r)) + \phi_1(r; \varepsilon) + \varepsilon c, \quad \Lambda^2 = b\Omega(\varepsilon)$$

where

i) $c$ is a constant that depends on the initial conditions of the problem,

ii) $K_0(z)$ represents the zeroth-order modified Bessel function of the second kind,

iii) $\partial_r \phi_1 \in H_{r,\sigma}^k(\mathbb{R}^2)$, and

iv) $\Omega = \Omega(\varepsilon) = C(\varepsilon) 4e^{-2\gamma_c} \exp[2/a]$, with

$$a = -\varepsilon b \int_0^\infty g_c(r) r \, dr,$$

and $C(\varepsilon)$ a $C^1$ constant that depends on $\varepsilon$.

Remark 1.3. Notice that under Hypothesis 1.1, the viscous eikonal equation, (1), is invariant under rotations. As a result, we can look for solutions that are radially symmetric. This assumption is made mainly for convenience, and one can follow the steps in [13] to tackle the more general case of non-symmetric inhomogeneities.

Remark 1.4. If the inhomogeneity $g$ has strong algebraic decay, i.e., $g(r) \sim 1/r^m$ with $m > 2$, then we are back in the regime considered in [13]. In this case, the impurity has finite mass and there is no need to split this function into the sum of its core and far-field functions. In fact, one can set $g = g_c = g_f$ and the above theorem is equivalent to Theorem 1 in [13] with $a = -\varepsilon b \int_0^\infty g(r) r \, dr < \infty$.

Remark 1.5. While the exact form of the cut-off function $\chi_D$ appearing in the definition of $g_c$ is not important for the proof of existence, it does play a role when approximating the pattern’s frequency, $\Omega$. As our numerical simulations show, there is an optimal way of picking the parameter $D$ that allows one to obtain better estimates for the frequency, see Sect.6. If a non-optimal choice is made, one can improve the estimates for $\Omega$ by using higher-order approximations for the intermediate and far-field solutions when carrying out the matched asymptotics, see Sect.5.
1.1. Outline

The paper is organized as follows. In Sect. 2, we introduce a special class of weighted Sobolev spaces and summarize Fredholm properties for the Laplacian and related operators. In Sect. 3, we work with our model (1) and derive from it an equation that is valid at intermediate scales. We then prove existence of solutions to this equation that are bounded near the origin and that have appropriate growth conditions. Next, in Sect. 4 we work with the full model (1) and, treating the frequency as a parameter, find a first-order approximation to target pattern solutions. Then, in Sect. 5.1 we use matched asymptotics to determine the value of the frequency, Ω, selected by the system. More precisely, we show that Ω is a $C^1$ function of the parameter $\varepsilon$. This then allows us to prove existence of solutions using the implicit function theorem. The analysis is complemented by numerical simulations presented in Sect. 6 and a discussion in Sect. 7.

2. Preliminaries

In this section, two different classes of Sobolev spaces are introduced, weighted Sobolev spaces and Kondratiev spaces. We also look at Fredholm properties for the specific operators that will appear in later sections. We will see how these properties depend on the weighted spaces used to define the domain and range of these operators. Throughout this section, we use the symbol $\langle x \rangle = (1 + |x|^2)^{1/2}$, which appears in the definition of the norms for the weighted Sobolev spaces introduced.

2.1. Weighted Sobolev spaces

Let $s$ be a nonnegative integer, $p \in (1, \infty)$, and $\gamma$ a real number. We denote by $W^{s,p}_\gamma(\mathbb{R}^d)$ the space of functions formed by taking the completion of $C_\infty^0(\mathbb{R}^d, \mathbb{C})$ under the norm

$$\|u\|_{W^{s,p}_\gamma(\mathbb{R}^d)} = \sum_{|\alpha| \leq s} \|\langle x \rangle^\gamma D^\alpha u(x)\|_{L^p(\mathbb{R}^d)}.$$ 

When $p = 2$, we let $W^{s,2}_\gamma(\mathbb{R}^d) = H^s_\gamma(\mathbb{R}^d)$. In this case, these spaces are also Hilbert spaces, with inner product defined in the natural way by

$$\langle f, g \rangle = \sum_{|\alpha| \leq s} \int_{\mathbb{R}^d} f(x) \bar{g}(x) \langle x \rangle^{2\gamma} \, dx$$

where the overbar denotes the complex conjugate.

Notice in particular that depending on the sign of the weight $\gamma$, the functions in these spaces are either allowed to grow ($\gamma < 0$), or forced to decay ($\gamma > 0$). We also have natural embeddings, with $W^{s,p}_\gamma(\mathbb{R}^d) \subset W^{s,p}_\sigma(\mathbb{R}^d)$ provided $\gamma > \sigma$, and $W^{s,p}_\gamma(\mathbb{R}^d) \subset W^{k,p}_\gamma(\mathbb{R}^d)$ whenever $s > k$. For $1 < p < \infty$ we can also identify the dual, $(W^{s,p}_\gamma(\mathbb{R}^d))^*$, with the space $W^{-s,q}_\gamma(\mathbb{R}^d)$, where $p$ and $q$ are conjugate exponents.

Kondratiev Spaces: With $s, p,$ and $\gamma$ as in the previous section, we define Kondratiev spaces as the completion of $C_0^\infty(\mathbb{R}^d, \mathbb{C})$ functions under the norm

$$\|u\|_{M^{s,p}_\gamma(\mathbb{R}^d)} = \sum_{|\alpha| \leq s} \|\langle x \rangle^{\gamma+|\alpha|} D^\alpha u(x)\|_{L^p(\mathbb{R}^d)}$$

and denote them by the symbol $M^{s,p}_\gamma(\mathbb{R}^d)$.

Again, we see that these spaces are Hilbert spaces when $p = 2$, with inner product given by

$$\langle f, g \rangle = \sum_{|\alpha| \leq s} \int_{\mathbb{R}^d} f(x) \bar{g}(x) \langle x \rangle^{2(\gamma+|\alpha|)} \, dx.$$
We also have the following natural embeddings. One can check that \( M_{\gamma}^{s,p}(\mathbb{R}^d) \subset M_{\sigma}^{s,p}(\mathbb{R}^d) \) whenever \( s > k \), and \( M_{\gamma}^{s,p}(\mathbb{R}^d) \subset M_{\gamma}^{s,p}(\mathbb{R}^d) \) provided \( \gamma > \sigma \). In addition, as in the case of standard Sobolev spaces, one can identify the dual space \( (M_{\gamma}^{s,p}(\mathbb{R}^d))^* \) with \( M_{-\gamma,q}^{s,q}(\mathbb{R}^d) \), where \( p \) and \( q \) are conjugate exponents.

As was the case with the weighted Sobolev spaces defined above, Kondratiev spaces encode growth or decay depending on the sign of \( \gamma \). However, in contrast to \( W_{\gamma}^{s,p}(\mathbb{R}^d) \), Kondratiev spaces enforce a specific algebraic growth or algebraic decay depending on the value of \( \gamma \). In addition, we have the following which summarizes decay properties for functions in \( M_{\gamma}^{s,p}(\mathbb{R}^d) \) in terms of \( \gamma \).

**Lemma 2.1.** Let \( d, m \in \mathbb{N} \) with \( m > (d-1)/2 \). Then, for \( \gamma > -d/2 \), and for all \( f \in M_{\gamma}^{m+1,2}(\mathbb{R}^d) \) there is a constant \( C > 0 \) such that

\[
|f(x)| \leq C \|\nabla f\|_{M_{\gamma}^{m+1,2}(\mathbb{R}^d)} \|f\|_{M_{\gamma}^{m+1,2}(\mathbb{R}^d)} |x|^{-\gamma-d/2},
\]

whenever \( |x| \) is large and where \( \delta = (d-1)/2m \).

**Proof.** Let \((\theta, r)\) represent spherical coordinates in \( \mathbb{R}^d \), with \( r \) being the radial direction and \( \theta \) representing the coordinates in the unit sphere, \( \Sigma \). Then,

\[
\int_{\Sigma} |f(\theta, R)|^2 d\theta \leq \int_{\Sigma} \left( \int_{0}^{R} |\partial_r f(\theta, s)| ds \right)^2 d\theta
\]

\[
\|f(\cdot, R)\|_{L^2(\Sigma)} \leq \left[ \int_{\Sigma} \left( \int_{0}^{R} s^{\alpha} s^{\gamma+1} |\partial_r f(\theta, s)| s^{(d-1)/2} ds \right)^2 d\theta \right]^{1/2},
\]

where \( \alpha = -\gamma + 1 + (1-d)/2 \) and \( R \) is fixed. Since the inner integral is squared, we also switched the bounds of integration. Then, applying Minkowski’s inequality for integrals [6, Theorem 6.19] followed by Cauchy Schwartz,

\[
\|f(\cdot, R)\|_{L^2(\Sigma)} \leq \int_{0}^{R} s^{\alpha} \left[ \int_{\Sigma} s^{2(\gamma+1)} |\partial_r f(\theta, s)|^2 s^{(d-1)} d\theta \right] ds
\]

\[
\leq \left[ \int_{0}^{R} s^{2\alpha} \right]^{1/2} \left[ \int_{\Sigma} \int_{0}^{R} s^{2(\gamma+1)} |\partial_r f(\theta, s)|^2 s^{(d-1)} d\theta ds \right]^{1/2}
\]

\[
\leq C(\gamma, d) R^{\alpha+1/2} \|\nabla f\|_{L_{\gamma+1}^{2}(\mathbb{R}^d)}
\]

where the last line holds provided \( 2\alpha + 1 < 0 \). We therefore have the following inequality

\[
\|f(\cdot, R)\|_{L^2(\Sigma)} \leq C(\gamma, d) R^{-\gamma+d/2} \|\nabla f\|_{L_{\gamma+1}^{2}(\mathbb{R}^d)}, \quad \gamma > -d/2.
\]

Suppose now that \( f \in M_{\gamma}^{m+1,2}(\mathbb{R}^d) \). Using the chain rule, we find that \( |\partial_r D_\Sigma^k f(\cdot, R)| \leq |D^{k+1} f| R^k \), where by \( D_\Sigma^k \) we mean derivatives with respect to unit sphere variables. As a result, similar calculations as the one above show that for integer values of \( k \in [1, m] \), we have that

\[
\|D_\Sigma^k f(\cdot, R)\|_{L^1(\Sigma)} \leq R^{-\gamma-d/2} \|D^{k+1} f\|_{L_{\gamma+1}^{2}(\mathbb{R}^d)}, \quad \gamma > -d/2.
\]

Next, we use a result from Adams and Fournier [2, Theorem 5.9], which states that given a domain \( \Sigma \) of dimension \( n \), and conjugate exponents, \( p, q \), satisfying \( p > 1 \) and \( mp > n, 1 \leq q \leq p \), there exist a constant \( C \) such that for all \( u \in W_{\gamma}^{m,p}(\Sigma) \)

\[
\|u\|_{L^\infty(\Sigma)} \leq C \|u\|_{W_{\gamma}^{m,p}(\Sigma)} \|u\|_{L^q(\Sigma)}^{(1-\delta)},
\]
where \( \delta = np/(np + (mp - n)q) \). Choosing \( p = q = 2, n = d - 1, m > (d - 1)/2 \), we obtain \( \delta = (d - 1)/2m \) and
\[
\| f(\cdot, R) \|_\infty \leq C(m) R^{-\gamma - d/2} \| f \|_{M_{m+1.2}^\gamma(R^d)} \| \nabla f \|_{L_{2+1}^1(R^d)},
\]
provided \( f \in M_{m+1.2}^\gamma(R^d) \).

\[ \square \]

**Remark 2.2.** Although in the definitions presented above, the spaces \( M_{r,\gamma}^s(R^d) \) and \( W_{\gamma}^{s,p}(R^d) \) consist of complex-valued functions, in what follows we will assume that all functions are real-valued.

### 2.2. Fredholm operators

In this section and throughout the paper, we use the notation \( M_{r,\gamma}^s(R^d) \) and \( H_{r,\gamma}^s(R^d) \) to denote the subspaces of radially symmetric functions in \( M_{r,\gamma}^s(R^d) \) and \( H_{r,\gamma}^s(R^d) \), respectively.

The next lemma shows that for \( \lambda > 0 \), the operator \( \partial_r + \frac{1}{r} + \lambda \) is invertible in appropriate spaces.

**Lemma 2.3.** Let \( \gamma \in \mathbb{R}, \lambda > 0 \), and \( k \in \mathbb{N} \). Then, the operator \( (\partial_r + \frac{1}{r} + \lambda) : H_{r,\gamma}^k(R^2) \rightarrow H_{r,\gamma}^{k-1}(R^2) \) has a bounded inverse.

**Proof.** A short calculation shows that the kernel of this operator is spanned by the function \( e^{-\lambda r}/r \), which is singular near the origin and is therefore not in \( L^{2}_{r,\gamma}(R^2) \). At the same time, the adjoint of this operator is given by \( -\partial_{r} + \lambda : L_{r,\gamma}^{2}(R^2) \rightarrow H_{r,\gamma}^{-1}(R^2) \), and we find that the cokernel is spanned by the function \( e^{\lambda r} \), which again is not in the space \( L^{2}_{r,\gamma}(R^2) \) no matter what the value of \( \gamma \) is. As a result, the kernel and cokernel of this operator are trivial.

To prove the result, we are left with showing that that the inverse operator
\[
L^{2}_{r,\gamma}(R^2) \rightarrow H_{r,\gamma}^{1}(R^2) \quad f(r) \mapsto u(r) = \frac{1}{r} \int_{0}^{r} e^{\lambda(s-r)} f(s) \, ds
\]
is bounded. To show that \( \| u \|_{L^{2}_{r,\gamma}(R^2)} \leq \| f \|_{L^{2}_{r,\gamma}(R^2)} \), we use the inequality
\[
\| u \|_{L^{2}_{r,\gamma}(R^2)} \leq \| u \|_{L^{2}_{r,\gamma}(B_1)} + \| u \|_{L^{2}_{r,\gamma}(R^2 \setminus B_1)},
\]
where \( B_1 \) is the unit ball in \( R^2 \). Then,
\[
\| u \|_{L^{2}_{r,\gamma}(R^2 \setminus B_1)} = \left[ \int_{r}^{\infty} \left( \int_{0}^{r} f(s) e^{\lambda(s-r)} \, ds \right) \frac{dr}{r} \right]^{1/2}
\]
\[
\leq 2^{\gamma + 1/2} \left[ \int_{1}^{\infty} \left( \int_{0}^{r} f(s) \gamma + 1/2 e^{-\lambda(r-s)} (r-s)^{\gamma+1/2} \, ds \right) \frac{dr}{r} \right]^{1/2},
\]
\[
\leq 2^{\gamma + 1/2} \left[ \int_{0}^{r} \left( \int_{0}^{r} |f(r-z)| (r-z)^{\gamma+1/2} e^{-\lambda z} \, dz \right) \frac{dz}{r} \right]^{1/2},
\]
\[
\leq 2^{\gamma + 1/2} \int_{0}^{\gamma + 1/2} e^{-\lambda z} |z|^{\gamma+1/2} \left( \int_{0}^{\infty} |f(r+z)|^2 (r+z)^{2\gamma} \, dz \right)^{1/2} \, dz,
\]
\[
\leq C_1(\gamma, \lambda) \| f \|_{L^{2}_{r,\gamma}(R^2)}
\]
where the second line follows from the fact that $s/r < 1$ and the relation $(s)^{-\sigma} (r)^{\sigma} \leq 2^{\sigma |r|} (r - s)^{|\sigma|}$. The inequality on the third line comes from using the change of coordinates $z = r - s$ and extending the outer limits of integration to zero, while the fourth line follows from an application of Minkowski’s inequality for integrals, [6, Theorem 6.19]. In the final result, we let $C_1(\gamma, \lambda) = 2^{\gamma + 1/2} \int_0^\infty e^{-\lambda z(z)} |\gamma + 1/2| \, dz$.

To prove the relation $\|u\|_{L^2_{r,\gamma}(B_1)} \leq C \|f\|_{L^2_{r,\gamma}(\mathbb{R}^2)}$, we bound

$$|u(r)| \leq \frac{1}{r} \int_0^r \left| f(s) e^{-\lambda (r-s)} s \right| \, ds$$

$$\leq \frac{1}{r} \int_0^r |f(s)| s \, ds$$

$$\leq \frac{1}{r} \left( \int_0^r |f(s)|^2 s \, ds \right)^{1/2} \left( \int_0^r s \, ds \right)^{1/2}$$

$$\leq \frac{1}{r} \left( \int_0^\infty |f(s)|^2 (s^{2\gamma} s \, ds \right)^{1/2} \left( \frac{r}{\sqrt{2}} \right)$$

$$\leq \frac{1}{\sqrt{2}} \|f\|_{L^2_{r,\gamma}(\mathbb{R}^2)}$$

where the third line follows from Hölder’s inequality. We then obtain that $\|u(r)\|_{L^2_{r,\gamma}(B_1)} \leq C_2(\gamma) \|f\|_{L^2_{r,\gamma}(\mathbb{R}^2)}$ with $C_2(\gamma) = \frac{1}{\sqrt{2}} \left( \int (r)^{2\gamma} r \, dr \right)^{1/2}$, and consequently

$$\|u(r)\|_{L^2_{r,\gamma}(\mathbb{R}^2)} \leq (C_1(\gamma, \lambda) + C_2(\gamma)) \|f\|_{L^2_{r,\gamma}(\mathbb{R}^2)}.$$  

Next, to show that the derivative $\partial_r u \in L^2_{r,\gamma}(\mathbb{R}^2)$, we use the equation to write $\partial_r u = f - \lambda u - \frac{u}{r}$, and thus obtain

$$\|\partial_r u\|_{L^2_{r,\gamma}(\mathbb{R}^2 \setminus B_1)} \leq (1 + (1 + \lambda)C_1(\gamma, \lambda)) \|f\|_{L^2_{r,\gamma}(\mathbb{R}^2)}.$$  

To bound $\|\partial_r u\|_{L^2_{r,\gamma}(B_1)}$, notice first that

$$\frac{|u(r)|}{r} \leq \frac{1}{r^2} \int_0^r |f(s)| s \, ds \leq \frac{1}{r^2} \int_0^r |f(s) - f(y)| s \, ds + \frac{1}{2} |f(y)|,$$

where we pick $y \in [0, r]$. Letting $\tilde{B}(y, 2r) = B(y, 2r) \cap B(0, r) \subset \mathbb{R}^2$, where $B(y, 2r)$ is the ball centered at $y$ of radius $2r$, the inequality becomes

$$\frac{|u(r)|}{r} \leq \frac{|\tilde{B}(y, 2r)|}{2\pi r^2} \left( \frac{1}{|\tilde{B}(y, 2r)|} \int_{\tilde{B}(y, 2r)} |f(s) - f(y)| s \, ds \, d\theta \right) + \frac{1}{2} |f(y)|.$$

Here, $|\tilde{B}(y, 2r)|$ denotes the measure of the set $\tilde{B}(y, 2r)$. Since $L^2(B(0, M)) \subset L^1(B(0, M))$ for any ball $B(0, M) \subset \mathbb{R}^2$, with finite radius $M$, we have that $f$ is in $L^1_{\text{loc}}(B(0, M))$. By the Lebesgue differentiation theorem [6, Theorem 3.21], the expression in parenthesis approaches zero as $r \to 0$, while the fraction in
front remains bounded since $|\tilde{B}(y,2r)| > 2\pi r^2$. Therefore, close to the origin, the function $|u(r)|/r$ is bounded by $f(r)$, and using again the equation $\partial_r u = f - \lambda u - u/r$, we find that

$$\|\partial_r u\|_{L^2_{r,\gamma}(B_1)} \leq (2 + (1 + \lambda)C_2(\gamma)) \|f\|_{L^2_{r,\gamma}(\mathbb{R}^2)}.$$ 

It then follows that

$$\|\partial_r u\|_{L^2_{r,\gamma}(\mathbb{R}^2)} \leq (3 + (1 + \lambda)(C_1(\gamma, \lambda) + C_2(\gamma))) \|f\|_{L^2_{r,\gamma}(\mathbb{R}^2)}.$$ 

The above calculations then show that the map $L_\lambda = \partial_r + \frac{1}{r} + \lambda : H^1_{r,\gamma}(\mathbb{R}^2) \rightarrow L^2_{r,\gamma}(\mathbb{R}^2)$ is invertible. Moreover, we also obtain that the operator norm of its inverse satisfies,

$$\|L_\lambda^{-1}\| \leq 3 + (2 + \lambda)(C_1(\gamma, \lambda) + C_2(\gamma)).$$

To extend the result to the more general operator $\partial_r + \frac{1}{r} + \lambda : H^k_{r,\gamma}(\mathbb{R}^2) \rightarrow H^{k-1}_{r,\gamma}(\mathbb{R}^2)$, one can proceed by induction: Assuming that $f$ and $u$ are in $H^k_{r,\gamma}(\mathbb{R}^2)$, one shows that $\partial^k u$ is in $L^2_{r,\gamma}(\mathbb{R}^2)$ using the relation $\partial^k u = \partial^{k-1}(f - u - u/r)$. The fact that $\partial^{k-1}(u/r)$ is in the correct space follows by a similar argument as the one done above to prove $u/r \in L^2_{r,\gamma}(\mathbb{R}^2)$.

To simplify notation, we define $L(\lambda) = \partial_r + \frac{1}{r} + \lambda$ and prove in the next lemma that its inverse, defined over appropriate weighted spaces, is continuously differentiable with respect to the parameter $\lambda$.

**Lemma 2.4.** Let $\lambda > 0$ and $k \in \mathbb{N} \cup \{0\}$ and consider the operator defined by $L(\lambda)u = \partial_r u + \frac{1}{r} u + \lambda u$. Then, its inverse,

$$L^{-1}(\lambda) : H^k_{r,\gamma}(\mathbb{R}^2) \rightarrow H^{k+1}_{r,\gamma}(\mathbb{R}^2)$$

is $C^1$ with respect to the parameter $\lambda$.

**Proof.** Lemma 2.3 shows that $L^{-1}(\lambda)$, with the specified domain and range, is a bounded operator for all $\lambda \in (0, \infty)$. To prove the continuity of this operator with respect to $\lambda$, we must show that given $f \in H^k_{r,\gamma}(\mathbb{R}^2)$,

$$\sup_{\|f\|_{H^k_{r,\gamma}}=1} \| (L^{-1}(\lambda + h) - L^{-1}(\lambda)) f \|_{H^{k+1}_{r,\gamma}} \leq Ch.$$ 

Using the notation $\phi(\lambda) = L^{-1}(\lambda)f$, we notice that

$$(L^{-1}(\lambda + h) - L^{-1}(\lambda)) f = \phi(\lambda + h) - \phi(\lambda) = -hL^{-1}(\lambda) \phi(\lambda + h) = -hL^{-1}(\lambda) L^{-1}(\lambda + h)f,$$

from which the desired result follows. This last expression also shows that for $\lambda > 0$, the derivative of $L^{-1}(\lambda)$ with respect to this parameter is given by

$$\partial_\lambda L^{-1}(\lambda) : H^k_{r,\gamma}(\mathbb{R}^2) \rightarrow H^{k+1}_{r,\gamma}(\mathbb{R}^2)$$

Since the derivative $\partial_\lambda L^{-1}(\lambda)$ is the composition of two continuous operators, it follows that it is itself continuous with respect to $\lambda$.

Finally, the next proposition establishes Fredholm properties for the radial operators $\Delta_n : M^{2,2}_{r,\gamma-2}(\mathbb{R}^2) \rightarrow L^2_{r,\gamma}(\mathbb{R}^2)$,

$$\Delta_n = \partial_{rr} + \frac{1}{r} \partial_r - \frac{n^2}{r^2}, \quad n \in \mathbb{N} \cup \{0\}.$$
The results follows from [19], where it is shown that the Laplace operator \( \Delta : M^{2,p}_{\gamma-2}(\mathbb{R}^2) \to L^p_{\gamma}(\mathbb{R}^2) \) is Fredholm, and the fact that when \( p = 2 \), one can decompose the space \( M^{2,2}_{\gamma-2}(\mathbb{R}^2) \) into a direct sum \( \oplus m^2_{n,\gamma-2} \) where
\[
 m^2_{n,\gamma-2} = \{ u \in M^{2,2}_{\gamma-2}(\mathbb{R}^2) : u(r, \theta) = v_n(r)e^{in\theta}, \ v_n \in M^{2,2}_{r,\gamma-2}(\mathbb{R}^2) \}, \quad n \in \mathbb{Z}.
\]
Notice that because the functions in \( M^{2,2}_{\gamma-2}(\mathbb{R}^2) \) are real valued, we have that \( v_{-n} = \bar{v}_n \). For a detailed proof of this next result, see [13, Lemma 3.1].

**Proposition 2.5.** Let \( \gamma \in \mathbb{R}\setminus\mathbb{Z} \), and \( n \in \mathbb{Z} \). Then, the operator \( \Delta_n : M^{2,2}_{r,\gamma-2}(\mathbb{R}^2) \to L^2_{r,\gamma}(\mathbb{R}^2) \) given by
\[
 \Delta_n \phi = \partial_{rr} \phi + \frac{1}{r} \partial_r \phi - \frac{n^2}{r^2} \phi.
\]
is a Fredholm operator and,
1. for \( 1 - |n| \leq \gamma \leq |n| + 1 \), the map is invertible;
2. for \( \gamma > |n| + 1 \), the map is injective with cokernel spanned by \( r^{|n|} \);
3. for \( \gamma < 1 - |n| \), the map is surjective with kernel spanned by \( r^{|n|} \).

On the other hand, the operator is not Fredholm for integer values of \( \gamma \).

### 3. Intermediate approximations to the viscous Eikonal equation

As mentioned in the introduction, our interest in the viscous eikonal equation
\[
 \tilde{\phi}_t = \Delta \tilde{\phi} - b|\nabla \tilde{\phi}|^2 - \varepsilon g(x), \quad x \in \mathbb{R}^2,
\]
stems from its role as a model equation for the phase dynamics of target patterns and spiral waves in oscillatory systems. We are therefore interested in solutions of the form \( \tilde{\phi}(x, t) = \phi(x) - \Omega t \), which then satisfy the steady-state equation,
\[
 \Delta \phi - b|\nabla \phi|^2 - \varepsilon g(x) + \Omega = 0 \quad x \in \mathbb{R}^2. \tag{4}
\]
Because the gradient, \( \nabla \phi \), then approximates the pattern’s wavenumber, target patterns then correspond to those \( \phi \) which in addition fulfill the boundary conditions, \( \nabla \phi \to k \) as \( |x| \to \infty \). Consequently, we look for solutions to equation (4) that bifurcate from zero when \( \varepsilon > 0 \), and whose gradients are bounded at infinity.

Notice that the condition on the gradient, \( \nabla \phi \), provides enough information to derive an equation that is valid at intermediate scales. Indeed, assuming a regular perturbation for both \( \phi \) and \( \Omega \), one obtains at order \( O(\varepsilon) \) the equation,
\[
 \partial_{rr} \phi_1 + \frac{1}{r} \partial_r \phi_1 - g = -\Omega_1.
\]
A short calculation then shows that in order to obtain solutions with bounded derivatives, the parameter \( \Omega_1 \) must be zero. Continuing this perturbation analysis, one checks that this condition must be satisfied at all orders of \( \varepsilon \). In other words, the frequency, \( \Omega \), must be small beyond all orders of this parameter. Consequently, at intermediate scales the system is well approximated by the intermediate equation
\[
 \Delta_0 \phi - b(\partial_r \phi)^2 - \varepsilon g_c - \varepsilon g_f = 0. \tag{5}
\]
Notice that we have explicitly written the inhomogeneity as the sum of two functions satisfying Hypothesis 1.1. This choice of notation will be used next in Sect. 3.1, where we construct a first-order approximation for the above equation. We then use this information to prove existence of solutions to Eq. (5) in Sect. 3.2.
Remark 3.2. Because $\gamma_c$ to denote the Euler–Mascheroni constant and the symbols $\chi, \chi_M \in C^\infty(\mathbb{R}^2)$ to denote smooth radial cut-off functions satisfying
\[
\chi(x) = \begin{cases} 
0 & \text{if } |x| < 1 \\
1 & \text{if } |x| > 2,
\end{cases} \quad \chi_M(x) = \begin{cases} 
0 & \text{if } |x| < 1 \\
1 & 2 < |x| < M, \\
0 & 2M < |x|.
\end{cases}
\]
where $M$ is a positive constant. Notice, in particular, that $\chi_M$ has compact support.

3.1. First-order approximation

We construct a first-order approximation, $\phi$, which is the sum of two functions, $\phi_0$ and $\phi_1$. We take
\[
\phi_0 = -\frac{1}{b} \chi_M \log(1 + a \log r + \varepsilon K(r)),
\]
a choice that is motivated by the Hopf–Cole transform, $\phi = -\frac{1}{b} \log \Psi$, which turns the eikonal equation (5) into the steady-state Schrödinger equation with potential $\varepsilon g$. The value of the constant $M$ appearing in the definition of the cut-off function $\chi_M$ is taken so that the expression $\log(1 + a \log r + \varepsilon K(r))$, always remains bounded. The constant $a$ is a parameter that is determined when constructing the second part to the approximation, while the function $K(r)$ satisfies
\[
\Delta_0 K + bg_f = 0.
\]

The fact that we can solve this last equation follows from our assumptions on the inhomogeneity. Recall that $g_f \in H^k_\sigma(\mathbb{R}^2)$ for $0 < \sigma < 1$. Proposition 2.5 then shows that the radial Laplacian, $\Delta_0$, is a surjective operator with a one-dimensional kernel spanned by $\{1\}$. We can therefore use Lyapunov–Schmidt reduction to solve this equation and find a family of solutions
\[
K(r) = K_\sigma(r) + c, \quad c \in \mathbb{R}.
\]
Since $c$ is arbitrary, without loss of generality we pick $c = 0$. We also find that the solution, $K$, belongs to the space $M^{2,2}_{r,\sigma-2}(\mathbb{R}^2)$. In fact, one can check that $K$ has more regularity and is in the space $R^k_\sigma = \{u \in M^{2,2}_{r,\sigma-2}(\mathbb{R}^2) : D^2 u \in H^k_\sigma(\mathbb{R}^2)\}$.

Remark 3.1. Notice that the function $\phi_0$ is as regular as the function $K$ and that as a result the derivative $\partial_r \phi_0$ is in the space $H^{k+1}_\sigma(\mathbb{R}^2)$. In addition, this function is bounded and has compact support.

Remark 3.2. Because $g_f$ is in $H^k_\sigma(\mathbb{R}^2)$ with $\sigma \in (0,1)$ and $k \geq 2$, we then have the following decay properties for the solution to Eq. (7):

- if $g_f$ decays like $1/r^m$ in the far field, with $1 < m < 2$, then $K \sim O(r^{2-m})$ at infinity, while
- if $g_f$ decays like $1/r^2$ in the far field, then $K \sim O((\log r)^2)$ at infinity.

Next, we define the second function, $\phi_1$, as the solution to the equation
\[
\Delta_0 \phi_1 - \frac{a}{b} \Delta_0 (\chi \log r) - \varepsilon g_c = 0.
\]
Here, the constant $a$ is the same as the one appearing in the definition of $\phi_0$, and the function $g_c$ is in the space $H^k_\gamma(\mathbb{R}^2)$ with $1 < \gamma$, by assumption.

To justify the existence of $\phi_1$, we use again Proposition 2.5 which shows that for values of $\gamma > 1$, the operator $\Delta_0 : M^{2,2}_{r,\gamma-2}(\mathbb{R}^2) \to L^2_{r,\gamma}(\mathbb{R}^2)$ is Fredholm with index $-1$, and cokernel spanned by $\{1\}$. Because the projection of $\Delta_0 (\chi \log r)$ onto the cokernel is non-trivial, i.e.,
\[
\int_{\mathbb{R}^2} \Delta_0 (\chi \log r) \, dx = 2\pi,
\]
the bordering lemma stated at the end of this subsection then shows that the operator
\[ M_{r,γ}^{2,2}(\mathbb{R}^2) \times \mathbb{R} \to L_{r,γ}^2(\mathbb{R}^2) \]
\[
(φ, a) \mapsto Δ_0φ_1 - \frac{a}{b}Δ_0(χ log r)
\]
is invertible. Therefore, the equation for φ₁ is indeed solvable. In addition, projecting onto the constant functions we also find that
\[
a = -εb \int_0^∞ g_c(r) \, r \, dr.
\]  
Finally, since \( g_c \in H^k(\mathbb{R}^2) \), it follows that our solution φ₁ is in the space \( R^k γ \), defined as in (8).

**Lemma 3.3.** [Bordering Lemma] Let \( X \) and \( Y \) be Banach spaces, and consider the operator
\[
S = \begin{bmatrix} A & B \\ C & D \end{bmatrix} : X × \mathbb{R}^p \to Y × \mathbb{R}^q,
\]
with bounded linear operators \( A : X \to Y \), \( B : \mathbb{R}^p \to Y \), \( C : X \to \mathbb{R}^q \), \( D : \mathbb{R}^p \to \mathbb{R}^q \). If \( A \) is Fredholm of index \( i \), then \( S \) is Fredholm of index \( i + p - q \).

**Proof.** One can write \( S \) as the sum of a block diagonal operator with the indicated index, \( i + p - q \), and a compact operator consisting of the off-diagonal elements. Since compact perturbations do not alter the index of a Fredholm operator, the result then follows. □

### 3.2. Existence of solutions to intermediate approximation

Using the first-order approximation, \( φ₀ + φ₁ \), defined in the previous subsection we now prove the existence of solutions to Eq. (5) using the implicit function theorem.

Inserting the ansatz
\[
φ = φ₀ + φ₁ + φ₂
\]
into Eq. (5), one obtains the following expression for \( φ₂ \),
\[
Δ_0φ₂ - b \left( 2∂ᵣ φ₀(∂ᵣ φ₁ + ∂ᵣ φ₂) + (∂ᵣ φ₁ + ∂ᵣ φ₂)^2 \right) + G₁ = 0,
\]  
where the term \( G₁ \) is given by
\[
G₁ = Δ_0φ₀ - b(∂ᵣ φ₀)^2 - εg_f + \frac{a}{b}Δ_0(χ log r)
\]

To continue the analysis, we let \( ψ = ∂ᵣ φ₂ \) in Eq. (10) and add and subtract the term \( λψ \). We assume that the parameter \( λ \) is sufficiently small, positive, and fixed. The result is,
\[
∂ᵣ ψ + \frac{1}{r}ψ + λψ - b \left( 2∂ᵣ φ₀(∂ᵣ φ₁ + ψ) + (∂ᵣ φ₁ + ψ)^2 \right) + G₁ - λψ = 0.
\]

Letting \( L_λ = ∂ᵣ + \frac{1}{r} + λ \), we may precondition this last equation with \( L_λ^{-1} \) and write
\[
ψ + L_λ^{-1} \left[ -b \left( 2∂ᵣ φ₀(∂ᵣ φ₁ + ψ) + (∂ᵣ φ₁ + ψ)^2 \right) + G₁ - λψ \right] = 0.
\]

Because the above expression is equivalent to the intermediate equation (5), if we find a solution \( ψ \) to (11), we immediately obtain a corresponding solution to (5) of the form
\[
φ(r; ε) = φ₀(r; ε) + φ₁(r; ε) + φ₂(r; ε) + εc,
\]
where \( ∂ᵣ φ₂ = ψ \) and \( c \) is a constant of integration.

To use the implicit function theorem, we view the left-hand side of (11) as an operator \( F : H^k_γ(\mathbb{R}^2) × \mathbb{R}_+ \to H^k_γ(\mathbb{R}^2) \) for some appropriate \( δ ∈ \mathbb{R} \) and show that it is well defined, smooth with respect to \( ε \),
and that its Fréchet derivative $D_{\psi}F(0; 0) : H^k_\delta(\mathbb{R}^2) \to H^k_\delta(\mathbb{R}^2)$ is invertible. The result is the following theorem.

**Theorem 2.** Let $k \geq 2$, $\sigma \in (0, 1)$ and consider functions $g \in H^k_{r, \sigma}(\mathbb{R}^2)$ satisfying Hypothesis 1.1. Then, the intermediate equation

$$\Delta_0 \phi - b(\partial_r \phi)^2 - \varepsilon g = 0,$$

has a family of solutions

$$\phi(r; \varepsilon) = -\frac{1}{b} \chiM \log \left(1 + a \log r + \varepsilon K\right) + \phi_1(r; \varepsilon) + \phi_2(r; \varepsilon) + \varepsilon c, \quad c \in \mathbb{R},$$

that bifurcates from zero at $\varepsilon = 0$ and is $C^1$ in $\varepsilon \in [0, \infty)$. Moreover, letting $\gamma \in (1, \infty)$, and $\sigma \in (0, 1)$ be defined as in Hypothesis 1.1, the family of solutions satisfies:

- $\phi_1 \in \{u \in M^{2, \infty}_{r, \gamma - 2}(\mathbb{R}^2) : D^2 u \in H^k_\delta(\mathbb{R}^2)\}$
- $\partial_r \phi_2 \in H^k_\delta(\mathbb{R}^2)$, where $\delta = \min(\gamma - 1, \sigma)$.
- $\Delta_0 K = g_f$, with $K \in R^k_\gamma$, and
- $a = -\varepsilon b \int_0^\infty g_c(r) \, r \, dr$.

**Proof.** As already mentioned, the result follows from finding solutions to Eq. (11) using the implicit function theorem. We therefore consider the left-hand side of this equation as an operator $F : H^k_\delta(\mathbb{R}^2) \times \mathbb{R}_+ \to H^k_\delta(\mathbb{R}^2)$, with $\delta = \min(\gamma - 1, \sigma) > 0$.

Since the operator’s dependence on $\varepsilon$ comes from the three functions $\partial_r \phi_0, \partial_r \phi_1$, and $G_1$ and these functions are all smooth with respect to $\varepsilon$ on the interval $[0, \infty)$, then the same result holds for the operator $F$.

To show that the Fréchet derivative $D_{\psi}F(0; 0) = Id - \lambda \mathcal{L}^{-1}_\lambda$ is invertible, we recall the results from Sect. 2. In particular, Lemma 2.3 shows that $\mathcal{L}^{-1}_\lambda : H^k_\delta(\mathbb{R}^2) \to H^{k+1}_\delta(\mathbb{R}^2)$ is bounded. Since the embedding $H^{k+1}_\delta(\mathbb{R}^2) \subset H^k_\delta(\mathbb{R}^2)$ is continuous, it then follows that $D_{\psi}F(0; 0) : H^k_\delta(\mathbb{R}^2) \to H^k_\delta(\mathbb{R}^2)$ is a small perturbation of the identity operator, and is thus invertible for a sufficiently small $\lambda$.

To complete the proof, we need to show that the operator $F$ is well defined. Taking into account the results of Lemma 2.3, this is equivalent to showing that the expression

$$N(\psi, \varepsilon) = -b \left(2 \partial_r \phi_0(\partial_r \phi_1 + \psi) + (\partial_r \phi_1 + \psi)^2\right) + G_1 - \lambda \psi,$$

defines a bounded operator $N : H^k_\delta(\mathbb{R}^2) \times \mathbb{R}_+ \to H^{k-1}_\delta(\mathbb{R}^2)$.

We start with the term $\partial_r \phi_0(\partial_r \phi_1 + \psi)$. From the definition of $\phi_1$, we know that this is a function in $R^k_\gamma$ with $\gamma > 1$. In particular,

$$\partial_r \phi_1 \in \{u \in M^{2, \infty}_{r, \gamma - 2}(\mathbb{R}^2) : D u \in H^k_\delta(\mathbb{R}^2)\} \subset H^{k+1}_{\gamma - 1}(\mathbb{R}^2).$$

Because $\psi \in H^k_\delta(\mathbb{R}^2)$ and $\delta = \min(\gamma - 1, \sigma) > 0$, it then follows that the sum $(\partial_r \phi_1 + \psi)$ is also in this space. Since $\partial_r \phi_0$ has compact support and $k+1$ bounded derivatives (see Remark 3.1), then the product $\partial_r \phi_0(\partial_r \phi_1 + \psi)$ is also well defined in $H^k_\delta(\mathbb{R}^2)$.

Next, since $(\partial_r \phi_1 + \psi)$ is in $H^k_\delta(\mathbb{R}^2)$, with $\delta > 0$ and $k \geq 2$, Lemma 3.4 shows that $(\partial_r \phi_1 + \psi)^2$ is in $H^{k-1}_\delta(\mathbb{R}^2)$. Finally, Lemma 3.5 at the end of this section shows that $G_1 \in H^k_\delta(\mathbb{R}^2)$, and because $\delta = \min(\gamma - 1, \sigma) > 0$, this term is also well defined.

Since the operator $F$ satisfies the assumptions of the implicit function theorem, we obtain a family of solutions $\psi(r; \varepsilon)$ that bifurcates from zero and is smooth with respect to $\varepsilon$. Because $\psi = \partial_r \phi_2$, we arrive at the family

$$\phi(r; \varepsilon) = \phi_0(r; \varepsilon) + \phi_1(r; \varepsilon) + \phi_2(r; \varepsilon) + \varepsilon c, \quad c \in \mathbb{R}.$$

This finishes the proof of the theorem.
Lemma 3.4. Let $\psi \in H^k_\gamma(\mathbb{R}^2)$ with $\gamma > 0$ and $k \geq 2$. Then, $\psi^2 \in H^{k-1}_\gamma(\mathbb{R}^2)$.

Proof. To simplify the analysis, we let $D^j$ denote any $j$th-order derivative, and we only prove that $D^{k-1}(\psi^2)$ is in $L^2_{\gamma,\gamma}(\mathbb{R}^2)$, since a similar analysis shows that lower derivatives are in this same space. Because $k \geq 2$ and $\gamma > 0$, it follows by Sobolev embeddings that $D^j \psi \in H^2(\mathbb{R}^2) \subseteq C_B(\mathbb{R}^2)$ for $0 \leq j \leq k-2$. Then, writing

$$D^{k-1}(\psi^2) = \sum_{j=0}^{k-1} \binom{k-1}{j} D^{k-1-j} \psi \partial^j \psi$$

we see that this derivative can be written as a product of a bounded function and a function that is in $L^2_{\gamma,\gamma}(\mathbb{R}^2)$. Hence, $\psi^2 \in H^{k-1}_\gamma(\mathbb{R}^2)$. \qed

The next lemma shows that $G_1$ is in $H^k_\gamma(\mathbb{R}^2)$ with $\sigma \in (0,1)$.

Lemma 3.5. Let $k \geq 2$, $\sigma \in (0,1)$ and take $g_f \in H^k_\gamma(\mathbb{R}^2)$. Consider the function $\phi_0$ constructed from $g_f$ and described above in (6). Then, the expression

$$G_1 = \Delta_0 \phi_0 - b(\partial_r \phi_0)^2 - \varepsilon g_f + \frac{a}{b} \Delta_0 (\chi \log r)$$

is also in $H^k_\gamma(\mathbb{R}^2)$.

Proof. Using the notation $\phi_0 = \chi_M \tilde{\phi}_0$, we first expand $G_1$

$$G_1 = \Delta_0 \phi_0 - b(\partial_r \phi_0)^2 - \varepsilon g_f + \frac{a}{b} \Delta_0 (\chi \log r)$$

$$G_1 = \left( \tilde{\phi}_0 \Delta_0 \chi_M + 2 \chi_M \tilde{\phi}_0 \right) - b(2 \chi_M \tilde{\phi}_0 \partial_r \tilde{\phi}_0 + b(\partial_r \tilde{\phi}_0)^2) (\chi_M - \chi_M^2) \right)$$

$$- \frac{1}{b} \chi_M \left[ \frac{\Delta_0(a \log r + \varepsilon K)}{1 + a \log r + \varepsilon K} \right] - \varepsilon g_f + \frac{a}{b} \Delta_0 (\chi \log r)$$

$$G_1 = \left( \tilde{\phi}_0 \Delta_0 \chi_M + 2 \chi_M \tilde{\phi}_0 \right) - b(2 \chi_M \tilde{\phi}_0 \partial_r \tilde{\phi}_0 + b(\partial_r \tilde{\phi}_0)^2) (\chi_M - \chi_M^2) \right)$$

$$- \frac{1}{b} \chi_M \left[ \frac{\Delta_0 a \log r}{1 + a \log r + \varepsilon K} \right] - \frac{1}{b} \chi_M \left[ \frac{\varepsilon \Delta_0 K}{1 + a \log r + \varepsilon K} \right] - \varepsilon g_f + \frac{a}{b} \Delta_0 (\chi \log r)$$

Because the log $r$ is a fundamental solution of the Laplacian and since $\chi_M$ is zero near the origin, then the term

$$\chi_M \left[ \frac{\Delta_0 a \log r}{1 + a \log r + \varepsilon K} \right] = 0.$$

Similarly, because the function $K$ is a solution to $\Delta_0 K + bg_f = 0$, then we may write

$$- \frac{1}{b} \chi_M \left[ \frac{\varepsilon \Delta_0 K}{1 + a \log r + \varepsilon K} \right] - \varepsilon g_f = - \varepsilon g_f (1 - \chi_M) - \varepsilon g_f \chi_M \left[ \frac{(a \log r + \varepsilon K)}{1 + a \log r + \varepsilon K} \right].$$

Therefore,

$$G_1 = \left( \tilde{\phi}_0 \Delta_0 \chi_M + 2 \chi_M \tilde{\phi}_0 \right) - b(2 \chi_M \tilde{\phi}_0 \partial_r \tilde{\phi}_0 + (\partial_r \tilde{\phi}_0)^2) (\chi_M - \chi_M^2) \right)$$

$$- \varepsilon g_f (1 - \chi_M) - \varepsilon g_f \chi_M \left[ \frac{(a \log r + \varepsilon K)}{1 + a \log r + \varepsilon K} \right] + \frac{a}{b} \Delta_0 (\chi \log r)$$

From the definition of $\chi_M$, it is clear that all terms involving a derivative of this function are localized and have compact support. Because the value of $M$ in the definition of $\chi_M$ was chosen to vanish whenever the expression $1 + a \log r + \varepsilon K$ is $\leq 0$, we see that the term in brackets is also bounded and with compact
support. In addition, this term is as regular as the function $K \in R^k$ and is therefore in $H^{k+2}(\mathbb{R}^2)$. On the other hand, the function $(1−\chi_M)\varepsilon_f$ behaves like $g_f$ at infinity, and as a result, it is in the same space as the inhomogeneity. Finally, because $\Delta_0 \log r = 0$ on $\mathbb{R}^2 \setminus \{0\}$, the function $\Delta_0(\chi \log r)$ is localized and smooth. Taking all this into account, we may conclude that $G_1$ is in the space $H^k_{\sigma}(\mathbb{R}^2)$. \hfill \square

4. Far-field approximation to the viscous Eikonal equation

In this section, we consider again the full equation

$$\Delta_0 \phi - b(\partial_r \phi)^2 - \varepsilon g(r) + \Omega = 0 \quad r = |x| \in [0, \infty),$$

but assume that the value of $\Omega$ is fixed and different from zero. As in Sect. 3, we first find an appropriate expression for the far-field behavior of the solution and a first-order approximation for this new equation. We then use this result to prove existence of solutions using the implicit function theorem.

Because the inhomogeneity is algebraically decaying, for large values of $r$ the relevant terms in the equation are

$$\Delta_0 \phi - b(\partial_r \phi)^2 + \Omega = 0$$

To find a first-order approximation, we can again use the Hopf–Cole transform, $\phi(r) = -(1/b) \log(K)$, rewriting the equation as

$$\partial_r K + \frac{1}{r} \partial_r K - \Lambda^2 K = 0 \quad \Lambda^2 = b\Omega.$$

Notice that this is either a Bessel, or the modified Bessel equation, depending on the sign of $b\Omega$. Because we are interested in solutions, $\phi(r)$, that are real, we pick $b\Omega > 0$ so that the solution to this last equation is $K = K_0$, the modified Bessel function of the second kind. In particular, because $K_0(z) \sim O(e^{-z})$ as $z \to \infty$ (see Table 1), this implies that $\partial_r \phi(r)$ is bounded in the far field, as desired.

We therefore consider the ansatz

$$\phi(r) = \phi_0(r) + \phi_1(r)$$

where $\phi_0$ is given by

$$\phi_0(r) = -\frac{1}{b} \chi(\Lambda r) \log(K_0(\Lambda r)), \quad \Lambda^2 = b\Omega > 0.$$

Here, again $\chi$ represents a cut-off function that removes the singular behavior of the log function near the origin. Inserting this expression into Eq. (12) gives

$$\Delta_0 \phi_1 - 2b \partial_r \phi_0 \partial_r \phi_1 - b(\partial_r \phi_1)^2 + (\Delta_0 \phi_0 - b(\partial_r \phi_0)^2) - \varepsilon g = 0.$$

Since the terms appearing in the parenthesis represent a localized function with compact support, they do not contribute to the behavior of the solution for large values of $r$. Thus, the far-field behavior of the solution is determined by

$$\Delta_0 \phi_1 - 2b \partial_r \phi_0 \partial_r \phi_1 - b(\partial_r \phi_1)^2 - \varepsilon g = 0. \quad (13)$$

Letting $\psi = \partial_r \phi_1$ and adding and subtracting the term $2\Lambda \psi$ give us

$$\partial_r \psi + \frac{1}{r} \psi + 2\Lambda \psi + [-2b \partial_r \phi_0 \psi - b\psi^2 - \varepsilon g - 2\Lambda \psi] = 0.$$

We can then precondition this equation by $L_{2\Lambda}^{-1}$, since by Lemma 2.3 we know that this operator is bounded for all values of $\Lambda > 0$, if its domain is $H^k_{r,\sigma}(\mathbb{R}^2)$. Thus, the equation can be written as

$$F(\psi; \varepsilon) = \text{Id} + L_{2\Lambda}^{-1} \left[-2b \partial_r \phi_0 \psi - b\psi^2 - \varepsilon g - 2\Lambda \psi \right] = 0. \quad (14)$$
In what follows, we will show that the operator $F : H^k_{r,\sigma}(\mathbb{R}^2) \times \mathbb{R} \rightarrow H^k_{r,\sigma}(\mathbb{R}^2)$ satisfies the conditions of the implicit function theorem and prove the following theorem.

**Theorem 3.** Take $g \in H^k_{r,\sigma}$ with $k \geq 2$ and let $\sigma \in (0,1)$. Then, there exist a positive constant $\Lambda_0$ such that for any fixed $\Lambda \in (0,\Lambda_0)$, there is an $\varepsilon_0 > 0$, and $C^1$ family of solutions, $\phi = \phi(r; \varepsilon)$, to Eq. (13) that bifurcates from zero at $\varepsilon = 0$ and is valid for $\varepsilon \in (-\varepsilon_0, \varepsilon_0)$. Moreover, this family has the form

$$
\phi(r; \varepsilon) = -\frac{1}{b} \chi(\Lambda r) \log(K_0(\Lambda r)) + \phi_1(r; \varepsilon) + \varepsilon c
$$

where

i) $K_0(z)$ represents the zeroth-order modified Bessel function of the second kind,

ii) $\partial_r \phi_1 \in H^k_{r,\sigma}(\mathbb{R}^2)$,

iii) and $c \in \mathbb{R}$ is an arbitrary constant.

**Proof.** Because finding solutions to Eq. (13) is equivalent to finding the zeros of the operator $F : H^k_{r,\sigma}(\mathbb{R}^2) \times \mathbb{R} \rightarrow H^k_{r,\sigma}(\mathbb{R}^2)$ defined in (14), we check that $F$ satisfies the assumptions of the implicit function theorem.

It is clear that $F(0; 0) = 0$ and that this operator is smooth with respect to the parameter $\varepsilon$. To check that the Fréchet derivative, $D_{\psi}F(0; 0) : H^k_{r,\sigma}(\mathbb{R}^2) \rightarrow H^k_{r,\sigma}(\mathbb{R}^2)$, given by

$$
D_{\psi}F(0; 0) = Id + L^{-1}_{2\Lambda}[-2b\partial_r \phi_0 - 2\Lambda],
$$

is invertible, notice that the term in the brackets can be written as the product of a bounded function times the constant $2\Lambda$. Indeed, this can be checked by expanding this term,

$$
-2b\partial_r \phi_0 - 2\Lambda = -2b\partial_r \left[-\frac{1}{b} \chi(\Lambda r) \log[K_0(\Lambda r)] \right] - 2\Lambda
$$

$$
= 2\Lambda \left[ \chi'(\Lambda r) \log[K_0(\Lambda r)] + \chi(\Lambda r) \frac{K_0'(\Lambda r)}{K_0(\Lambda r)} - 1 \right]
$$

and using the fact that the ratio $\frac{K_0'(z)}{K_0(z)} = -1$ as $r \rightarrow \infty$ and that $\chi'$ has compact support. Since the operator $L^{-1}_{2\Lambda} : H^k_{r,\sigma}(\mathbb{R}^2) \rightarrow H^{k+1}_{r,\sigma}(\mathbb{R}^2)$ is bounded, it follows that there is a small number $\Lambda_0 > 0$ such that if $\Lambda \in (0,\Lambda_0)$, the derivative $D_{\psi}F(0; 0)$ is a small perturbation of the identity and is therefore invertible.

We are left with showing that the operator $F$ is well defined. Taking into account again that the map $L^{-1}_{2\Lambda} : H^{k-1}_{r,\sigma}(\mathbb{R}^2) \rightarrow H^k_{r,\sigma}(\mathbb{R}^2)$ is bounded, this is equivalent to showing that the terms

$$
-2b\partial_r \phi_0 \psi - b\psi^2 - \varepsilon g - 2\Lambda \psi
$$

define a bounded operator $N : H^k_{r,\sigma}(\mathbb{R}^2) \times \mathbb{R} \rightarrow H^{k-1}_{r,\sigma}(\mathbb{R}^2)$.

First, notice that by assumption, the impurity $g$ is in the desired space. As for the elements involving the variable $\psi \in H^k_{r,\sigma}(\mathbb{R}^2)$, because the derivative $\partial_r \phi_0$ is a bounded function, we can easily check that they are both in the space $H^{k-1}_{r,\sigma}(\mathbb{R}^2)$. Finally, since $\sigma > 0$, Lemma 3.4 shows that the product $\psi^2$ is in $H^{k-1}_{r,\sigma}(\mathbb{R}^2)$.

This proves that the operator $F$ satisfies the conditions of the implicit function theorem and proves the existence of a family of solutions solving $F(\psi; \varepsilon) = 0$. Going back to the definition of $\psi = \partial_r \phi_1$, we see that the above result also gives us a family of solutions $\phi(r; \varepsilon) = \phi_0(r) + \phi_1(r; \varepsilon) + \varepsilon c$ solving the far-field equation (12), where $\phi_1 \in H^{k+1}_{r,\sigma}(\mathbb{R}^2)$ and $c$ is for now an arbitrary constant which is the result of integrating $\psi$. This proves the result of the theorem. \qed
5. Existence of target patterns

5.1. Matching

To determine an expression for the eigenvalue $\Omega$, we must match the intermediate and far-field approximations of the wavenumber, $\partial_r \phi$. For convenience, we recall their expressions,

$$\phi_{\text{far}}(r; \varepsilon, \Lambda) = -\frac{1}{b} \chi(\Lambda r) \log(K_0(\Lambda r)) + \phi_1(r; \varepsilon) + \varepsilon c$$

$$\phi_{\text{int}}(r; \varepsilon) = -\frac{1}{b} \chi_M \log \left(1 + a \log r + \varepsilon K(r)\right) + \phi_1(r; \varepsilon) + \phi_2(r; \varepsilon) + \varepsilon c,$$

As before, $K_0$ denotes the modified Bessel function of the first kind, while the function $K$ satisfies

$$\Delta_0 K + bg_f = 0.$$

Notice that the remaining terms, $\phi_1$, $\phi_2$, all have derivatives that decay algebraically at infinity. In particular,

1. The function $\bar{\phi}_1$ defined in Sect. 3.1 is in the space $\tilde{M}_2 \subset M^{2, 2}_{r, r^{-\gamma - 2}}$, with $\gamma > 1$. From Lemma 2.1, it follows that $|\bar{\phi}_1| < |x|^{-\gamma + 1}$. In particular, if the inhomogeneity $g_c \sim O(r^{-(d+2)})$, with $d > 0$, we have that $\partial_r \bar{\phi}_1 \sim O(r^{-(d+1)})$.

2. From Theorems 2 and 3, we know that the functions $\partial_r \phi_1$ and $\partial_r \phi_2$ are in the space $H^2_b(\mathbb{R}^2)$, where $\delta \in (0, 1)$. It then follows from Sobolev embeddings that these functions are bounded. In addition, because $\delta > 0$, they must decay algebraically.

To do the matching, recall from the analysis in Sect. 3 that the parameter $\Lambda^2 = b\Omega$ is assumed to be small beyond all orders of $\varepsilon$. This justifies the scaling $r = \eta(\varepsilon)r_n$, where $r_n$ is a constant and $\eta(\varepsilon) = \varepsilon/\Lambda$. As a result, $\Lambda r \to 0$ as $\varepsilon \to 0$, while $r \to \infty$, and we find that for small value of $\eta$ we are in the region where both approximations are valid. Moreover, since $\varepsilon \sim o(\eta(\varepsilon))$ there is always an open interval where the two approximations can be matched, even as $\varepsilon \to 0$. Because in this region the functions $\chi_M = \chi = 1$, we obtain

$$\partial_r \phi_{\text{far}}(r; \varepsilon, \Lambda) \sim -\frac{1}{b} \left[ \frac{\Lambda K_0'(\Lambda r)}{K_0(\Lambda r)} + \partial_r \phi_1 \right]$$

$$\partial_r \phi_{\text{int}}(r; \varepsilon) \sim -\frac{1}{b} \left[ \frac{a/r + \varepsilon \partial_r K(r)}{1 + a \log r + \varepsilon K(r)} \right] + \partial_r \phi_1 + \partial_r \phi_2.$$

Setting the derivatives equal to each other, $\partial_r \phi_{\text{far}} = \partial_r \phi_{\text{int}}$, we find that

$$\left( \Lambda K_0'(\Lambda r) - bK_0(\Lambda r) \partial_r \phi_1 \right) (1 + a \log r + \varepsilon K(r))$$

$$= K_0(\Lambda r) \left[ (a/r + \varepsilon \partial_r K(r)) - b(\partial_r \phi_1 + \partial_r \phi_2)(1 + a \log r + \varepsilon K(r)) \right]$$

$$\left( - \Lambda \left( \frac{1}{\Lambda r} + O(\Lambda r) \right) - bK_0(\Lambda r) \partial_r \phi_1 \right) (1 + a \log r + \varepsilon K(r))$$

$$= \left( - \log(\Lambda/2) - \gamma_e - \log(r) + O((\Lambda r)^2) \right)$$

$$\times \left[ (a/r + \varepsilon \partial_r K(r)) - b(\partial_r \phi_1 + \partial_r \phi_2)(1 + a \log r + \varepsilon K(r)) \right],$$

where in the second line we use the fact that $K_0'(z) = -K_1(z)$ and the expansions from Table 1.

We now proceed with the matched asymptotic analysis to determine the value of $\Lambda$. Notice that due to the relation $\Omega b = \Lambda^2$, this will also allow us to obtain an expression for the frequency. The method is as follows: We first divide the above expression by different gage functions in order to select terms of similar order in $\varepsilon$. We then cancel any duplicate terms, let $\varepsilon$ go to zero, and select the value of any undefined constant so that the remaining terms add up to zero.
Because we are interested only in finding the value of the constant $\Lambda$, we can simplify these computations by noticing that terms of the form $O(\Lambda^2)$ will go to zero, as $\varepsilon \to 0$, faster than any other term. Thus, they are not of the same order in $\varepsilon$ as elements that involve $\Lambda$. This follows from the scalings picked and the algebraic decay rate of the the function $\partial_r \phi_1$. We may therefore consider instead the expression

$$-\frac{1}{r} (1 + a \log r + \varepsilon K(r)) = (- \log(\Lambda/2) - \gamma_e - \log(r))$$

$$\times \left(\frac{a}{r} + \varepsilon \partial_r K(r) - b(\partial_r \phi_1 + \partial_r \phi_2)(1 + a \log r + \varepsilon K(r))\right).$$

It is worth pointing out here that, in contrast to more standard matched asymptotic analyses, the elements in Eq. (15) are not of order $O(\varepsilon^n)$, $n \in \mathbb{N}$. Moreover, we find that dominants terms depend on the yet to be determined approximations $\phi_i$, $i = 1, 2$. Thus, we will not be able to match them exactly, but we can justify that the process can be done.

First, looking at the right-hand side, one notices that the dominant term is $-\varepsilon \log(r) \partial_r K$. Because $r$ depends on $\varepsilon$, we may use this function as a gage function. Dividing by $-\varepsilon \log(r) \partial_r K$ and letting $\varepsilon \to 0$, or equivalently $r \to \infty$, we are left with matching,

$$0 = 1 - b \frac{\partial_r \phi_1 + \partial_r \phi_2}{\varepsilon \partial_r K}(1 + a \log r + \varepsilon K).$$

By picking the value of $\delta \in (0, 1)$ so that the higher-order correction term, $\phi_2 \in L^2_\delta(\mathbb{R}^2)$, is in the same space as both, $K$ and $\partial_r K$, we see that it is possible to match these terms. Expression (15) then becomes

$$-\frac{1}{r} (1 + a \log r + \varepsilon K(r)) = (- \log(\Lambda/2) - \gamma_e)$$

$$\times \left(\frac{a}{r} + \varepsilon \partial_r K(r) - b(\partial_r \phi_1 + \partial_r \phi_2)(1 + a \log r + \varepsilon K(r))\right)$$

$$- \frac{a}{r} \log r.$$

Second, canceling the term $\frac{a}{r} \log r$, using $\varepsilon K/r$ as a gage function, and letting $r \to \infty$, we obtain

$$-1 = (- \log(\Lambda/2) - \gamma_e) \left[ \frac{\partial_r K}{r K} - \frac{b r}{\varepsilon K}(\partial_r \phi_1 + \partial_r \phi_2)(1 + a \log r + \varepsilon K(r)) \right].$$

Since $\phi_i$, $i = 1, 2$ represent all higher-order terms, and not just one function, one can again justify that these terms can be matched. As a result, Eq. (15) now reads

$$-\frac{1}{r} = (- \log \Lambda/2 - \gamma_e) \frac{a}{r}.$$

Finally, solving for $\Lambda$, we see that $\Lambda = 2e^{-\gamma_e} \exp(1/a)$. Using the relation $\Lambda^2 = b \Omega$, we also obtain that

$$\Omega = \frac{4}{b} e^{-2\gamma_e} \exp\left(\frac{1}{a}\right).$$
In particular, from the definition of $a$, i.e., $a = -be \int_0^{\infty} g_c(r) r \, dr$, we may conclude that both $\Lambda$ and $\Omega$ are smooth functions of $\varepsilon$, for all $\varepsilon \in (0, \varepsilon_M)$, with $\varepsilon_M$ a positive constant. In addition, notice that as $\varepsilon$ approaches zero, the value of $\Lambda$ and $\Omega$ also goes to zero.

**Remark 5.1.** Notice that:

1. We need the constant $a < 0$ in order for $\Lambda = b\Omega$ to satisfy our initial assumption of being small beyond all orders of $\varepsilon$. If $\varepsilon > 0$, this condition is guaranteed from formula (9) and the assumption that $g$ is a positive function.

2. Notice also that if $\varepsilon \int g < 0$, the gradient $\partial_r \phi_{\text{int}}$ would also be negative and we would not be able to match the two approximations. This is in line with previous results which show that target pattern solutions (or thanks to the Hopf–Cole transform, $\phi = -\frac{1}{b} \log(\Psi)$, ground states of the Schrödinger eigenvalue problem, $\Delta \Psi + \varepsilon g \Psi$) do not exist when the inhomogeneity (potential) satisfies $\varepsilon \int g < 0$. See [21] for a proof of this result.

3. Because we rigorously proved the existence of solutions to the intermediate and far-field approximations, we know that we can obtain approximations for $\phi_{\text{far}}$ and $\phi_{\text{inter}}$ to any desired order. Thus, by matching these higher-order approximations, we can obtain better estimates for the parameter $\Lambda$. In particular, if we consider $a = \varepsilon a_1 + \varepsilon^2 a_2$ and find the corresponding expressions for $\phi_{\text{far}}, \phi_{\text{inter}}$ and $a_2$, the above matching process leads to $\Lambda = C(\varepsilon)2e^{-\gamma_c} \exp(1/\varepsilon a_1)$, with $C(\varepsilon) = \exp(1/a - 1/\varepsilon a_1)$. In addition, by defining $\Lambda(0) = \partial_r \Lambda(0) = 0$, we obtain that this estimate is also $C^1$ with respect to $\varepsilon$ on $[0, \varepsilon_M)$, for some $\varepsilon_M > 0$.

### 5.2. Existence of solutions

In this subsection, we combine the results of the previous subsections and prove Theorem 1, which is stated in introduction and reproduced below for convenience.

**Theorem.** Let $k \geq 2$ and $\sigma \in (0, 1)$ and consider a function $g \in H^k_{r, \sigma}(\mathbb{R}^2)$ satisfying Hypothesis 1.1. Then, there exist a constant $\varepsilon_0 > 0$ and a $C^1([0, \varepsilon_0])$ family of eigenfunctions $\phi = \phi(r; \varepsilon)$ and eigenvalues $\Omega = \Omega(\varepsilon)$ that bifurcate from zero and solve the equation

$$
\Delta_0 \phi - b(\partial_r \phi)^2 - \varepsilon g(r) + \Omega = 0 \quad r = |x| \in [0, \infty).
$$

Moreover, this family has the form

$$
\phi(r; \varepsilon) = -\frac{1}{b} \chi(\Lambda r) \log(K_0(\Lambda r)) + \phi_1(r; \varepsilon) + \varepsilon c, \quad \Lambda^2 = b\Omega(\varepsilon)
$$

where

- i) $c$ is a constant that depends on the initial conditions of the problem,
- ii) $K_0(z)$ represents the zeroth-order modified Bessel function of the second kind,
- iii) $\partial_r \phi_1 \in H^k_{r, \sigma}(\mathbb{R}^2)$, and
- iv) $\Omega = \Omega(\varepsilon) = C(\varepsilon)4e^{-2\gamma_c} \exp[2/a]$, with

$$
a = -\varepsilon b \int_0^{\infty} g_c(r) r \, dr,
$$

and $C(\varepsilon)$ a $C^1$ constant that depends on $\varepsilon$.

**Proof.** The proof mimics the analysis done for the far-field approximation, except that now we consider the full equation (16). As above, we use the ansatz $\phi(r) = \phi_0(r) + \phi_1(r)$, with $\phi_0$ given by

$$
\phi_0(r) = -\frac{1}{b} \chi(\Lambda r) \log(K_0(\Lambda r)), \quad \Lambda^2 = b\Omega > 0.
$$
In contrast to the analysis from Sect. 4, here we treat the parameter $\Lambda$ as a $C^1$ function of $\varepsilon$, a result that follows from the matched asymptotic analysis of Sect. 5.1. Thus, given any $\varepsilon > 0$ there is a corresponding value of $\Lambda$ that defines an approximation, $\phi_0$, and a frequency, $\Omega = \Lambda^2/b$, both of which satisfy the equation $(\Delta_0 \phi_0 - b(\partial_r \phi_0)^2 + \Omega) = 0$ in the far field.

Inserting this ansatz into Eq. (16) gives
\[
\Delta_0 \psi_1 - 2b\partial_r \phi_0 \partial_r \psi_1 - b(\partial_r \psi_1)^2 + (\Delta_0 \phi_0 - b(\partial_r \phi_0)^2 + \Omega) - \varepsilon g = 0.
\]
Letting $\psi = \partial_r \phi_1$, adding and subtracting the term $2\Lambda \psi$, and preconditioning the result by $L_{2\Lambda}^{-1}$ give the following equivalent formulation of Eq. (16),
\[
F(\psi; \varepsilon) = \text{Id} + L_{2\Lambda}^{-1} \left[ -2b\partial_r \phi_0 \psi - b\psi^2 + (\Delta_0 \phi_0 - b(\partial_r \phi_0)^2 + \Omega) - \varepsilon g - 2\Lambda \psi \right] = 0.
\]
Our goal is to show that the operator $F : H^k_{r,\sigma}(\mathbb{R}^2) \times \mathbb{R} \to H^k_{r,\sigma}(\mathbb{R}^2)$ satisfies the conditions of the implicit function theorem.

By Remark 5.1, $\Lambda(0) = \partial_r \Lambda(0) = 0$, so that the operator $F$ is $C^1([0, \varepsilon_M])$ with respect to $\varepsilon$, for some $\varepsilon_M > 0$. Moreover, thanks to the cut-off function in the definition of $\phi_0$, i.e., $\chi = \chi(\Lambda r)$, we find that the terms $(\Delta_0 \phi_0 - b(\partial_r \phi_0)^2 + \Omega)$ tend to zero as $\varepsilon$ goes to zero. Therefore, $F(0; 0) = 0$. In addition, because the elements in the parenthesis are smooth and have compact support, they belong to the space $H^k_{r,\sigma}((\mathbb{R}^2)$, for any natural number $k$ and any real number $\sigma$. A similar analysis as in the proof of Theorem 3 then shows that the rest of the terms in $F$ belong to the space $H^k_{r,\sigma}(\mathbb{R}^2)$, with $k \geq 2$, and $\sigma \in (0, 1)$. As a result, the operator $F$ is also well defined. Since its Fréchet derivative, $D_{\psi}F(0; 0)$, is now the identity map on $H^k_{r,\sigma}(\mathbb{R}^2)$, we may apply the implicit function theorem to conclude the existence of solutions $\psi = \partial_r \phi \in H^k_{r,\sigma}(\mathbb{R}^2)$. The results of Theorem 1 then follow in a similar way as those done in Sect. 4. □

6. Simulations

In this section, we numerically explore the effects of adding large inhomogeneities, $\varepsilon g$, as perturbations to the eikonal equation, i.e.,
\[
\phi_t = \Delta \phi - |\nabla \phi|^2 - \varepsilon g.
\]
To run the simulations, we model the equation on a square domain with periodic boundary conditions and employ a spectral RK4 method based on [14], using a mesh size $h = 100/512$ and a time step $dt = 0.5$. The numerical scheme is continued until a steady state is reached. Different domain lengths were tested, $(L = \{100, 120, 140, 160, 180, 200\})$, resulting in the same approximations for $\phi$. Thus, a domain of length $L = 100$ was chosen to run all numerical experiments for computational efficiency.

Simulations confirm our analytical results, finding that for inhomogeneities that take the form
\[
g = \left( \frac{A}{(1 + r^2)^p} \right), \quad p \in (1/2, 1], \quad A \in \mathbb{R},
\]
the solutions to the eikonal equation grow linearly at infinity. This is depicted in Fig. 1a where the gradient, $\nabla \phi$, is plotted for different values of the parameter $\varepsilon$. Notice that because we are using periodic boundary conditions, the value of $\nabla \phi$ goes to zero at the boundary of the domain. As predicted by the analysis of the previous sections, we find that the wavenumber, $k = \lim_{|x| \to \infty} \nabla \phi \sim \Lambda/b$, and as a result the frequency, $\Omega = \Lambda^2/b$, is small beyond all orders of $\varepsilon$. To confirm this result, we approximate the wavenumber by evaluating the gradient $\nabla \phi$ at large values of $|x|$. In Fig. 1b, we plot the relation $\log(\Lambda) = A \log(k) - B$ vs. $a$, where $a$ represents the mass of $g_e = (1 - \chi)g$, which we take as a substitute for $\varepsilon$, since $a = -b \int g_e(r)r \, dr$. Notice how in the figure the data points taken from the simulations follow a straight line, confirming that $\Lambda \sim \exp(1/a)$.

Finally, to determine how the decay rate, $p$, affects the wavenumber, we ran simulations for values of $p \in (0.5, 3]$. Notice that using the notation from Hypothesis 1.1, where $g \sim 1/r^m$, this is equivalent to
Fig. 1. Numerical simulation of the time-dependent eikonal equation (18) with \( g \) as in (19), \( p = 0.8 \), initial condition \( \phi = 0 \), and various values of \( \varepsilon \sim a = 0.15 \times (3 : 2 : 20) \), where \( a = \int_0^3 g(r) r \, dr \). 

A) Plots of the gradient of the steady-state solution. 

B) Plot of \( \frac{1}{\text{log}(k(a)) - 1} \) vs. \( a \), where for large \( |x| \) the gradient \( \nabla \phi \) approximates the wavenumber, \( k(a) \). Circles represent data from the simulation, while dotted line is the linear fit considering values of \( m \in (1, 6) \). These results are summarized in Fig. 2a. They show that the wavenumber decreases as the decay rate of the inhomogeneity, \( p \), increases. The figure also compares the numerical approximation to the wavenumber, \( k \), which we plot using stars, with the analytical result \( k \sim \exp(1/a) \). 

In particular, following Theorem 1 we use

\[
a = \begin{cases} 
-\varepsilon b \int_0^3 g(r) r \, dr & \text{for } p \in (1/2, 1) \\
-\varepsilon b \int_0^\infty g(r) r \, dr & \text{for } p \in (1, 3).
\end{cases}
\]

For values of \( p \in (0.5, 1) \sim m \in (1, 2) \), we are in the regime considered in this paper, where the impurity \( g \) does not have finite mass and is thus a large inhomogeneity. In this case, we assume a value of \( D \sim 3 \) in the definition of \( g_c \) specified in the introduction, see Eq. (3). We then calculate the mass of this function by integrating from 0 to 3, since this provided the best fit to the data (see Remark 1.5). This approximation is plotted using a dashed line. On the other hand, when \( p \in (1, 3) \) we are in the regime where the impurity, \( g \), has finite mass and the results from [13] apply (see also Theorem 1 and Remark 1.4). In this case, the value of \( a = -\varepsilon b \int_0^\infty g(r) r \, dr \). This approximation is plotted using a solid line.

Notice that both approximations for the wavenumber, \( k \), do a good job of following the data in the respective regions of the \( p \)-axis where they are valid, i.e., \( 0.5 < p < 1 \) for the dashed line and \( p > 2 \) for the solid line. However, the estimates for \( p \in (1, 2) \) using the mass of \( g \) (solid line) are not accurate, even though they follow the results from Theorem 1 in [13], or equivalently, Theorem 1 together with Remark 1.4 stated in this paper. This is not unreasonable given that the frequency of the pattern, \( \Omega \), and as a result its wavenumber, \( k \), are both small beyond all orders of the parameter \( \varepsilon \). In particular, when \( p \to 1 \) we have that \( a = -\varepsilon b \int g(r) r \, dr \to -\infty \). Because \( a \sim O(\varepsilon) \), the estimates for \( \Omega \sim \exp(1/a) \) become worse and worse, and in this case, one needs to approximate \( a \) to higher orders in \( \varepsilon \) to obtain better estimates.
Figure 2a then suggests that the interval $1 < p < 2$ is a transitional regime, where one can numerically obtain a better fit to the data by using a cut-off function to better approximate the value of $a$.

Finally, we also confirm numerically that for values of $p \leq 0.5$ the inhomogeneity no longer produces target patterns, but rather solutions with $\nabla \phi \sim O(r)$ at infinity, see Fig. 2b. This is not a tight bound on the growth rate of $\nabla \phi$ and is just a very rough estimate based on our numerical experiments.

7. Discussion

In this paper, we showed that large defects can generate target patterns in oscillatory media. Under the assumption of weak coupling, we modeled such systems using a viscous eikonal equation, and represented the defect as a localized inhomogeneity. In contrast to previous results, which assume that the inhomogeneity is strongly localized, in this paper we relaxed this assumption and described impurities as functions with algebraic decay of order $O(1/|x|^m)$, $1 < m \leq 2$.

Our main motivation for studying this problem came from the universality of the viscous eikonal equation as a model for the phase dynamics of coherent structures in oscillatory media. In particular, our interest stems from the fact that this same equation can be used to describe the phase dynamics of spiral waves in oscillatory media with nonlocal coupling. In this context, the large inhomogeneity no longer represents a defect, but instead encodes information about variations in the amplitude of the pattern.

A second motivation came from the fact that the steady-state viscous eikonal equation is conjugate to a Schrödinger eigenvalue problem. Indeed, it is well known that the Hopf–Cole transformation maps target pattern solutions to bound states of the corresponding Schrödinger operator and that the frequency of target pattern solutions then corresponds to the energy of these states. In this context, the results presented here expand the conditions on the Schrödinger potential that allow for such bound states to exist. In particular, we show that Schrödinger operators with potentials that decay sufficiently fast at infinity can have bound states even when the mass of the potential $\int g(r) r \, dr$ is not finite.
In particular, our analysis provides a first-order approximation for target pattern solutions and for their frequency. In agreement with simulations, we show that, just as in the case of small defects, the frequency is small beyond all orders of the small parameter used to describe the strength of the impurity. As a result, solutions do not follow a regular expansion. Therefore, to obtain our results we first found intermediate and far-field approximations to the steady-state viscous eikonal equation. Then, using a matched asymptotic analysis we were able to determine the value of the frequency selected by the system. This approach is similar in spirit to the one used to prove existence of target patterns and spiral waves in reaction–diffusion equations using spatial dynamics, [15,20]. There, the modeling equations are viewed as a system of ordinary differential equation in the radial variable, and a center manifold reduction is used to obtain a vector field describing the amplitude of these patterns. Coherent structures then correspond to heteroclinic solutions, connecting a fixed point at infinity with solutions that are bounded near the origin. Our matching process is then equivalent to showing that the center-stable manifold of the fixed point intersects transversely the solution curve that lives in the center manifold.

Finally, the analysis presented in this paper is complemented by simulations of the viscous eikonal equation. Our numerical experiments are in good agreement with simulations. They confirm that the wavenumber and therefore the frequency of target patterns do not follow a regular expansion on the small parameter $\varepsilon$ representing the strength of the impurity $g$. They also confirm that when $m \leq 1$, the solutions to the viscous eikonal equation no longer represent target patterns, since in this case the gradient $\nabla \phi$ does not approach a constant as $|x| \to \infty$.

**Declarations**

**Conflict of interest** The author declares that she has no conflict of interest.

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**Appendix**

In [10], it was shown that the following amplitude equation governs the dynamics of one-armed spiral waves in nonlocal oscillatory media,

$$0 = \beta \Delta_1 w + \lambda w + \alpha |w|^2 w + N(w, \varepsilon), \quad r \in [0, \infty).$$

Here, $w$ is a radial and complex-valued function, and

$$\beta = (\sigma - \varepsilon \lambda), \quad \lambda, \alpha \in \mathbb{C}, \quad N \sim O(|\varepsilon||w|^4).$$

It was also established in [10] that the constant $\lambda_I$ is an unknown parameter that needs to be determined when solving the equation.

In this section, a multiple-scale analysis is used to derive a steady-state viscous eikonal equation from the above expression. We will see that this eikonal equation is of the form considered in this paper and that it involves an inhomogeneity that decays at order $O(1/|x|^2)$.

To accomplish this task, we first let $w = A \tilde{w}$, with $A^2 = -\lambda^*_R/\alpha_R$. This change of variables is done for convenience and leads to the following equation,

$$0 = \beta \Delta_1 \tilde{w} + \lambda \tilde{w} + (-\lambda_R + i\tilde{\alpha}_I) |\tilde{w}|^2 \tilde{w} + N(\tilde{w}, \varepsilon), \quad \tilde{\alpha}_I = -\alpha \lambda^*_R/\alpha_R.$$
Letting \( \tilde{w} = \rho e^{i\phi} \) and separating the real and imaginary parts of the above expression, one finally obtains the system

\[
0 = \beta_R \left[ \Delta_1 \rho - (\partial_r \phi)^2 \rho \right] - \beta_I \left[ \Delta_0 \rho + 2\partial_r \phi \partial_r \rho \right] + \lambda_R \rho - \lambda_R \rho^3 + \Re \left[ N(\tilde{w}; \varepsilon) e^{-i\phi} \right] \tag{20}
\]

\[
0 = \beta_R \left[ \Delta_0 \rho + 2\partial_r \phi \partial_r \rho \right] + \beta_I \left[ \Delta_1 \rho - (\partial_r \phi)^2 \rho \right] + \lambda_I \rho + \tilde{\alpha}_I R^3 + \Im \left[ N(\tilde{w}; \varepsilon) e^{-i\phi} \right]. \tag{21}
\]

Next, we proceed with a perturbation analysis following \([4]\). We rescale the variable \( r \) by defining \( S = \delta r \), where \( \delta \) is assumed to be a small positive parameter. We also use the following expressions for the unknown functions:

\[
\rho = \rho_0 + \delta^2 (R_0 + \delta R_1), \quad \rho_0 = \rho_0(r), \quad R_i = R_i(\delta r) \quad i = 0, 1
\]

\[
\phi = \phi_0 + \delta \phi_1, \quad \phi_i = \phi_i(\delta r) \quad i = 0, 1.
\]

And for the parameter, we choose \( \lambda_I = -\tilde{\alpha}_I + \delta^2 \tilde{\lambda}_I \), with \( \tilde{\alpha} \) as above and \( \tilde{\lambda}_I \) a free parameter.

Inserting the above ansatz into Eqs. \((20)\) and \((21)\), we obtain a set of equations in powers of \( \delta \). To write these equations more compactly, we use the subscript \( S \) to distinguish operators that are applied to functions that depend on this variable, i.e., \( \Delta_{0,S} \). The absence of this subscript indicates that the operator is applied to a function of the original variable \( r \).

At order \( O(1) \), we find that \( \rho_0 \) must satisfy,

\[
0 = \beta_R \Delta_1 \rho_0 + \lambda_R \rho_0 - \lambda_R \rho_0^3,
\]

\[
0 = \beta_I \Delta_1 \rho_0 - \tilde{\alpha}_I \rho_0 + \tilde{\alpha}_I \rho_0^3.
\]

At the next order, \( O(\delta^2) \), we find two equations involving \( R_0 \) and \( \phi_0 \),

\[
0 = -\beta_I \rho_0 \Delta_{0,S} \phi_0 - 2\beta_I \partial_S \phi_0 \partial_S \rho_0 - \beta_R \rho_0 (\partial_S \phi_0)^2 + \lambda_R R_0 (1 - 3\rho_0^2),
\]

\[
0 = \beta_R \rho_0 \Delta_{0,S} \phi_0 + 2\beta_R \partial_S \phi_0 \partial_S \rho_0 - \beta_I \rho_0 (\partial_S \phi_0)^2 + \tilde{\alpha}_I R_0 (3\rho_0^2 - 1) - \tilde{\lambda}_I \rho_0.
\]

For our purposes, it is enough to stop at this stage and not list higher-order terms.

We first focus on the order \( O(1) \) system. The first equation can be solved, provided \( \beta_R, \lambda_R > 0 \). This equation falls into a broader family of o.d.e. which were solved in \([16]\). In this reference, the authors showed that such equations possess a unique solution \( \rho_* \) satisfying

\[
\rho_* \rightarrow 1 \quad \text{as} \quad r \rightarrow \infty, \quad \rho_*(r) \sim br \quad \text{when} \quad r \sim 0.
\]

Of course, the solution \( \rho_* \) would not satisfy the second equation in the system. So, we let

\[
G = \beta_I \Delta_1 \rho_* - \tilde{\alpha}_I \rho_* + \tilde{\alpha}_I \rho_*^3 = \left( \frac{\beta_I}{\beta_R} \lambda_R + \tilde{\alpha}_I \right) \rho_* (\rho_*^2 - 1),
\]

and add these terms to the order \( O(\delta^2) \) system.

Going back to the order \( O(\delta^2) \) system, we first notice that because \( \rho_0 = \rho_* \sim br = bS/\delta \) near the origin, then the terms that involve this variable are in fact ‘large’ when compared to the terms that do not. Concentrating only on these large terms, we find that in the first equation we can solve for \( R_0 \) in terms of the variable \( \phi_0 \). Inserting this result into the second equation gives us the viscous eikonal equation,

\[
\Delta_{0,S} \phi_0 - b(\partial_S \phi_0)^2 + \Omega - cg = 0
\]

as expected, where

\[
b = \frac{\beta_I \lambda_R - \beta_R \tilde{\alpha}_I}{\tilde{\alpha}_I \beta_R + \lambda_R \beta_I}, \quad \Omega = \frac{\tilde{\lambda}_I \lambda_R}{\tilde{\alpha}_I \beta_I + \lambda_R \beta_R}, \quad c = -\frac{\beta_I \lambda_R + \tilde{\alpha}_I \beta_R}{\beta_R (\tilde{\alpha}_I \beta_I + \lambda_R \beta_R)}, \quad g = (1 - \rho_*^2).
\]

Numerical simulations show that the perturbation \( g \) decays at order \( O(1/r^2) \) as \( r \) goes to infinity, see Fig. 3. To obtain these results, we solved the boundary value problem

\[
0 = \partial_{rr} \rho + \frac{1}{r} \partial_r \rho - \frac{1}{r^2} \rho + \rho - \rho^3, \quad \rho(\infty) = 1, \quad \rho(0) = 0, \tag{22}
\]
Fig. 3. Solution to the boundary value problem (22) treating the equation as a system of o.d.e. and using a shooting method with condition

\[ \rho(r) \sim br \quad \text{when} \quad r \sim 0. \]

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