# Eigenvalue bifurcation in doubly nonlinear problems with an application to surface plasmon polaritons 

Tomáš Dohnal© and Giulio Romani®


#### Abstract

We consider a class of generally non-self-adjoint eigenvalue problems which are nonlinear in the solution as well as in the eigenvalue parameter ("doubly" nonlinear). We prove a bifurcation result from simple isolated eigenvalues of the linear problem using a Lyapunov-Schmidt reduction and provide an expansion of both the nonlinear eigenvalue and the solution. We further prove that if the linear eigenvalue is real and the nonlinear problem $\mathcal{P T}$-symmetric, then the bifurcating nonlinear eigenvalue remains real. These general results are then applied in the context of surface plasmon polaritons (SPPs), i.e. localized solutions for the nonlinear Maxwell's equations in the presence of one or more interfaces between dielectric and metal layers. We obtain the existence of transverse electric SPPs in certain $\mathcal{P} \mathcal{T}$-symmetric configurations.


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## 1. Introduction

We study the nonlinear problem

$$
\begin{equation*}
L(x, \omega) \varphi:=A \varphi-W(x, \omega) \varphi=f(x, \omega, \varphi), \quad x \in \mathbb{R}^{d} \tag{1}
\end{equation*}
$$

where $A: L^{2}\left(\mathbb{R}^{d}, \mathbb{C}\right) \supset D(A) \rightarrow L^{2}\left(\mathbb{R}^{d}, \mathbb{C}\right)$ is a densely defined, closed (possibly non-self-adjoint) operator with a non-empty resolvent set. Throughout the paper the space $L^{2}\left(\mathbb{R}^{d}, \mathbb{C}\right)$ and all other function spaces are complex vector spaces, i.e. defined over the complex field $\mathbb{C}$. The potential $W$ is generally nonlinear in the spectral parameter $\omega$ and typically complex valued. The function $f$ is nonlinear in both $\omega$ and $\varphi$ and is asymptotically equivalent to a monomial
near $\varphi=0$. Moreover, we suppose that $f$ is Lipschitz continuous in a neighbourhood of an eigen-pair $\left(\omega_{0}, \varphi_{0}\right)$, where $\omega_{0} \in \mathbb{C}$ is a simple isolated eigenvalue of $L(x, \cdot)$. We prove the bifurcation from $\omega_{0}$ using a fixed point argument and a Lyapunov-Schmidt decomposition. Bifurcation from simple eigenvalues is a well studied problem even in the non-selfadjoint case $[6,7,15,18]$. In particular, bifurcation in complex Banach spaces (as relevant in our problem) is investigated in $[7,15]$ by means of a Lyapunov-Schmidt reduction coupled with topological degree techniques. However, our result includes also an asymptotic expansion of $(\omega, \varphi)$ depending on the behaviour of $f$ for small $\varphi$ and for $\omega$ near $\omega_{0}$. More precisely, we find a solution $(\omega, \varphi) \in \mathbb{C} \times D(A)$ of the form

$$
\omega=\omega_{0}+\varepsilon \nu+\varepsilon^{1+\tau} \sigma, \quad \varphi=\varepsilon^{\alpha} \varphi_{0}+\varepsilon^{\alpha+1} \phi+\varepsilon^{\alpha+1+\tau} \psi
$$

where $\varepsilon>0$ is small, $\omega$ is the spectral parameter, $\alpha$ is related to the degree of homogeneity of $f(x, \omega, \cdot)$ near $0, \tau$ is a positive parameter, and $\nu, \sigma \in \mathbb{C}$ as well as $\phi, \psi \in D(A) \cap\left\langle\varphi_{0}^{*}\right\rangle^{\perp}$ are uniquely determined. Moreover, $\nu$ is explicit, see (16), and $\phi$ satisfies the linear equation (20).

In [11] the bifurcation was proved (and an asymptotic expansion of $(\omega, \varphi)$ was provided) for

$$
\begin{equation*}
(A-\omega) \varphi=\varepsilon f(\varphi) \tag{2}
\end{equation*}
$$

with $A$ as above. This problem clearly has a linear dependence on the spectral parameter $\omega$. The coefficient $\varepsilon$ in (2) is the bifurcation parameter and one studies the bifurcation from an eigenvalue $\omega_{0}$ at $\varepsilon=0$. As the bifurcation parameter appears explicitly in the equation, the form of the asymptotic expansion of $(\omega, \varphi)$ is unique. Note that (2) can be rescaled to $(A-\omega) \psi=f(\psi)$ only for the case of homogeneous nonlinearities $f$. Therefore, our result extends that of [11] to the case of more general nonlinearities $f$ and a nonlinear dependence of both $f$ and $W$ on the spectral parameter.

An important application of non-selfadjoint problems which are nonlinear in $\omega$ is the propagation of electromagnetic waves in dispersive media, in particular in structures that include a metal. Interfaces of two different media can support localized waves. A typical example is a surface plasmon polariton (SPP) at the interface of a dielectric and a metal, see e.g. [19, 20] or, when more layers of dielectrics and/or metals are considered, $[14,24,26]$. The general case is, of course, described by Maxwell's equations. Assuming the absence of free charges, we have

$$
\begin{equation*}
\mu_{0} \partial_{t} \mathcal{H}=-\nabla \times \mathcal{E}, \quad \epsilon_{0} \partial_{t} \mathcal{D}=\nabla \times \mathcal{H}, \quad \nabla \cdot \mathcal{D}=\nabla \cdot \mathcal{H}=0 \tag{3}
\end{equation*}
$$

where $\mathcal{E}$ and $\mathcal{H}$ is the electric and magnetic field respectively, $\mathcal{D}=\mathcal{D}(\mathcal{E})$ is the electric displacement field and $\epsilon_{0}$ and $\mu_{0}$ are respectively the permittivity and the permeability of the free space. The displacement field $\mathcal{D}$ is generally nonlinear in $\mathcal{E}$ and non-local in time. For odd (e.g. Kerr) nonlinearities and a monochromatic field $(\mathcal{E}, \mathcal{H}, \mathcal{D})(x, y, z, t)=(E, H, D)(x, y, z) e^{\mathrm{i} \omega t}+$ c.c. (with a real frequency $\omega$ ) a nonlinear eigenvalue problem in $(\omega,(E, H)$ ) is obtained if higher harmonics are neglected, for details see Sect. 4. Equation (3) as well as the eigenvalue problem have to be accompanied by interface conditions if an
interface of two media is present. Assuming that the interface is planar and parallel to the $y z$-plane, the interface conditions are

$$
\begin{equation*}
\llbracket E_{2} \rrbracket=\llbracket E_{3} \rrbracket=\llbracket D_{1} \rrbracket=0, \llbracket H \rrbracket=0, \tag{4}
\end{equation*}
$$

where we define (for the interface located at $\left.x=x_{0}\right), \llbracket E_{2} \rrbracket:=E_{2}\left(x \rightarrow x_{0}^{+}\right)-$ $E_{2}\left(x \rightarrow x_{0}^{-}\right)$, etc., see Sect. 4.

A simple example in the cubically nonlinear case is obtained in structures independent of the $y, z$-variables by choosing the transverse electric (TE) ansatz

$$
\begin{equation*}
E(x, y, z)=(0,0, \varphi(x))^{T} e^{\mathrm{i} k y} \tag{5}
\end{equation*}
$$

with $k \in \mathbb{R}$. It leads to the scalar nonlinear problem

$$
\begin{equation*}
\varphi^{\prime \prime}+W(x, \omega) \varphi+\Gamma(x, \omega)|\varphi|^{2} \varphi=0, \quad x \in \mathbb{R} \tag{6}
\end{equation*}
$$

with functions $W, \Gamma: \mathbb{R}^{2} \rightarrow \mathbb{C}$. The interface conditions here boil down to the continuity condition on $\varphi$ and $\varphi^{\prime}$, see Sect. 4.

The bifurcation result provides a curve $\varepsilon \mapsto(\omega(\varepsilon), \varphi(\varepsilon))$ with $\omega(0)=\omega_{0}$ and $\varphi(0)=0$. Even if $\omega_{0}$ is real, the curve can lie in $\mathbb{C} \backslash \mathbb{R}$ for all $\varepsilon \neq 0$. In order for $\omega$ to correspond to the (real) frequency of an electromagnetic field, one needs to ensure that the curve lies in $\mathbb{R}$. As we show in Sect. 3, this is possible by restricting the fixed point argument to a symmetric subspace, namely the $\mathcal{P} \mathcal{T}$-symmetric subspace. $\mathcal{P} \mathcal{T}$-symmetry has been studied extensively in quantum mechanics, see e.g. [4,5]. Recently, a number of physics papers have studied nonlinear $\mathcal{P} \mathcal{T}$-symmetric problems from a phenomenological point of view mainly with emphasis on localized solutions, e.g. [5, 21, 27]. In the context of SPPs, where metals normally lead to a lossy propagation, $\mathcal{P} \mathcal{T}$-symmetry has been applied to obtain lossless propagation, see $[2,3]$. Mathematically, the restriction of a fixed point argument to a $\mathcal{P} \mathcal{T}$-symmetric (or more generally antilinearly symmetric) subspace has been used to obtain real nonlinear eigenvalues, see, e.g., [ $9,11,22$ ].

This article is organized as follows. In Sect. 2 we state and prove our main bifurcation result (Theorem 2.1). The realness of the nonlinear eigenvalue is ensured in the case of $\mathcal{P} \mathcal{T}$-symmetry in Sect. 3. Applications to SPPs are then given in Sect. 4, where 2- and 3-layer-configurations are investigated.

## 2. Bifurcation of nonlinear eigenvalues

In this section we study problem (1), where $A, W$ and $f$ satisfy assumptions (A1)-(f4) below. We prove the existence of a branch of solutions starting from an eigenpair $\left(\omega_{0}, \varphi_{0}\right) \in \mathbb{C} \times D(A)$, i.e. $L\left(x, \omega_{0}\right) \varphi_{0}=0$, such that
i) $\omega_{0}$ is algebraically simple: $\operatorname{ker}\left(L\left(\cdot, \omega_{0}\right)^{2}\right)=\operatorname{ker}\left(L\left(\cdot, \omega_{0}\right)\right)=\left\langle\varphi_{0}\right\rangle$,
ii) $\omega_{0}$ is isolated: there exists $\rho>0$ such that $0 \in \sigma(L(\cdot, \omega))$ for $\omega \in B_{\rho}\left(\omega_{0}\right) \subset$ $\mathbb{C}$ if and only if $\omega=\omega_{0}$.
Notation: Henceforth, the norm and the inner product in the underlying Hilbert space $L^{2}\left(\mathbb{R}^{d}, \mathbb{C}\right)$ will be denoted by $\|\cdot\|$ and $\langle\cdot, \cdot\rangle$ respectively. Moreover,
$\|\cdot\|_{\infty}$ stands for the $L^{\infty}\left(\mathbb{R}^{d}\right)$ norm and $\|\cdot\|_{A}$ for the graph norm, i.e. $\|u\|_{A}:=$ $\|u\|+\|A u\|$. Note that $D(A)=D(L)$ due to assumptions (A2) and (W1).

We define $\varphi_{0}^{*}$ as the ${ }^{1}$ eigenfunction of $\left(L\left(\cdot, \omega_{0}\right)\right)^{*}$ and we suppose that the normalizations of the eigenfunctions $\varphi_{0}$ and $\varphi_{0}^{*}$ are chosen such that

$$
\left\|\varphi_{0}\right\|=1, \quad\left\langle\varphi_{0}, \varphi_{0}^{*}\right\rangle=1 .
$$

The latter normalization is allowed since the simplicity of $\omega_{0}$ ensures that $\left\langle\varphi_{0}, \varphi_{0}^{*}\right\rangle \neq 0$. Indeed, if $\left\langle\varphi_{0}, \varphi_{0}^{*}\right\rangle=0$, then $\varphi_{0} \in \operatorname{Ker}\left(L\left(\cdot, \omega_{0}\right)\right) \cap \operatorname{Ran}\left(L\left(\cdot, \omega_{0}\right)\right)$ since $\operatorname{Ker}\left(L\left(\cdot, \omega_{0}\right)^{*}\right)^{\perp}=\operatorname{Ran}\left(L\left(\cdot, \omega_{0}\right)\right)$. From $L\left(\cdot, \omega_{0}\right) \psi=\varphi_{0}$ for some $\psi \in$ $\mathrm{D}(A)$ we get $L\left(\cdot, \omega_{0}\right)^{2} \psi=0$ and the algebraic simplicity in i) implies $\psi=c \varphi_{0}$ and hence $\varphi_{0}=0$, which is a contradiction.
We consider the following assumptions on the operator $A$ and the potential $W$ : there exists $\delta>0$ such that
A1) $A: D(A) \rightarrow L^{2}\left(\mathbb{R}^{d}, \mathbb{C}\right)$ is a densely defined, closed operator with a nonempty resolvent set;
A2) $D(A) \hookrightarrow L^{\infty}\left(\mathbb{R}^{d}, \mathbb{C}\right)$, where the embedding is continuous;
W1) $W: \mathbb{R}^{d} \times \mathbb{C} \rightarrow \mathbb{C}$ satisfies that $W(x, \cdot)$ is holomorphic on $B_{\delta}\left(\omega_{0}\right) \subset \mathbb{C}$ for a.e. $x \in \mathbb{R}^{d}$ and there exists $c>0$ such that

$$
\|W(\cdot, \omega)\|_{\infty} \leq c \quad \forall \omega \in B_{\delta}\left(\omega_{0}\right)
$$

and the technical assumption
Wt) $\left\langle\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}, \varphi_{0}^{*}\right\rangle \neq 0$.
Regarding the nonlinearity $f=f(x, \omega, \varphi)$, we assume that there are $\delta, \alpha, \lambda_{0}>$ 0 such that
f1) $f(\cdot, \omega, \varphi(\cdot)) \in L^{2}\left(\mathbb{R}^{d}, \mathbb{C}\right)$ for all $\omega \in B_{\delta}\left(\omega_{0}\right)$ and $\varphi \in B_{\delta}(0) \subset D(A)$,
f2) there exists a constant $K_{f}^{\omega}\left(\varphi_{0}\right)>0$ such that for any $\lambda \in\left(0, \lambda_{0}\right)$ there holds

$$
\begin{aligned}
& \left\|f\left(\cdot, \omega_{1}, \lambda \varphi\right)-f\left(\cdot, \omega_{2}, \lambda \varphi\right)\right\| \\
& \quad \leq K_{f}^{\omega}\left(\varphi_{0}\right) \lambda^{\frac{1}{\alpha}+1}\left|\omega_{1}-\omega_{2}\right|, \quad \forall \omega_{1}, \omega_{2} \in B_{\delta}\left(\omega_{0}\right), \forall \varphi \in B_{\delta}\left(\varphi_{0}\right) ;
\end{aligned}
$$

f3) there exists a constant $K_{f}^{\varphi}\left(\omega_{0}\right)>0$ such that for any $\lambda \in\left(0, \lambda_{0}\right)$ there holds

$$
\begin{aligned}
& \left\|f\left(\cdot, \omega, \lambda \varphi_{1}\right)-f\left(\cdot, \omega, \lambda \varphi_{2}\right)\right\| \\
& \quad \leq K_{f}^{\varphi}\left(\omega_{0}\right) \lambda^{\frac{1}{\alpha}+1}\left\|\varphi_{1}-\varphi_{2}\right\|, \quad \forall \varphi_{1}, \varphi_{2} \in B_{\delta}\left(\varphi_{0}\right), \forall \omega \in B_{\delta}\left(\omega_{0}\right) ;
\end{aligned}
$$

f4) for any $\omega \in B_{\delta}\left(\omega_{0}\right)$ there exists $g_{\omega} \in L^{\infty}\left(\mathbb{R}^{d}, \mathbb{C}\right) \backslash\{0\}$ such that

$$
\begin{equation*}
f(\cdot, \omega, \varphi)-g_{\omega}(\cdot)|\varphi|^{\frac{1}{\alpha}} \varphi=\mathcal{O}\left(|\varphi|^{\frac{1}{\alpha}+1+\beta}\right) \quad \text { as } \quad \mathbb{C} \ni \varphi \rightarrow 0 \tag{7}
\end{equation*}
$$

for some $\beta>0$, uniformly wrt $\omega \in B_{\delta}\left(\omega_{0}\right)$. Moreover, we assume that $g_{\omega}$ is Lipschitz in $\omega$ for $\omega \in B_{\delta}\left(\omega_{0}\right)$ uniformly wrt $x \in \mathbb{R}^{d}$, i.e. there exists a constant $K_{g}>0$ such that

$$
\left\|g_{\omega_{1}}-g_{\omega_{2}}\right\|_{\infty} \leq K_{g}\left|\omega_{1}-\omega_{2}\right| \quad \forall \omega_{1}, \omega_{2} \in B_{\delta}\left(\omega_{0}\right) .
$$

[^0]Theorem 2.1. Suppose that $\omega_{0}$ is an algebraically simple and isolated eigenvalue of $L$ with eigenfunction $\varphi_{0}$ and that $A, W$ and $f$ satisfy assumptions (A1)-(f4). Let also $\tau \in(0, \min \{1, \alpha \beta\}]$. Then there exists $\varepsilon_{0}>0$ such that for any $\varepsilon \in\left(0, \varepsilon_{0}\right)$ there is a unique solution $(\omega, \varphi)$ of the nonlinear eigenvalue problem (1) with the constraint $\left\langle\varphi, \varphi_{0}^{*}\right\rangle=\varepsilon^{\alpha}$. It can be written as

$$
\begin{equation*}
\omega=\omega_{0}+\varepsilon \nu+\varepsilon^{1+\tau} \sigma, \quad \varphi=\varepsilon^{\alpha} \varphi_{0}+\varepsilon^{\alpha+1} \phi+\varepsilon^{\alpha+1+\tau} \psi \tag{8}
\end{equation*}
$$

with $\nu, \sigma \in \mathbb{C}$ and $\phi, \psi \in D(A) \cap\left\langle\varphi_{0}^{*}\right\rangle^{\perp}$.
Before going into the details of the proof, let us give some remarks and examples of applications.

Remark 1. The proof of Theorem (2.1) is constructive. Therefore, we also know which problems the constants $\nu, \sigma$ and the functions $\phi, \psi$ satisfy. In particular, $\nu$ and $\phi$ are uniquely determined by (16), (20) and $(\sigma, \psi)$ uniquely solves the nonlinear system (17), (21) in $B_{r_{1}}(0) \times B_{r_{2}}(0) \subset \mathbb{C} \times D(A)$, for suitable $r_{1}, r_{2}=O(1)$ as $\varepsilon \rightarrow 0$.

Remark 2. When $D(A)$ is the Sobolev space $H^{s}\left(\mathbb{R}^{d}, \mathbb{C}\right)$, assumption (A2) is equivalent to the requirement $s>d / 2$.

Remark 3. Note that assumption (W1) ensures that

$$
\begin{equation*}
\left\|\partial_{\omega}^{k} W\left(\cdot, \omega_{0}\right)\right\|_{\infty} \leq c<\infty \quad \forall k \in \mathbb{N} . \tag{9}
\end{equation*}
$$

Indeed, by Taylor's theorem for holomorphic functions (see Chapt. 4, Sect. 3.1 in [1]) we have

$$
\begin{equation*}
W(\cdot, \omega)=\sum_{j=0}^{n-1} \partial_{\omega}^{j} W\left(\cdot, \omega_{0}\right)\left(\omega-\omega_{0}\right)^{j}+T_{n}(\cdot, \omega)\left(\omega-\omega_{0}\right)^{n} \tag{10}
\end{equation*}
$$

where

$$
T_{n}(\cdot, \omega)=\frac{1}{2 \pi \mathrm{i}} \int_{\partial B_{r}\left(\omega_{0}\right)} \frac{W(\cdot, z)}{\left(z-\omega_{0}\right)^{n}(z-\omega)} d z
$$

for all $\omega \in B_{r}\left(\omega_{0}\right)$ and any $0<r<\delta$. One can easily estimate

$$
\begin{align*}
\left\|T_{n}(\cdot, \omega)\right\|_{\infty} & \leq \frac{\operatorname{ess} \sup _{x \in \mathbb{R}^{d}}^{r^{n-1}\left(r-\left|\omega-\omega_{0}\right|\right)} \max _{z \in \partial B_{r}\left(\omega_{0}\right)}|W(x, z)|}{\text { for all } \omega \in B_{r^{\prime}}\left(\omega_{0}\right) \text { if } 0<r^{\prime}<r} \\
& \leq \frac{2 M}{r^{n}} \quad \text { for all } \omega \in B_{r / 2}\left(\omega_{0}\right), \tag{11}
\end{align*}
$$

where

$$
M:=\sup _{\omega \in B_{\delta}\left(\omega_{0}\right)}\|W(\cdot, \omega)\|_{\infty}
$$

Note that $\omega \in B_{r / 2}\left(\omega_{0}\right)$ is satisfied if $\varepsilon>0$ is small enough. To estimate $\left\|\partial_{\omega}^{k} W\left(\cdot, \omega_{0}\right)\right\|_{\infty}, k \in \mathbb{N}$, one proceeds by induction using (10), assumption (W1), and (11).

Remark 4. Assumption (Wt) is the classical transversality condition for the bifurcation from a simple eigenvalue [6, Theorem 1], [8, Theorem 28.6] in the case

$$
\begin{align*}
& f(x, \cdot, \cdot) \in C^{1}(\mathbb{C} \times \Omega), \partial_{\omega} \partial_{\varphi} f(x, \cdot, \cdot) \in C(\mathbb{C} \times \Omega) \\
& \quad \text { for some neighborhood } \Omega \subset D(A) \text { of zero. } \tag{12}
\end{align*}
$$

However, note that in our setting the nonlinearity $f(\cdot, \omega, \varphi) \sim g_{\omega}(\cdot)|\varphi|^{\frac{1}{\alpha}} \varphi$ for $\varphi \rightarrow 0$ (and with $\varphi \in D(A)$ with a suitable $D(A)$, e.g. $D(A)=H^{2}(\mathbb{R}, \mathbb{C})$ for $d=1$ ) is differentiable only at zero. This is due to the fact that our function spaces are defined over the complex field.

Nevertheless, given the differentiability property, if (12) holds, (Wt) is equivalent to $F_{\omega \varphi}\left(\omega_{0}, 0\right)\left[\varphi_{0}\right] \notin \operatorname{Ran}\left(F_{\varphi}\left(\omega_{0}, 0\right)\right)$, where $F(\omega, \varphi):=L(\cdot, \omega) \varphi$ $f(\cdot, \omega, \varphi)$. To see this, first note that

$$
F_{\varphi}\left(\omega_{0}, 0\right)=A-W\left(\cdot, \omega_{0}\right)
$$

since $\partial_{\varphi} f(\cdot, \omega, 0)=0$ for all $\omega \in B_{\delta}\left(\omega_{0}\right)$ by (f4), and

$$
F_{\omega \varphi}\left(\omega_{0}, 0\right)\left[\varphi_{0}\right]=-\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}-\partial_{\omega} \partial_{\varphi} f\left(\cdot, \omega_{0}, 0\right) \varphi_{0}=-\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}
$$

By the closed range theorem is $\left(A-W\left(\cdot, \omega_{0}\right)\right) u=-\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}$ solvable if and only if $(\mathrm{Wt})$ is violated. Here we have used the fact that $A-W\left(\cdot, \omega_{0}\right)$ is a Fredholm operator (in particular, $\operatorname{Ran}(A)$ is closed), see [16, Theorem IV.5.28].
Remark 5. (Example of the potential and the nonlinearity.) Assumptions (A1)(f4) are satisfied for instance by equation (6) with $D(A)=H^{2}(\mathbb{R}, \mathbb{C})$ provided (W1) and (Wt) hold and $\Gamma$ satisfies

$$
\left\|\Gamma\left(\cdot, \omega_{1}\right)-\Gamma\left(\cdot, \omega_{2}\right)\right\| \leq L\left|\omega_{1}-\omega_{2}\right|
$$

for some $L>0$ and all $\omega_{1,2} \in B_{\delta}\left(\omega_{0}\right)$. An example corresponding to the Drude model for metals (see Sect. 4) is

$$
W(x, \omega)=\frac{\omega^{2}}{c^{2}}\left(1-\frac{\omega_{p}^{2}}{\omega^{2}+\mathrm{i} \gamma \omega}\right)-k^{2}, \quad \Gamma(x, \omega)=3 \frac{\omega^{2}}{c^{2}} \hat{\chi}^{(3)}(x, \omega)
$$

with parameters $\gamma, \omega_{p}, k \in \mathbb{R}$ and $\hat{\chi}^{(3)}$ a bounded function, Lipschitz continuous in $\omega$. This choice will be important for modelling SPPs. Note that, with such potential $W$, the operator $L(\cdot, \omega):=\partial_{x}^{2}+W(\cdot, \omega)$ is non-selfadjoint and also nonlinear in $\omega$.

Remark 6. Let us discuss the role of the parameter $\tau$ in the expansion (8). Clearly, $\tau$ determines the accuracy of the expansion given by the first two terms. According to Theorem 2.1, if $\alpha \beta \leq 1$, then the optimal value is $\tau=\alpha \beta$, which is proportional to the difference of the degree of the lowest degree term in $f$ and the next term. Notice also that higher order terms in the nonlinearity do not play any role in the choice of $\tau$. To give an example, consider the nonlinearities in the table below.

As explained in Sect. 4, in the applications of Theorem 2.1 to time harmonic electromagnetic waves, the relevant nonlinearities are odd. In the case of a cubic nonlinearity $(\alpha=1 / 2)$ the first correction term in $f$ is of the kind $|\varphi|^{4} \varphi$ and therefore we have $\beta=2$ such that we are in the optimal case $\tau=1$.

| $f(x)$ | $\alpha$ | $\beta$ | $\tau_{\max }$ | $\omega$ |
| :--- | :--- | :--- | :--- | :--- |
| $\|\varphi\|^{2} \varphi$ | $\frac{1}{2}$ | $+\infty$ | 1 | $\omega_{0}+\varepsilon \nu+\varepsilon^{2} \sigma$ |
| $\|\varphi\|^{2} \varphi+\|\varphi\|^{4} \varphi$ | $\frac{1}{2}$ | 2 | 1 | $\omega_{0}+\varepsilon \nu+\varepsilon^{2} \sigma$ |
| $\|\varphi\|^{2} \varphi+\|\varphi\|^{3} \varphi$ | $\frac{1}{2}$ | 1 | $\frac{1}{2}$ | $\omega_{0}+\varepsilon \nu+\varepsilon^{\frac{3}{2}} \sigma$ |
| $\|\varphi\|^{2} \varphi+\|\varphi\|^{3} \varphi+\|\varphi\|^{4} \varphi$ | $\frac{1}{2}$ | 1 | $\frac{1}{2}$ | $\omega_{0}+\varepsilon \nu+\varepsilon^{\frac{3}{2}} \sigma$ |

The strategy of the proof of Theorem 2.1 may be summarized as follows. We employ a Lyapunov-Schmidt reduction making use of spectral projections (here the simplicity of the eigenvalue $\omega_{0}$ is used) to decompose the problem into a system in which one rather easily determines $\nu$ and $\phi$. Then, the rest becomes a system for the unknowns $\sigma$ and $\psi$, which will be solved by means of a nested fixed-point argument. In particular, the assumption that $\omega_{0}$ be isolated is exploited to invert the operator $L\left(\cdot, \omega_{0}\right)$ restricted to $\left\langle\varphi_{0}^{*}\right\rangle^{\perp}$.

## Proof of Theorem 2.1. Lyapunov-Schmidt decomposition.

In this initial step we reformulate problem (1) with the ansatz in (8) as a system of two equations using the Lyapunov-Schmidt decomposition.

Let us first introduce the projections $P_{0}: u \mapsto\left\langle u, \varphi_{0}^{*}\right\rangle \varphi_{0}$ and $Q_{0}:=$ $I d-P_{0}$. Clearly, $P_{0}: L^{2}\left(\mathbb{R}^{d}, \mathbb{C}\right) \rightarrow\left\langle\varphi_{0}\right\rangle$ and $Q_{0}: L^{2}\left(\mathbb{R}^{d}, \mathbb{C}\right) \rightarrow\left\langle\varphi_{0}^{*}\right\rangle^{\perp}$. Using our constraint $\left\langle\varphi, \varphi_{0}^{*}\right\rangle=\varepsilon^{\alpha}$, it is easy to see that $P_{0}\left(\phi+\varepsilon^{\tau} \psi\right)=0$. Applying then $P_{0}$ to our Eq. (1), we get

$$
\begin{align*}
\left\langle L(\cdot, \omega) \varphi, \varphi_{0}^{*}\right\rangle= & \left\langle f(\cdot, \omega, \varphi), \varphi_{0}^{*}\right\rangle \\
= & \left.\left.\left\langle f(\cdot, \omega, \varphi)-\varepsilon^{\alpha+1} g_{\omega_{0}}(\cdot)\right| \varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}, \varphi_{0}^{*}\right\rangle  \tag{13}\\
& \left.+\left.\varepsilon^{\alpha+1}\left\langle g_{\omega_{0}}(\cdot)\right| \varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}, \varphi_{0}^{*}\right\rangle .
\end{align*}
$$

Notice that the first term is of higher-order in $\varepsilon$, mainly because of our assumption (f4), see the forthcoming computations (27)-(31).

On the other hand, Taylor expanding $W(\cdot, \omega)$ in $\omega_{0}$ up to order two and using assumption (W1), we have

$$
\begin{aligned}
& \left\langle L(\cdot, \omega) \varphi, \varphi_{0}^{*}\right\rangle \\
& \quad=\left\langle L\left(\cdot, \omega_{0}\right) \varphi, \varphi_{0}^{*}\right\rangle+\left\langle\left(L(\cdot, \omega)-L\left(\cdot, \omega_{0}\right)\right) \varphi, \varphi_{0}^{*}\right\rangle \\
& \quad=\left\langle\varphi, L\left(\cdot, \omega_{0}\right)^{*} \varphi_{0}^{*}\right\rangle-\left\langle\left(W(\cdot, \omega)-W\left(\cdot, \omega_{0}\right)\right) \varphi, \varphi_{0}^{*}\right\rangle \\
& \quad=-\left\langle\left(\partial_{\omega} W\left(\cdot, \omega_{0}\right)\left(\omega-\omega_{0}\right)+\frac{1}{2} \partial_{\omega}^{2} W\left(\cdot, \omega_{0}\right)\left(\omega-\omega_{0}\right)^{2}+I(\cdot, \omega)\right) \varphi, \varphi_{0}^{*}\right\rangle,
\end{aligned}
$$

where

$$
I(x, \omega):=\frac{\left(\omega-\omega_{0}\right)^{3}}{2 \pi \mathrm{i}} \int_{\partial B_{r}\left(\omega_{0}\right)} \frac{W(x, z)}{\left(z-\omega_{0}\right)^{3}(z-\omega)} d z
$$

for any $r<\delta$, see (10). Inserting now the expansions of $\omega$ and $\varphi$ from (8), i.e.

$$
\omega=\omega_{0}+\varepsilon \nu+\varepsilon^{1+\tau} \sigma \quad \text { and } \quad \varphi=\varepsilon^{\alpha} \varphi_{0}+\varepsilon^{\alpha+1} \phi+\varepsilon^{\alpha+1+\tau} \psi
$$

we obtain

$$
\begin{align*}
-\left\langle L(\cdot, \omega) \varphi, \varphi_{0}^{*}\right\rangle= & \left.\varepsilon^{\alpha+1} \nu\left\langle\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}, \varphi_{0}^{*}\right\rangle+\varepsilon^{\alpha+1+\tau} \sigma\left\langle\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right), \varphi_{0}^{*}\right\rangle \\
& +\varepsilon^{\alpha+2}\left\langle\left(\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \phi+\frac{\nu^{2}}{2} \partial_{\omega}^{2} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right), \varphi_{0}^{*}\right\rangle  \tag{14}\\
& +\left\langle v(\cdot, \omega, \varphi), \varphi_{0}^{*}\right\rangle
\end{align*}
$$

where $v$ collects all other terms of higher-order in $\varepsilon$, namely

$$
\begin{align*}
v(\cdot, \omega, \varphi):= & \varepsilon^{\alpha+2+\tau} \sigma \partial_{\omega} W\left(\cdot, \omega_{0}\right) \phi+\varepsilon^{\alpha+2+\tau}\left(\nu+\varepsilon^{\tau} \sigma\right) \partial_{\omega} W\left(\cdot, \omega_{0}\right) \psi \\
& +\varepsilon^{\alpha+3} \frac{\nu^{2}}{2} \partial_{\omega}^{2} W\left(\cdot, \omega_{0}\right)\left(\phi+\varepsilon^{\tau} \psi\right)  \tag{15}\\
& +\varepsilon^{\alpha+2+\tau} \frac{1}{2} \partial_{\omega}^{2} W\left(\cdot, \omega_{0}\right)\left(2 \nu \sigma+\varepsilon^{\tau} \sigma^{2}\right)\left(\varphi_{0}+\varepsilon \phi+\varepsilon^{1+\tau} \psi\right) \\
& +\varepsilon^{\alpha} I(\cdot, \omega)\left(\varphi_{0}+\varepsilon \phi+\varepsilon^{1+\tau} \psi\right)
\end{align*}
$$

Comparing now (13) and (14), the terms of order $\alpha+1$ in $\varepsilon$ match if and only if we take

$$
\begin{equation*}
\nu:=-\frac{\left.\left.\left\langle g_{\omega_{0}}(\cdot)\right| \varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}, \varphi_{0}^{*}\right\rangle}{\left\langle\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}, \varphi_{0}^{*}\right\rangle}, \tag{16}
\end{equation*}
$$

which is well-defined thanks to assumption (Wt). From the rest of (13)-(14) we obtain

$$
\begin{align*}
\varepsilon^{\alpha+1+\tau} & \sigma\left\langle\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}, \varphi_{0}^{*}\right\rangle \\
= & -\varepsilon^{\alpha+2}\left\langle\left(\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \phi+\frac{\nu^{2}}{2} \partial_{\omega}^{2} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right), \varphi_{0}^{*}\right\rangle  \tag{17}\\
& \left.-\left\langle v(\cdot, \omega, \varphi), \varphi_{0}^{*}\right\rangle-\left.\left\langle f(\cdot, \omega, \varphi)-\varepsilon^{\alpha+1} g_{\omega_{0}}(\cdot)\right| \varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}, \varphi_{0}^{*}\right\rangle .
\end{align*}
$$

Let us now apply $Q_{0}$ to (1). On the one hand we have

$$
\begin{aligned}
Q_{0} L(\cdot, \omega) Q_{0} \varphi= & -Q_{0}\left(W(\cdot, \omega)-W\left(\cdot, \omega_{0}\right)\right)\left(\varepsilon^{\alpha+1} \phi+\varepsilon^{\alpha+1+\tau} \psi\right) \\
& +Q_{0} L\left(\cdot, \omega_{0}\right)\left(\varepsilon^{\alpha+1} \phi+\varepsilon^{\alpha+1+\tau} \psi\right)
\end{aligned}
$$

and on the other hand

$$
\begin{aligned}
Q_{0} L(\cdot, \omega) Q_{0} \varphi & =Q_{0} L(\cdot, \omega) \varphi-\varepsilon^{\alpha} Q_{0}\left(L(\cdot, \omega)-L\left(\cdot, \omega_{0}\right)\right) \varphi_{0}-\varepsilon^{\alpha} Q_{0} L\left(\cdot, \omega_{0}\right) \varphi_{0} \\
& =Q_{0} f(\cdot, \omega, \varphi)+Q_{0}\left(W(\cdot, \omega)-W\left(\cdot, \omega_{0}\right)\right)\left(\varepsilon^{\alpha} \varphi_{0}\right)
\end{aligned}
$$

Therefore, we obtain

$$
\begin{aligned}
& Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0}\left(\varepsilon^{\alpha+1} \phi+\varepsilon^{\alpha+1+\tau} \psi\right) \\
& \quad=Q_{0} f(\cdot, \omega, \varphi)+\varepsilon^{\alpha} Q_{0}\left(W(\cdot, \omega)-W\left(\cdot, \omega_{0}\right)\right)\left(\varphi_{0}+\varepsilon \phi+\varepsilon^{1+\tau} \psi\right)
\end{aligned}
$$

An expansion of the last term as in (14)-(15) yields

$$
\begin{align*}
& \varepsilon^{\alpha+1} Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0} \phi+\varepsilon^{\alpha+1+\tau} Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0} \psi \\
& \quad=\varepsilon^{\alpha+1} Q_{0}\left(g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right)+Q_{0}\left(f(\cdot, \omega, \varphi)-\varepsilon^{\alpha+1} g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right) \\
& \quad+Q_{0}\left(\varepsilon^{\alpha+1} \nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}+\varepsilon^{\alpha+2}\left(\frac{\nu^{2}}{2} \partial_{\omega}^{2} W\left(\cdot, \omega_{0}\right) \varphi_{0}+\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \phi\right)\right. \\
& \left.\quad+\varepsilon^{\alpha+1+\tau} \sigma \partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}+v(\cdot, \omega, \varphi)\right) \tag{18}
\end{align*}
$$

Rewriting (16) as

$$
\begin{aligned}
& \left.g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}-\left.\left\langle g_{\omega_{0}}(\cdot)\right| \varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}, \varphi_{0}^{*}\right\rangle \varphi_{0} \\
& \quad=g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}+\nu\left\langle\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}, \varphi_{0}^{*}\right\rangle \varphi_{0}
\end{aligned}
$$

and using (18), we obtain

$$
\begin{align*}
& \varepsilon^{\alpha+1} Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0} \phi+\varepsilon^{\alpha+1+\tau} Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0} \psi \\
& \quad=\varepsilon^{\alpha+1}\left(g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}+\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right) \\
& \quad+Q_{0}\left(f(\cdot, \omega, \varphi)-\varepsilon^{\alpha+1} g_{\omega_{0}}\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right)+\varepsilon^{\alpha+1+\tau} \sigma Q_{0}\left(\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right)  \tag{19}\\
& \quad+\varepsilon^{\alpha+2} Q_{0}\left(\frac{\nu^{2}}{2} \partial_{\omega}^{2} W\left(\cdot, \omega_{0}\right) \varphi_{0}+\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \phi\right)+Q_{0}(v(\cdot, \omega, \varphi))
\end{align*}
$$

Again, imposing that the terms of the lowest-order in $\varepsilon$ match, we get a linear equation for $\phi$ :

$$
\begin{equation*}
Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0} \phi=g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}+\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0} \tag{20}
\end{equation*}
$$

Notice that, with our choice of $\nu$, Eq. (20) is uniquely solvable in $Q_{0} D(A)=$ $D(A) \cap\left\langle\varphi_{0}^{*}\right\rangle^{\perp}$ by the closed range theorem. Indeed, the operator on the left hand side is Fredholm (see [16, Theorem IV.5.28], where the fact that $\omega_{0}$ is a simple isolated eigenvalue in the above sense is used) and the right hand side is orthogonal to the kernel of the adjoint operator, i.e. to $\varphi_{0}^{*}$, due to (16).
The rest of (19) produces the following equation for $(\sigma, \psi)$ :

$$
\begin{align*}
\varepsilon^{\alpha+1+\tau} Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0} \psi= & Q_{0}\left(f(\cdot, \omega, \varphi)-\varepsilon^{\alpha+1} g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right) \\
& +\varepsilon^{\alpha+1+\tau} \sigma Q_{0}\left(\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right) \\
& +\varepsilon^{\alpha+2} Q_{0}\left(\frac{\nu^{2}}{2} \partial_{\omega}^{2} W\left(\cdot, \omega_{0}\right) \varphi_{0}+\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \phi\right) \\
& +Q_{0}(v(\cdot, \omega, \varphi)) \\
= & : R_{1}+\varepsilon^{\alpha+1+\tau} R_{2}+\varepsilon^{\alpha+2} R_{3}+R_{4}=: R(\sigma, \psi) \tag{21}
\end{align*}
$$

## Fixed Point Argument

In order to solve our initial problem (1), we now need to solve system (17), (21) for $(\sigma, \psi) \in \mathbb{C} \times D(A)$. Inserting then $\sigma$ and $\psi$ into (8) produces a solution $(\omega, \varphi)$ of (1).

We proceed by a fixed point argument. Note that although a direct fixed point argument for $(\sigma, \psi)$ is possible, we opt for a nested version, where we first solve for $\psi$ as a function of $\sigma$ and subsequently solve for $\sigma$. This approach is arguably more transparent.

Let us first address Eq. (21) and write it as a fixed point equation for $\psi$, exploiting our assumptions on the eigenvalue $\omega_{0}$, which is assumed simple and isolated. This actually means that $Q_{0} L\left(\omega_{0}\right) Q_{0}: Q_{0} D(A) \rightarrow Q_{0} L^{2}\left(\mathbb{R}^{d}\right)$ is boundedly invertible in $Q_{0} D(A)$ and its norm is bounded by a constant $C\left(\omega_{0}\right)$ :

$$
\begin{align*}
& \psi=\varepsilon^{-(\alpha+1+\tau)}\left[Q_{0} L\left(\omega_{0}\right) Q_{0}\right]^{-1} R(\sigma, \psi):=G(\sigma, \psi), \\
& \text { with } \quad\|G(\sigma, \psi)\|_{A} \leq C\left(\omega_{0}\right) \varepsilon^{-(\alpha+1+\tau)}\|R(\sigma, \psi)\| . \tag{22}
\end{align*}
$$

In our nested fixed point argument for (17), (22) we first solve (22) for $\psi \in$ $B_{r_{2}}(0) \subset D(A)$ for all $\sigma \in B_{r_{1}}(0) \subset \mathbb{C}$ fixed with $r_{1}>0$ arbitrary. To this aim we need to show that for each $\sigma \in B_{r_{1}}(0)$ there exists $r_{2}>0$ so that ${ }^{2}$
(i) $\psi \in B_{r_{2}}(0) \Rightarrow G(\psi) \in B_{r_{2}}(0)$,
(ii) $\exists \rho \in(0,1):\left\|G\left(\psi_{1}\right)-G\left(\psi_{2}\right)\right\|_{A} \leq \rho\left\|\psi_{1}-\psi_{2}\right\|_{A}$ for all $\psi_{1}, \psi_{2} \in B_{r_{2}}(0)$
if $\varepsilon>0$ is small enough.
Then, having obtained $\psi=\psi(\sigma) \in B_{r_{2}}(0)$, we shall solve Eq. (17) for $\sigma$ and finally find a suitable $r_{1}$. Note that the fixed point argument for equation (17) requires the Lipschitz continuity of $\sigma \mapsto \psi(\sigma)$, which we verify below.

To ensure $(i)$, we need to estimate $\|R(\psi)\|$. The second and the third term in (21) are easy to handle. Henceforth, we track the dependence of all constants on $\sigma$ and $\psi$ via $r_{1}, r_{2}$.

$$
\begin{gather*}
\left\|R_{2}\right\| \leq|\sigma|\left\|\partial_{\omega} W\left(\cdot, \omega_{0}\right)\right\|_{\infty} \leq C r_{1}  \tag{23}\\
\left\|R_{3}\right\| \leq \max \left\{\frac{|\nu|^{2}}{2}, \nu\|\phi\|\right\}\left(\left\|\partial_{\omega} W\left(\cdot, \omega_{0}\right)\right\|_{\infty}+\left\|\partial_{\omega}^{2} W\left(\cdot, \omega_{0}\right)\right\|_{\infty}\right)=C \tag{24}
\end{gather*}
$$

using (9). Let us now deal with $R_{4}$. Inspecting (15), we obtain

$$
\begin{align*}
\left\|R_{4}\right\| \leq & \|v(\cdot, \omega, \varphi)\| \leq C\left[\varepsilon^{\alpha+2+\tau}\left(r_{1}+\left(1+\varepsilon^{\tau} r_{1}\right)\|\psi\|\right)+\varepsilon^{\alpha+3}\left(1+\varepsilon^{\tau}\|\psi\|\right)\right. \\
& +\varepsilon^{\alpha+2+\tau} r_{1}\left(1+\varepsilon^{\tau} r_{1}\right)\left(1+\varepsilon^{1+\tau}\|\psi\|\right) \\
& \left.+\varepsilon^{\alpha+3}\left(1+\varepsilon^{\tau} r_{1}\right)\left(1+\varepsilon^{1+\tau}\|\psi\|\right)\right]  \tag{25}\\
\leq & C\left[\varepsilon^{\alpha+2+\tau}\left(r_{1}+\|\psi\|\right)+\varepsilon^{\alpha+3}\right]+\varepsilon^{\alpha+2+2 \tau} h_{1}\left(r_{1},\|\psi\|\right)
\end{align*}
$$

where $h_{1}$ is polynomial in $r_{1}$, linear in $\|\psi\|$ and satisfies $h_{1}(0,\|\psi\|)=0$. Actually, all terms appearing in (15) are easy to estimate, so here we just briefly justify the one for the integral rest $I(x, \omega)$, for later use too. Indeed, using (11)

[^1]for $n=3$, we have
\[

$$
\begin{equation*}
\|I(\cdot, \omega)\|_{\infty} \leq \frac{2 M}{r^{3}}\left|\omega-\omega_{0}\right|^{3}=\frac{2 M}{r^{3}} \varepsilon^{3}\left|\nu+\varepsilon^{\tau} \sigma\right|^{3} \leq \frac{2 M}{r^{3}} \varepsilon^{3}\left(1+\varepsilon^{\tau} r_{1}\right) \tag{26}
\end{equation*}
$$

\]

for any $0<r<\delta$ and all $\omega \in B_{r / 2}\left(\omega_{0}\right)$, i.e. for all $\varepsilon>0$ small enough. Recall that $M=\sup _{\omega \in B_{\delta}\left(\omega_{0}\right)}\|W(\cdot, \omega)\|_{\infty}$. Finally, we need to estimate $R_{1}$, which involves the nonlinearity. We split it as

$$
\begin{align*}
f(\cdot, \omega, \varphi)-\varepsilon^{\alpha+1} g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}= & \left(f(\cdot, \omega, \varphi)-g_{\omega}|\varphi|^{\frac{1}{\alpha}} \varphi\right) \\
& +g_{\omega}\left(|\varphi|^{\frac{1}{\alpha}} \varphi-\varepsilon^{\alpha+1}\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right)  \tag{27}\\
& +\varepsilon^{\alpha+1}\left(g_{\omega}-g_{\omega_{0}}(\cdot)\right)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}
\end{align*}
$$

and estimate term by term. First,

$$
\begin{align*}
\left\|f(\cdot, \omega, \varphi)-g_{\omega}|\varphi|^{\frac{1}{\alpha}} \varphi\right\| & \leq C\left(\omega_{0}\right)\left\||\varphi|^{\frac{1}{\alpha}+1+\beta}\right\| \leq C\|\varphi\|_{\infty}^{\frac{1}{\alpha}+\beta}\|\varphi\| \leq C\|\varphi\|_{A}^{\frac{1}{\alpha}+1+\beta} \\
& \leq C \varepsilon^{\alpha+1+\alpha \beta}\left\|\varphi_{0}+\varepsilon \phi+\varepsilon^{1+\tau} \psi\right\|_{A}^{\frac{1}{\alpha}+1+\beta}  \tag{28}\\
& \leq C \varepsilon^{\alpha+1+\alpha \beta}\left(1+\varepsilon^{1+\tau} r_{2}\right)^{\frac{1}{\alpha}+\beta}\left(1+\varepsilon^{1+\tau}\|\psi\|_{A}\right) \\
& \leq C_{1}\left(r_{2}\right) \varepsilon^{\alpha+1+\alpha \beta}\left(1+\varepsilon^{1+\tau}\|\psi\|_{A}\right)
\end{align*}
$$

for $\varepsilon>0$ small enough. In Eq. (28) we used the embedding $D(A) \hookrightarrow L^{\infty}\left(\mathbb{R}^{d}\right)$. The dependence $r_{2} \mapsto C_{1}\left(r_{2}\right)$ is of power type; in detail,
$C_{1}\left(r_{2}\right)=C_{0}\left(1+\varepsilon^{1+\tau} r_{2}\right)^{1 / \alpha+\beta}=C_{0}\left(1+\left(\frac{1}{\alpha}+\beta\right) \varepsilon^{1+\tau} r_{2}\left(1+\varepsilon^{1+\tau} r_{*}\right)^{1 / \alpha+\beta-1}\right)$
for some $r_{*} \in\left[0, r_{2}\right]$. For each $r_{2}>0$ there exists $\varepsilon_{0}=\varepsilon_{0}\left(r_{2}\right)$ such that

$$
\left(\frac{1}{\alpha}+\beta\right) \varepsilon^{1+\tau} r_{2}\left(1+\varepsilon^{1+\tau} r_{*}\right)^{1 / \alpha+\beta-1} \leq 1
$$

for all $\varepsilon \in\left(0, \varepsilon_{0}\right)$. In conclusion

$$
\begin{equation*}
C_{1}\left(r_{2}\right) \leq 2 C_{0}=2 C_{1}(0) \tag{29}
\end{equation*}
$$

for all $\varepsilon \in\left(0, \varepsilon_{0}\right)$.
Next,
$\left\|g_{\omega}\left(|\varphi|^{\frac{1}{\alpha}} \varphi-\varepsilon^{\alpha+1}\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right)\right\|$

$$
\begin{equation*}
\leq\left\|g_{\omega}\right\|_{\infty} \varepsilon^{\alpha+1}\left\|\left|\varphi_{0}+\varepsilon \phi+\varepsilon^{1+\tau} \psi\right|^{\frac{1}{\alpha}}\left(\varphi_{0}+\varepsilon \phi+\varepsilon^{1+\tau} \psi\right)-\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right\| \tag{30}
\end{equation*}
$$

$$
\leq C \varepsilon^{\alpha+1}\left\|\varepsilon \phi+\varepsilon^{1+\tau} \psi\right\| \leq C \varepsilon^{\alpha+2}\left(1+\varepsilon^{\tau}\|\psi\|\right)
$$

where we used the fact that the map $z \mapsto|z|^{\frac{1}{\alpha}} z$ is locally Lipschitz as well as estimates of the type $\left\||\varphi|^{1 / \alpha} \phi\right\| \leq\|\varphi\|_{\infty}^{1 / \alpha}\|\phi\| \leq c\|\varphi\|_{D(A)}^{1 / \alpha}\|\phi\|=C$. Finally, since $g_{\omega}$ is Lipschitz in $\omega$ by (f4),

$$
\begin{align*}
\varepsilon^{\alpha+1}\left\|\left(g_{\omega}-g_{\omega_{0}}\right)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right\| & \leq \varepsilon^{\alpha+1} C\left\|g_{\omega}-g_{\omega_{0}}\right\|_{\infty} \leq \varepsilon^{\alpha+1} C K_{g}\left|\omega-\omega_{0}\right|  \tag{31}\\
& \leq C \varepsilon^{\alpha+2}\left|\nu+\varepsilon^{\tau} \sigma\right| \leq C \varepsilon^{\alpha+2}\left(1+\varepsilon^{\tau} r_{1}\right)
\end{align*}
$$

Hence from (27)-(31) we infer

$$
\begin{equation*}
\left\|R_{1}\right\| \leq C_{1}\left(r_{2}\right) \varepsilon^{\alpha+1+\alpha \beta}\left(1+\varepsilon^{1+\tau}\|\psi\|_{A}\right)+C \varepsilon^{\alpha+2}\left(1+\varepsilon^{\tau}\left(r_{1}+\|\psi\|\right)\right) \tag{32}
\end{equation*}
$$

and therefore, from (23)-(25) and (32), that

$$
\begin{align*}
\|G(\psi)\|_{A} \leq & C \varepsilon^{-(\alpha+1+\tau)}\|R(\psi)\| \\
\leq & C_{1}\left(r_{2}\right) \varepsilon^{\alpha \beta-\tau}\left(1+\varepsilon^{1+\tau}\|\psi\|_{A}\right)+C\left(\varepsilon^{1-\tau}+\varepsilon\|\psi\|+r_{1}\right)  \tag{33}\\
& +\varepsilon^{1+\tau} h_{1}\left(r_{1},\|\psi\|\right) .
\end{align*}
$$

Using (29) and $\tau \leq \alpha \beta$, we may further estimate

$$
\begin{aligned}
\|G(\psi)\|_{A} & \leq 2 C_{0}\left(1+\varepsilon^{1+\tau} r_{2}\right)+C\left(\varepsilon^{1-\tau}+\varepsilon r_{2}+r_{1}\right)+\varepsilon^{1+\tau} h_{1}\left(r_{1}, r_{2}\right) \\
& \leq 4 C_{0}+\tilde{C} r_{1}
\end{aligned}
$$

for $\varepsilon>0$ small enough. This follows because for $\varepsilon$ small enough we have $\varepsilon^{1+\tau} h_{1}\left(r_{1}, r_{2}\right)<C_{0}$. Setting now $r_{2}=r_{2}\left(r_{1}\right):=4 C_{0}+\tilde{C} r_{1}$, we get $(i)$ provided $\varepsilon \in\left(0, \bar{\varepsilon}_{0}\right)$ with some $\bar{\varepsilon}_{0}=\bar{\varepsilon}_{0}\left(r_{1}\right)>0$.

Let us now address the contraction property (ii). We take $\psi_{1}, \psi_{2} \in B_{r_{2}}(0)$ and define $\varphi_{1,2}:=\varepsilon^{\alpha}\left(\varphi_{0}+\varepsilon \phi+\varepsilon^{1+\tau} \psi_{1,2}\right)$. By (22), we need to estimate $\left\|R\left(\psi_{1}\right)-R\left(\psi_{2}\right)\right\|$, with $R$ defined in (21). First notice that $R_{2}$ and $R_{3}$ do not depend on $\psi$, so they will vanish in the difference and we have to handle just the nonlinear term and $v$.

$$
\begin{align*}
& \left\|v\left(\cdot, \omega, \varphi_{1}\right)-v\left(\cdot, \omega, \varphi_{2}\right)\right\| \\
& \quad \leq \varepsilon^{\alpha+2+\tau} C\left(1+\varepsilon^{\tau} r_{1}\right)\left\|\psi_{1}-\psi_{2}\right\|+\varepsilon^{\alpha+3+\tau}\left\|\psi_{1}-\psi_{2}\right\| \\
& \quad+C \varepsilon^{\alpha+3+2 \tau}\left(r_{1}+\varepsilon^{\tau} r_{1}^{2}\right)\left\|\psi_{1}-\psi_{2}\right\|+C \varepsilon^{\alpha+4+\tau}\left(1+\varepsilon^{\tau} r_{1}\right)\left\|\psi_{1}-\psi_{2}\right\|  \tag{34}\\
& \quad \leq \varepsilon^{\alpha+2+\tau} C_{2}\left(r_{1}\right)\left\|\psi_{1}-\psi_{2}\right\|
\end{align*}
$$

where $C_{2}$ is polynomial in $r_{1}$. Here, estimate (26) was used. Next, by (f3),

$$
\begin{align*}
& \left\|f\left(\cdot, \omega, \varphi_{1}\right)-f\left(\cdot, \omega, \varphi_{2}\right)\right\| \\
& \quad=K_{f}^{\varphi}\left(\omega_{0}\right) \varepsilon^{\alpha+1}\left\|\varphi_{0}+\varepsilon \phi+\varepsilon^{1+\tau} \psi_{1}-\left(\varphi_{0}+\varepsilon \phi+\varepsilon^{1+\tau} \psi_{2}\right)\right\|  \tag{35}\\
& \quad \leq K_{f}^{\varphi}\left(\omega_{0}\right) \varepsilon^{\alpha+2+\tau}\left\|\psi_{1}-\psi_{2}\right\| .
\end{align*}
$$

Hence, by (22), (34) and (35), we get

$$
\begin{aligned}
\left\|G\left(\psi_{1}\right)-G\left(\psi_{2}\right)\right\|_{A} & \leq \varepsilon^{-(\alpha+1+\tau)} C\left\|R\left(\psi_{1}\right)-R\left(\psi_{2}\right)\right\| \\
& \leq \varepsilon \max \left\{K_{f}^{\varphi}\left(\omega_{0}\right), C_{2}\left(r_{1}\right)\right\}\left\|\psi_{1}-\psi_{2}\right\|,
\end{aligned}
$$

which yields $(i i)$ if $\varepsilon$ is small enough. Therefore, applying Banach's fixed point theorem, we infer the existence of a solution $\psi$ of Eq. (21). More precisely, for any $r_{1}>0$ and for any $\sigma \in B_{r_{1}}(0) \subset \mathbb{C}$ there exists $r_{2}=r_{2}\left(r_{1}\right)$ and $\varepsilon_{0}>0$ such that for each $\varepsilon \in\left(0, \varepsilon_{0}\right)$ there is a unique $\psi=\psi(\sigma) \in B_{r_{2}}(0) \subset D(A)$ such that $(\sigma, \psi)$ solves (21).

Let us now address Eq. (17). Inserting $\psi=\psi(\sigma)$ and dividing by the factor $\varepsilon^{\alpha+1+\tau}\left\langle\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}, \varphi_{0}^{*}\right\rangle$, which is nonzero by (Wt), this becomes a
fixed point equation for $\sigma$ :

$$
\begin{align*}
\sigma= & \left(\varepsilon^{\alpha+1+\tau}\left\langle\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}, \varphi_{0}^{*}\right\rangle\right)^{-1} \\
& \cdot\left(-\varepsilon^{\alpha+2}\left\langle\left(\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \phi+\frac{\nu^{2}}{2} \partial_{\omega}^{2} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right), \varphi_{0}^{*}\right\rangle\right. \\
& \left.\left.-\left\langle v(\cdot, \omega, \varphi(\sigma)), \varphi_{0}^{*}\right\rangle-\left.\left\langle f(\cdot, \omega, \varphi(\sigma))-\varepsilon^{\alpha+1} g_{\omega_{0}}(\cdot)\right| \varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}, \varphi_{0}^{*}\right\rangle\right) \\
= & S(\sigma)
\end{align*}
$$

where $\varphi(\sigma)$ is given by (8) with $\psi=\psi(\sigma)$. We need to show that for some $r_{1}>0$
$\left(i^{\prime}\right) S: B_{r_{1}}(0) \rightarrow B_{r_{1}}(0)$,
( $i i^{\prime}$ ) there exists $\rho \in(0,1)$ so that

$$
\left|S\left(\sigma_{1}\right)-S\left(\sigma_{2}\right)\right| \leq \rho\left|\sigma_{1}-\sigma_{2}\right| \quad \text { for all } \quad \sigma_{1}, \sigma_{2} \in B_{r_{1}}(0)
$$

if $\varepsilon>0$ is small enough.
Notice that the first term on the right-hand side of (17') is independent of $\sigma$ and $\psi$, therefore its norm can be simply estimated by a constant. Decomposing the nonlinear term as in (27), according to estimates (28)-(31), we get

$$
\begin{align*}
\left\|f(\cdot, \omega, \varphi)-\varepsilon^{\alpha+1} g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right\| \leq & C_{1}\left(r_{2}\right) \varepsilon^{\alpha+1+\alpha \beta}\left(1+\varepsilon^{1+\tau} r_{2}\right) \\
& +C \varepsilon^{\alpha+2}\left(1+\varepsilon^{\tau}\left(|\sigma|+r_{2}\right)\right) \tag{36}
\end{align*}
$$

Moreover, similarly to (25) we have

$$
\begin{align*}
\|v(\cdot, \omega, \varphi)\| \leq & \varepsilon^{\alpha+2+\tau} C|\sigma|+\varepsilon^{\alpha+2+\tau}\left(1+\varepsilon^{\tau}|\sigma|\right) r_{2}+C \varepsilon^{\alpha+3+\tau} r_{2} \\
& +\varepsilon^{\alpha+2+\tau} C\left(|\sigma|+\varepsilon^{\tau}|\sigma|^{2}\right)\left(1+\varepsilon^{1+\tau} r_{2}\right)  \tag{37}\\
& +\varepsilon^{\alpha+3} C\left(1+\varepsilon^{\tau}|\sigma|\right)\left(1+\varepsilon^{1+\tau} r_{2}\right) \\
\leq & C\left(\varepsilon^{\alpha+2+\tau}\left(|\sigma|+r_{2}\right)+\varepsilon^{\alpha+3}\right)+\varepsilon^{\alpha+2+2 \tau} h_{1}\left(|\sigma|, r_{2}\right) .
\end{align*}
$$

Therefore, combining (17'), (36) and (37), we obtain

$$
\begin{aligned}
|S(\sigma)| \leq & C_{1}\left(r_{2}\right) \varepsilon^{\alpha \beta-\tau}\left(1+\varepsilon^{1+\tau} r_{2}\right)+C \varepsilon^{1-\tau}\left(1+\varepsilon^{\tau}\left(|\sigma|+r_{2}\right)\right) \\
& +C\left(\varepsilon\left(|\sigma|+r_{2}\right)+\varepsilon^{2-\tau}\right)+\varepsilon^{1+\tau} h_{1}\left(|\sigma|, r_{2}\right) .
\end{aligned}
$$

For each $r_{1}>0$ there is $\tilde{\varepsilon}_{0}=\tilde{\varepsilon}_{0}\left(r_{1}\right)>0$ such that

$$
|S(\sigma)| \leq \varepsilon^{\min \{\alpha \beta-\tau, 1-\tau\}}\left(2 C_{1}\left(r_{2}\right)+2 C\right)+1
$$

for all $|\sigma|<r_{1}$ and all $\varepsilon \in\left(0, \tilde{\varepsilon}_{0}\right)$. Recalling that $r_{2}=r_{2}\left(r_{1}\right)=4 C_{0}+\tilde{C} r_{1}$ and that $\tau \leq \min \{1, \alpha \beta\}$, this implies that

$$
|S(\sigma)| \leq 4 C_{0}+2 C+1 \quad \text { for all }|\sigma|<r_{1} \text { and } \varepsilon>0 \text { small enough, }
$$

i.e. $\varepsilon \in\left(0, \varepsilon_{0}\right)$ with $\varepsilon_{0}=\varepsilon_{0}\left(r_{1}\right)$. Setting $r_{1}:=4 C_{0}+2 C+1$, we have property ( $i$ ').

Finally, let us address the contraction property ( $i i^{\prime}$ '). For $\sigma_{1,2} \in B_{r_{1}}(0)$ we define $\psi_{1,2}=\psi\left(\sigma_{1,2}\right)$ and analogously $\omega_{1,2}$ and $\varphi_{1,2}$. We need to estimate $\left|S\left(\sigma_{1}\right)-S\left(\sigma_{2}\right)\right|$. Recalling the form of $v$ in (15), we get

$$
\begin{align*}
&\left\|v\left(\cdot, \omega_{1}, \varphi_{1}\right)-v\left(\cdot, \omega_{2}, \varphi_{2}\right)\right\| \\
& \leq C\left[\varepsilon^{\alpha+2+\tau}\left(\left|\sigma_{1}-\sigma_{2}\right|+\left\|\psi_{1}-\psi_{2}\right\|+\varepsilon^{\tau}\left\|\sigma_{1} \psi_{1}-\sigma_{2} \psi_{2}\right\|\right)\right. \\
&+\varepsilon^{\alpha+3+\tau}\left\|\psi_{1}-\psi_{2}\right\|+\varepsilon^{\alpha+2+\tau}\left(\left|\sigma_{1}^{2}-\sigma_{2}^{2}\right|+\varepsilon^{1+\tau}\left\|\sigma_{1}^{2} \psi_{1}-\sigma_{2}^{2} \psi_{2}\right\|\right)  \tag{38}\\
&\left.+\varepsilon^{\alpha}\left\|I\left(\cdot, \omega_{1}\right)-I\left(\cdot, \omega_{2}\right)\right\|_{\infty}+\varepsilon^{\alpha+1+\tau}\left\|I\left(\cdot, \omega_{2}\right)\right\|_{\infty}\left\|\psi_{1}-\psi_{2}\right\|\right] \\
& \leq C_{3}\left(r_{1}\right)\left(1+r_{2}\right) \varepsilon^{\alpha+2+\tau}\left[\left|\sigma_{1}-\sigma_{2}\right|+\left\|\psi_{1}-\psi_{2}\right\|\right]
\end{align*}
$$

with $C_{3}\left(r_{1}\right)$ cubic in $r_{1}$. Here the estimates are rather standard and quite similar to the ones used in (37), so we just point out how to deal with the rest term $I$. We have, as in Remark 3,

$$
\begin{align*}
I\left(\cdot, \omega_{1}\right)-I\left(\cdot, \omega_{2}\right)= & \frac{1}{2 \pi \mathrm{i}}\left(\left(\omega_{1}-\omega_{0}\right)^{3} \int_{\partial B_{r}\left(\omega_{0}\right)} \frac{W(\cdot, z)}{\left(z-\omega_{0}\right)^{3}\left(z-\omega_{1}\right)} d z\right. \\
& \left.-\left(\omega_{2}-\omega_{0}\right)^{3} \int_{\partial B_{r}\left(\omega_{0}\right)} \frac{W(\cdot, z)}{\left(z-\omega_{0}\right)^{3}\left(z-\omega_{2}\right)} d z\right) \tag{39}
\end{align*}
$$

First, we estimate

$$
\begin{align*}
& \frac{1}{2 \pi}\left\|\int_{\partial B_{r}\left(\omega_{0}\right)} \frac{W(\cdot, z)}{\left(z-\omega_{0}\right)^{3}\left(z-\omega_{1}\right)} d z\left(\left(\omega_{1}-\omega_{0}\right)^{3}-\left(\omega_{2}-\omega_{0}\right)^{3}\right)\right\|_{\infty} \\
& \quad \leq \varepsilon^{3} \frac{2 M}{r^{3}}\left|\left(\nu+\varepsilon^{\tau} \sigma_{1}\right)^{3}-\left(\nu+\varepsilon^{\tau} \sigma_{2}\right)^{3}\right|  \tag{40}\\
& \quad \leq \varepsilon^{3} \frac{2 M}{r^{3}}\left|\varepsilon^{\tau} \nu^{2}\left(\sigma_{1}-\sigma_{2}\right)+3 \varepsilon^{2 \tau} \nu\left(\sigma_{1}^{2}-\sigma_{2}^{2}\right)+\varepsilon^{3 \tau}\left(\sigma_{1}^{3}-\sigma_{2}^{3}\right)\right| \\
& \quad \leq C_{4}\left(r_{1}\right) \varepsilon^{3+\tau}\left|\sigma_{1}-\sigma_{2}\right|,
\end{align*}
$$

which holds for all $\varepsilon>0$ small enough with $C_{4}\left(r_{1}\right)$ quadratic in $r_{1}$ using an estimate analogous to (26). Second, we have

$$
\begin{align*}
& \frac{1}{2 \pi}\left\|\left(\omega_{2}-\omega_{0}\right)^{3}\left(\int_{\partial B_{r}\left(\omega_{0}\right)} \frac{W(\cdot, z)}{\left(z-\omega_{0}\right)^{3}\left(z-\omega_{1}\right)}-\frac{W(\cdot, z)}{\left(z-\omega_{0}\right)^{3}\left(z-\omega_{2}\right)} d z\right)\right\|_{\infty} \\
& \leq c \varepsilon^{3}\left|\nu+\varepsilon^{\tau} \sigma_{2}\right|^{3}\left\|\int_{\partial B_{r}\left(\omega_{0}\right)} \frac{W(\cdot, z)\left(\omega_{1}-\omega_{2}\right)}{\left(z-\omega_{0}\right)^{3}\left(z-\omega_{1}\right)\left(z-\omega_{2}\right)} d z\right\|_{\infty} \\
& \leq \tilde{C}_{5}\left(r_{1}\right) \varepsilon^{4+\tau}\left|\sigma_{1}-\sigma_{2}\right|\left\|\int_{\partial B_{r}\left(\omega_{0}\right)} \frac{W(\cdot, z)}{\left(z-\omega_{0}\right)^{4}\left(z-\omega_{2}\right)} d z\right\|_{\infty} \\
& \leq C_{5}\left(r_{1}\right) \varepsilon^{4+\tau}\left|\sigma_{1}-\sigma_{2}\right| \tag{41}
\end{align*}
$$

for $\varepsilon>0$ small enough, where in the second step we have used $\omega_{1}-\omega_{2}=$ $\varepsilon^{1+\tau}\left(\sigma_{1}-\sigma_{2}\right)$ and estimated $\left|z-\omega_{1}\right| \geq \frac{1}{2}\left|z-\omega_{0}\right|$ for all $z \in \partial B_{r}\left(\omega_{0}\right)$, which holds for $\varepsilon>0$ small enough. In the last step estimate (11) was used again.

The constants $C_{5}\left(r_{1}\right), \tilde{C}_{5}\left(r_{1}\right)$ are cubic in $r_{1}$. Consequently, by (39)-(41) we get

$$
\left\|I\left(\cdot, \omega_{1}\right)-I\left(\cdot, \omega_{2}\right)\right\|_{\infty} \leq \varepsilon^{3+\tau} C_{6}\left(r_{1}\right)\left|\sigma_{1}-\sigma_{2}\right|
$$

with $C_{6}\left(r_{1}\right)$ cubic in $r_{1}$.
Now we have to deal with the third term in (17') involving the nonlinearity. However, this is easily estimated using its Lipschitz behaviour in $\omega$ and $\varphi$ as in (f2)-(f3). Indeed,

$$
\begin{align*}
\left\|f\left(\cdot, \omega_{1}, \varphi_{1}\right)-f\left(\cdot, \omega_{2}, \varphi_{2}\right)\right\| \leq & \left\|f\left(\cdot, \omega_{1}, \varphi_{1}\right)-f\left(\cdot, \omega_{1}, \varphi_{2}\right)\right\| \\
& +\left\|f\left(\cdot, \omega_{1}, \varphi_{2}\right)-f\left(\cdot, \omega_{2}, \varphi_{2}\right)\right\| \tag{42}
\end{align*}
$$

where the first term is estimated as in (35), whereas

$$
\begin{align*}
& \left\|f\left(\cdot, \omega_{1}, \varphi_{2}\right)-f\left(\cdot, \omega_{2}, \varphi_{2}\right)\right\| \\
& \quad \leq \varepsilon^{\alpha+1} K_{f}^{\omega}\left(\varphi_{0}\right)\left|\omega_{1}-\omega_{2}\right|=\varepsilon^{\alpha+2+\tau} K_{f}^{\omega}\left(\varphi_{0}\right)\left|\sigma_{1}-\sigma_{2}\right| . \tag{43}
\end{align*}
$$

Therefore, combining (38) and (42)-(43), we infer

$$
\begin{equation*}
\left|S\left(\sigma_{1}\right)-S\left(\sigma_{2}\right)\right| \leq C_{3}\left(r_{1}\right)\left(1+r_{2}\right) \varepsilon\left(\left\|\psi_{1}-\psi_{2}\right\|+\left|\sigma_{1}-\sigma_{2}\right|\right) \tag{44}
\end{equation*}
$$

with $C_{3}\left(r_{1}\right)$ cubic in $r_{1}$. Hence, it remains to show now that the map $\sigma \mapsto \psi(\sigma)$ is Lipschitz continuous. Taking $\left|\sigma_{1,2}\right| \leq r_{1}$, we shall estimate the difference $\left\|\psi_{1}-\psi_{2}\right\|$ starting from the fixed-point equation (22) for $\psi$, where as before we define $\psi_{1,2}:=\psi\left(\sigma_{1,2}\right)$ and similarly for $\omega_{1,2}$ and $\varphi_{1,2}$. Indeed, exploiting the above estimates (38) and (42)-(43), we have

$$
\begin{align*}
& \left\|G\left(\sigma_{1}, \psi_{1}\right)-G\left(\sigma_{2}, \psi_{2}\right)\right\| \\
& \quad \leq C \varepsilon^{-(\alpha+1+\tau)}\left[\left\|f\left(\cdot, \omega_{1}, \varphi_{1}\right)-f\left(\cdot, \omega_{2}, \varphi_{2}\right)\right\|+C \varepsilon^{\alpha+1+\tau}\left|\sigma_{1}-\sigma_{2}\right|\right. \\
& \left.\quad+\left\|v\left(\cdot, \sigma_{1}, \varphi_{1}\right)-v\left(\cdot, \sigma_{2}, \varphi_{2}\right)\right\|\right]  \tag{45}\\
& \leq \\
& \quad C\left[\left(1+\varepsilon K_{f}^{\omega}\left(\varphi_{0}\right)\right)\left|\sigma_{1}-\sigma_{2}\right|+C_{2}\left(r_{1}\right)\left(1+r_{2}\right) \varepsilon\left(\left|\sigma_{1}-\sigma_{2}\right|+\left\|\psi_{1}-\psi_{2}\right\|\right)\right] \\
& \leq C\left[2\left|\sigma_{1}-\sigma_{2}\right|+C_{2}\left(r_{1}\right)\left(1+r_{2}\right) \varepsilon\left\|\psi_{1}-\psi_{2}\right\|\right]
\end{align*}
$$

if $\varepsilon>0$ is small enough.
This, together with (22), yields $\left\|\psi\left(\sigma_{1}\right)-\psi\left(\sigma_{2}\right)\right\|_{A} \leq 3 C\left|\sigma_{1}-\sigma_{2}\right|$ for $\sigma_{1,2} \in B_{r_{1}}(0)$ and $\varepsilon>0$ small enough. As a consequence we may conclude the fixed point argument for $\sigma$ because from (44) and (45) and from the fact that $r_{2}=r_{2}\left(r_{1}\right)$ we obtain

$$
\left|S\left(\sigma_{1}\right)-S\left(\sigma_{2}\right)\right| \leq \tilde{C}\left(r_{1}\right) \varepsilon\left|\sigma_{1}-\sigma_{2}\right|
$$

with $\tilde{C}\left(r_{1}\right)>0$, which is the desired contraction property for suitably small values of $\varepsilon$. Hence, for small $\varepsilon>0$ the fixed point argument yields the sought solution for the system (17)-(21) and the proof of Theorem 2.1 is complete.

## 3. Bifurcation of real nonlinear eigenvalues: the $\mathcal{P} \mathcal{T}$-symmetric case

In applications one often starts from a real eigenvalue of the linear problem and seeks a bifurcation branch $(\omega, \varphi)$ where the realness of the nonlinear eigenvalue is preserved, i.e. $\omega \in \mathbb{R}$. This is required also for the applications to SPPs in Sect. 4. However, the solution $\omega$ obtained by Theorem 2.1 is a-priori just complex. In this section we provide a symmetric situation in which the realness of $\omega$ is preserved in the bifurcation. This is based on [11, Section III], whereby we adapt the analysis for our more general context. For the sake of completeness and since Sect. 4 is based on these results, we present most of the details here again.

Definition 3.1. A function $\psi: \mathbb{R}^{d} \rightarrow \mathbb{C}(d \geq 1)$ is called $\mathcal{P} \mathcal{T}$-symmetric if $\mathcal{B} \psi(x):=\overline{\psi(-x)}=\psi(x)$ for all $x \in \mathbb{R}^{d}$. Moreover, an operator $L$ acting on a Hilbert space with domain $D(L)$ is $\mathcal{P} \mathcal{T}$-symmetric when it commutes with $\mathcal{B}$, namely $L \mathcal{B}=\mathcal{B} L$ in $D(L)$.

The above operator $\mathcal{B}$ is in fact the composition of the operator $\mathcal{P}$, the space reflection (parity), and $\mathcal{T}$, the complex conjugation, which corresponds the time-reversal in quantum mechanics.

Notice that the Schrödinger operator $-\Delta+W$ with a complex potential $W$ is $\mathcal{P} \mathcal{T}$-symmetric if and only if the real and the imaginary parts of $W$ satisfy respectively $\operatorname{Re} W(-x)=\operatorname{Re} W(x)$ and $\operatorname{Im} W(-x)=-\operatorname{Im} W(x)$ for all $x$. Moreover, general polynomial nonlinearities

$$
f(x, \varphi)=\sum_{p, q=0}^{N} a_{p q}(x) \varphi^{p} \bar{\varphi}^{q}
$$

are $\mathcal{P} \mathcal{T}$-symmetric if and only if the coefficients are so: $\overline{a_{p q}}(-x)=a_{p q}(x)$, see [11, Sections III-IV]. An example is $f(\varphi)=|\varphi|^{2 m+1} \varphi, m \in \mathbb{N}$.

Proposition 3.1. Let $A, W, f$ and $g_{\omega}$ be as in Theorem 2.1 and suppose that they are $\mathcal{P} \mathcal{T}$-symmetric. If $\omega_{0} \in \mathbb{R}$ is an algebraically simple eigenvalue, then for all $\varepsilon \in\left(0, \varepsilon_{0}\right)$ the nonlinear eigenpair $(\omega, \varphi)$ from Theorem 2.1 satisfies $\omega \in \mathbb{R}$ and $\mathcal{B} \varphi=\varphi$.

Remark 7. Under the assumptions on $A, W, f$, and $g_{\omega}$ as in Prop. 3.1 and under the simplicity assumption on $\omega_{0}$, the eigenfunction $\varphi_{0}$ may always be chosen $\mathcal{P} \mathcal{T}$-symmetric. Indeed, $L\left(\cdot, \omega_{0}\right)=A+W\left(\cdot, \omega_{0}\right)$ is $\mathcal{P} \mathcal{T}$-symmetric and, applying $\mathcal{B}$ to $L\left(\cdot, \omega_{0}\right) \varphi_{0}=0$, we get $L\left(\cdot, \omega_{0}\right)\left(\mathcal{B} \varphi_{0}\right)=0$ and we conclude by the simplicity of $\omega_{0}$.

Similarly, we obtain $\mathcal{B} \varphi_{0}^{*}=\varphi_{0}^{*}$ because $\mathcal{B} L^{*}\left(\cdot, \overline{\omega_{0}}\right)=L^{*}\left(\cdot, \overline{\omega_{0}}\right) \mathcal{B}$. Indeed,

$$
\begin{aligned}
&\left\langle v, L^{*} \phi\right\rangle=\langle L v, \phi\rangle=\overline{\langle\overline{(L v)}}(-\cdot), \bar{\phi}(-\cdot)\rangle \\
&=\overline{\langle\mathcal{B} L v, \mathcal{B} \phi\rangle}=\overline{\langle L \mathcal{B} v, \mathcal{B} \phi\rangle} \\
&=\overline{\left\langle\mathcal{B} v, L^{*} \mathcal{B} \phi\right\rangle}=\left\langle v, \overline{L^{*}}(-\cdot) \phi\right\rangle=\left\langle v, \mathcal{B} L^{*} \phi\right\rangle
\end{aligned}
$$

for all $v, \phi \in D(A)$.

Proof. According to our expansions (8), we need to prove that $\nu, \sigma$ are real and that $\phi, \psi$ are $\mathcal{P} \mathcal{T}$-symmetric. First, $\nu$ in (16) satisfies $P_{0}\left(\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right)=$ $P_{0}\left(g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right)$, where we recall that $P_{0}$ is the spectral projection onto the eigenspace $\left\langle\varphi_{0}\right\rangle$ and that with our assumptions $P_{0}$ commutes with $\mathcal{B}$. Indeed,

$$
\begin{aligned}
\mathcal{B} P_{0} u & =\mathcal{B}\left(\left\langle u, \varphi_{0}^{*}\right\rangle \varphi_{0}\right)=\overline{\left\langle u, \varphi_{0}^{*}\right\rangle} \mathcal{B} \varphi_{0}=\left\langle\bar{u}, \overline{\varphi_{0}^{*}}\right\rangle \varphi_{0}=\left\langle\bar{u}, \varphi_{0}^{*}(-\cdot)\right\rangle \varphi_{0} \\
& =\left\langle\overline{u(-\cdot)}, \varphi_{0}^{*}\right\rangle \varphi_{0}=\left\langle\mathcal{B} u, \varphi_{0}^{*}\right\rangle \varphi_{0}=P_{0} \mathcal{B} u .
\end{aligned}
$$

Therefore, on the one hand,

$$
\begin{aligned}
\mathcal{B} P_{0}\left(\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right) & =\bar{\nu} P_{0} \mathcal{B}\left(\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right)=\bar{\nu} P_{0}\left(\partial_{\omega} \mathcal{B} W\left(\cdot, \omega_{0}\right) \mathcal{B} \varphi_{0}\right) \\
& =\bar{\nu} P_{0}\left(\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right)
\end{aligned}
$$

On the other hand,

$$
\begin{aligned}
\mathcal{B} P_{0}\left(\nu \partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}\right) & =\mathcal{B} P_{0}\left(g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right)=P_{0}\left(\mathcal{B}\left(g_{\omega_{0}}(\cdot)\right) \mathcal{B}\left(\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right)\right) \\
& =P_{0}\left(g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right) .
\end{aligned}
$$

This yields $\nu \in \mathbb{R}$. Next, let us analyze $\phi$, i.e. the solution of (20) in $Q_{0} D(A)$. Notice that if $P_{0}$ commutes with $\mathcal{B}$, then the same holds for $Q_{0}=I d-P_{0}$, too. Therefore, applying a similar argument, we get that $\mathcal{B}(\phi)$ satisfies (20) too. Moreover, since $\mathcal{B}(\phi) \in D(A) \cap\left\langle\varphi_{0}^{*}\right\rangle^{\perp}$, we get $\mathcal{B}(\phi)=\phi$ because (20) has a unique solution in $P_{0} D(A)$. Let us now address $\sigma$ and $\psi$. We will show that the coupled fixed point problem (17),(21) preserves the realness of $\sigma$ and the $\mathcal{P} \mathcal{T}$-symmetry of $\psi$. First, given $\sigma \in \mathbb{R}$ with $\sigma \in\left(-r_{1}, r_{1}\right)$, we prove that

$$
\begin{equation*}
\mathcal{B} \psi=\psi \quad \Rightarrow \quad \mathcal{B} G(\sigma, \psi)=G(\sigma, \psi) \tag{46}
\end{equation*}
$$

where $G$ is defined in (22). If (46) holds, then the fixed point $\psi=\psi(\sigma)$ of $\psi=G(\sigma, \psi)$ lies in $B_{r_{2}}(0) \cap\left\{\eta \in L^{2}\left(\mathbb{R}^{d}\right) \mid \mathcal{B} \eta=\eta\right\}$. To this aim, first we notice that for $u \in Q_{0} L^{2}\left(\mathbb{R}^{d}\right)$ it holds

$$
\begin{aligned}
\mathcal{B}\left(Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0}\right)^{-1} u & =\left(Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0}\right)^{-1}\left(Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0}\right) \mathcal{B}\left(Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0}\right)^{-1} u \\
& =\left(Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0}\right)^{-1} \mathcal{B}\left(Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0}\right)\left(Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0}\right)^{-1} u \\
& =\left(Q_{0} L\left(\cdot, \omega_{0}\right) Q_{0}\right)^{-1} \mathcal{B} u .
\end{aligned}
$$

Next, recalling that $\varphi_{0}, \phi$, and $\psi$ are now $\mathcal{P} \mathcal{T}$-symmetric, we get for $R$ (defined in (21))

$$
\begin{aligned}
\mathcal{B} R(\psi)= & Q_{0}\left(\mathcal{B} f(\cdot, \omega, \varphi)-\varepsilon^{\alpha+1} \mathcal{B}\left(g_{\omega_{0}}(\cdot)\left|\varphi_{0}\right|^{\frac{1}{\alpha}} \varphi_{0}\right)\right) \\
& +\varepsilon^{\alpha+1+\tau} \sigma Q_{0}\left(\partial_{\omega} \mathcal{B}\left(W\left(\cdot, \omega_{0}\right)\right) \mathcal{B}\left(\varphi_{0}\right)\right) \\
& +\varepsilon^{\alpha+2} Q_{0}\left(\frac{\nu^{2}}{2} \partial_{\omega}^{2} \mathcal{B}\left(W\left(\cdot, \omega_{0}\right)\right) \mathcal{B}\left(\varphi_{0}\right)+\nu \partial_{\omega} \mathcal{B}\left(W\left(\cdot, \omega_{0}\right)\right) \mathcal{B}(\phi)\right) \\
& +Q_{0}(\mathcal{B} v(\cdot, \omega, \varphi)) \\
= & R(\psi)+Q_{0}(\mathcal{B} v(\cdot, \omega, \varphi)-v(\cdot, \omega, \varphi)) .
\end{aligned}
$$

Inspecting all terms in $v$ appearing in (15) and exploiting $\sigma \in \mathbb{R}$ (and therefore $\omega \in \mathbb{R}$ ), we obtain that $\mathcal{B} v(\cdot, \omega, \varphi)=v(\cdot, \omega, \varphi)$, so (46) is proved. Finally, with similar manipulations one proves that $\mathcal{B}(\psi)=\psi$ implies $\mathcal{B} S(\sigma)=S(\sigma)$, where
$S$ is defined in (17'), obtaining thus $S(\sigma) \in \mathbb{R}$ since $\mathcal{B} S(\sigma)=\overline{S(\sigma)}$. Therefore, for a given $\mathcal{P} \mathcal{T}$-symmetric $\psi(\sigma)$ the fixed point $\sigma$ of $S(\sigma)$ must be real. This completes the proof.

## 4. Applications to nonlinear surface plasmons

As mentioned in the introduction, we are interested in surface plasmon polaritons (SPPs) localized at one or more interfaces between different dielectric and metal layers. We consider the time harmonic and $z$-independent ansatz (5) for the Maxwell system (3). A simple nonlinear, non-local relation for the displacement field is

$$
\begin{align*}
\mathcal{D}(x, y, z, t)= & \mathcal{E}+\int_{\mathbb{R}} \chi^{(1)}(x, y, z, t-s) \mathcal{E}(s) \mathrm{d} s \\
& +\int_{\mathbb{R}} \chi^{(3)}(x, y, z, t-s)((\mathcal{E} \cdot \mathcal{E}) \mathcal{E})(s) \mathrm{d} s \tag{47}
\end{align*}
$$

where $\chi^{(1,3)}: \mathbb{R}^{3} \times \mathbb{R} \rightarrow \mathbb{R}$ and $\chi^{(1,3)}(\cdot, \tau)=0$ for $\tau<0$. The functions $\chi^{(1)}$ and $\chi^{(3)}$ are the linear and the cubic electric susceptibilities of the material, respectively. In general, $\chi^{(1)}$ and $\chi^{(3)}$ are tensors but in the isotropic case, which we assume, the relation in (47) with scalar $\chi^{(1)}$ and $\chi^{(3)}$ holds, see [17, Section 2d].

Substituting a monochromatic ansatz $\mathcal{E}(x, y, z, t)=E(x, y, z) e^{\mathrm{i} \omega t}+$ c.c. in (3) and neglecting higher harmonics (terms proportional to $e^{3 i \omega t}$ and $e^{-3 i \omega t}$ ), we get $\mathcal{D}(x, y, z, t)=D(x, y, z) e^{\mathrm{i} \omega t}+$ c.c. with

$$
D(x, y, z)=\left(1+\hat{\chi}^{(1)}(x, y, z, \omega)\right) E+\hat{\chi}^{(3)}(x, y, z, \omega)\left(2|E|^{2} E+(E \cdot E) \bar{E}\right)
$$

Here $|E|^{2}=E \cdot \bar{E}$ and $\hat{f}(\omega)$ is the Fourier-transform of $f$. Neglecting higher harmonics is a common approach in theoretical studies of weakly nonlinear optical waves [23]. Using (3), we get that both $H$ and $D$ are curl-fields such that the divergence conditions $\nabla \cdot D=0, \nabla \cdot H=0$ hold automatically. Defining $c:=\left(\epsilon_{0} \mu_{0}\right)^{-1 / 2}$, in the second order formulation we have $\frac{\omega^{2}}{c^{2}} D=\nabla \times \nabla \times E$, i.e.

$$
\begin{align*}
\nabla & \times \nabla \times E-\frac{\omega^{2}}{c^{2}}\left(1+\hat{\chi}^{(1)}(x, y, z, \omega)\right) E \\
& -\frac{\omega^{2}}{c^{2}} \hat{\chi}^{(3)}(x, y, z, \omega)\left(2|E|^{2} E+(E \cdot E) \bar{E}\right)=0 \tag{48}
\end{align*}
$$

Recall again only odd nonlinearities are allowed in $\mathcal{D}$ when studying time harmonic waves. Even nonlinearities do not produce terms proportional to $e^{\mathrm{i} \omega t}$.

We consider structures independent of $y$ and $z$, i.e. $\chi^{(1,3)}=\chi^{(1,3)}(x, \omega)$. Interfaces between layers are thus parallel to the $y z$-plane.

For the TE-ansatz in (5), with only one nontrivial component, Eq. (48) reduces to the scalar problem

$$
\begin{equation*}
\varphi^{\prime \prime}+W(x, \omega) \varphi+\Gamma(x, \omega)|\varphi|^{2} \varphi=0, \quad x \in \mathbb{R} \tag{49}
\end{equation*}
$$

with

$$
W(x, \omega):=\frac{\omega^{2}}{c^{2}}\left(1+\hat{\chi}^{(1)}(x, \omega)\right)-k^{2}, \quad \Gamma(x, \omega):=\frac{3 \omega^{2}}{c^{2}} \hat{\chi}^{(3)}(x, \omega)
$$

We study layers of $x$-periodic (including homogeneous) media. When a metal layer is homogeneous, we choose the simplest Drude model of the linear susceptibility in that layer [19]

$$
\begin{equation*}
\hat{\chi}^{(1)}(x, \omega)=-\frac{\omega_{p}^{2}}{\omega^{2}+\mathrm{i} \gamma \omega}, \tag{50}
\end{equation*}
$$

where $\omega_{p} \in \mathbb{R}^{+}$and $\gamma \in \mathbb{R}$. For dielectric layers we choose a periodic (possibly constant) and generally complex $\hat{\chi}^{(1)}(\cdot, \omega)$. The imaginary part of $\hat{\chi}^{(1)}$ is related to the loss or gain of energy of an electromagnetic wave propagating inside the medium. For most of the materials, and in particular for metals, it is negative, corresponding to a lossy material. However, in active (doped) materials the energy of an electromagnetic wave is amplified (energy gain), and therefore the imaginary part of $\hat{\chi}^{(1)}$ is positive. Materials with a real $\hat{\chi}^{(1)}(\cdot, \omega)$ are called conservative.

To make Eq. (49) dimensionless, as well as in order to use the physical values of the parameters involved in the numerical study in Sect. 4.1.2, we introduce the new rescaled spatial variable, frequency, and wave number

$$
\begin{equation*}
\tilde{x}:=\frac{\omega_{p}}{c} x, \quad \tilde{\omega}:=\frac{\omega}{\omega_{p}}, \quad \tilde{k}:=\frac{c}{\omega_{p}} k, \tag{51}
\end{equation*}
$$

where $\omega_{p}$ is the bulk plasma frequency of a prescribed metal layer. Defining then $\tilde{\varphi}(\tilde{x}):=\varphi(x)$, we obtain the same equation as (49) but in the tilde variables and with $W$ and $\Gamma$ respectively replaced by

$$
\begin{align*}
\tilde{W}(\tilde{x}, \tilde{\omega}) & :=\frac{c^{2}}{\omega_{p}^{2}} W\left(\frac{c}{\omega_{p}} \tilde{x}, \omega_{p} \tilde{\omega}\right)=\tilde{\omega}^{2}\left(1+\hat{\chi}^{(1)}\left(\frac{c}{\omega_{p}} \tilde{x}, \omega_{p} \tilde{\omega}\right)\right)-\tilde{k}^{2}  \tag{52}\\
\tilde{\Gamma}(\tilde{x}, \tilde{\omega}) & :=\frac{c^{2}}{\omega_{p}^{2}} \Gamma\left(\frac{c}{\omega_{p}} \tilde{x}, \omega_{p} \tilde{\omega}\right)=3 \tilde{\omega}^{2} \hat{\chi}^{(3)}\left(\frac{c}{\omega_{p}} \tilde{x}, \omega_{p} \tilde{\omega}\right) \tag{53}
\end{align*}
$$

Note that for the Drude model the susceptibility is

$$
\hat{\chi}^{(1)}\left(\frac{c}{\omega_{p}} \tilde{x}, \omega_{p} \tilde{\omega}\right)=-\frac{1}{\tilde{\omega}^{2}+\mathrm{i} \tilde{\gamma} \tilde{\omega}},
$$

where $\tilde{\gamma}:=\gamma / \omega_{p}$.
For the sake of a simpler notation henceforth we will simply write $x, \omega, \gamma$, $k, W$, and $\Gamma$ instead of $\tilde{x}, \tilde{\omega}, \tilde{\gamma}, \tilde{k}, \tilde{W}$, and $\tilde{\Gamma}$.

Due to the presence of material interface(s), solutions of the Maxwell's equations (3) are not smooth. However, they satisfy the interface conditions that the tangential component of $\mathcal{E}$, the normal component of $\mathcal{D}$ and the whole vector $\mathcal{H}$ be continuous across each interface, see Sect. 33-3 in [13]. For our interfaces parallel to the $y z$-plane we get (4). For the ansatz (5) the interface conditions reduce to a $C^{1}$-continuity condition on $\varphi$

$$
\begin{equation*}
\llbracket \varphi \rrbracket=\llbracket \partial_{x} \varphi \rrbracket=0 . \tag{IFCs}
\end{equation*}
$$

Our nonlinear problem (1) is thus Eq. (49) with $W$ and $\Gamma$ respectively replaced by (52)-(53) and coupled with the (IFCs).

To apply our bifurcation result to (49) with (IFCs), a real, linear eigenvalue is needed. We show that such an eigenvalue exists in some $\mathcal{P} \mathcal{T}$-symmetric choices of the layers. SPPs in $\mathcal{P} \mathcal{T}$-symmetric structures have been studied in the physics literature before by employing active materials, see e.g. [2,3,14, 25]. Nevertheless, we are not aware of a rigorous mathematical existence proof of the asymptotic expansion of nonlinear SPPs in the frequency dependent case.

Theorem 2.1 can be applied to (49) with finitely many interfaces (at $\left.x=x_{1}, \ldots, x_{N}\right)$ using the following natural choice of $A, D(A)$, and $f$ :

$$
A:=-\frac{d^{2}}{d x^{2}}, \quad f(x, \omega, \varphi):=\Gamma(x, \omega)|\varphi|^{2} \varphi
$$

and

$$
\begin{aligned}
D(A):= & \left\{\varphi \in L^{2}(\mathbb{R}):\left.\varphi\right|_{\left(x_{j}, x_{j+1}\right)} \in H^{2}\left(\left(x_{j}, x_{j+1}\right), \mathbb{C}\right) \text { for } j=0, \ldots N\right. \\
& \text { and } \left.\llbracket \varphi \rrbracket_{j}=\llbracket \partial_{x} \varphi \rrbracket_{j}=0, j=1, \ldots, N\right\},
\end{aligned}
$$

where $\llbracket \varphi \rrbracket_{j}:=\lim _{x \rightarrow x_{j}+} \varphi(x)-\lim _{x \rightarrow x_{j}-} \varphi(x)$ and $x_{0}:=-\infty, x_{N+1}:=\infty$. It is easy to see that $D(A)=H^{2}(\mathbb{R}, \mathbb{C})$ using the definition of the second weak derivative and the fact that at the interfaces any $\varphi \in D(A)$ is of class $C^{1}$. Note that assumptions (A1)-(A2), (f1)-(f4) are satisfied for any $\omega \in \mathbb{C} \backslash\{0,-\mathrm{i} \gamma\}$.

### 4.1. The linear eigenvalue problem

In order to apply Theorem 2.1, we need to find an eigen-pair $\left(\omega_{0}, \varphi_{0}\right)$ of the linear problem

$$
\begin{equation*}
\varphi^{\prime \prime}+W(x, \omega) \varphi=0, \quad x \in \mathbb{R} \tag{54}
\end{equation*}
$$

coupled with IFCs, with a simple and isolated $\omega_{0} \in \mathbb{C} \backslash\{0\}$. To ensure the realness of the frequency we need, in fact, $\omega_{0} \in \mathbb{R} \backslash\{0\}$ as well as the $\mathcal{P} \mathcal{T}$ symmetry of $W(\cdot, \omega)$, such that Proposition 3.1 can be applied. We shall see that the existence of a simple and isolated eigenvalue $\omega_{0} \in \mathbb{R} \backslash\{0\}$ strongly depends on the choice of the layers. As we show in Sect. 4.1.1, the choice of two layers $(N=1)$ of periodic materials with one being conservative and the other a non-conservative homogeneous material (i.e. with a complex $\hat{\chi}^{(1)}$ ) leads to no real eigenvalues $\omega_{0}$. On the other hand, in Sect. 4.1.2 we find two $\mathcal{P} \mathcal{T}$-symmetric settings with three homogeneous layers $(N=2)$ leading to the existence of an isolated simple eigenvalue $\omega_{0} \in \mathbb{R}$ in (54). These settings are: (active dielectric - conservative Drude metal-lossy dielectric) and a hypothetical setting of (Drude metal with gain-lossless dielectric-lossy Drude metal).
4.1.1. Two periodic layers. We consider first the case of two layers, each being either a periodic metal or a periodic dielectric, where we set the interface at $x=0$. Hence

$$
W(x, \omega)=W_{ \pm}(x, \omega) \quad \text { for } \quad \pm x>0
$$

where the functions $W_{ \pm}(x, \omega)=\omega^{2}\left(1+\hat{\chi}_{ \pm}^{(1)}(x, \omega)\right)-k^{2}$ are periodic in $x$ with periods $\nu_{ \pm}>0$. The governing linear problem (54) has two linearly
independent Bloch wave solutions $\psi_{1,2}^{ \pm}$on the half line $\pm x>0$ respectively. The Bloch wave theory for the Hill's equation (54) can be found in [12]. A necessary condition for the existence of an $L^{2}(\mathbb{R}, \mathbb{C})$-solution of (54) is

$$
0 \notin S:=\sigma\left(-\frac{d^{2}}{d x^{2}}-W_{+}(\cdot, \omega)\right) \cup \sigma\left(-\frac{d^{2}}{d x^{2}}-W_{-}(\cdot, \omega)\right)
$$

Otherwise (if $0 \in S$ ), on at least one of the half lines there is no decaying solution. If $0 \notin S$, the solutions $\psi_{1,2}^{ \pm}$have the form

$$
\psi_{j}^{+}(x)=p_{j}^{+}(x) e^{(-1)^{j} \lambda_{+} x}, \quad \psi_{j}^{-}(x)=p_{j}^{-}(x) e^{(-1)^{j} \lambda_{-} x}, \quad j=1,2
$$

where $\operatorname{Re}\left(\lambda_{ \pm}\right)>0, p_{j}^{ \pm}\left(x+\nu_{ \pm}\right)=p_{j}^{ \pm}(x)$ for all $x \in \mathbb{R}$ and both $j=1,2$. If, say, $W_{+}$is real, then $p_{j}^{+}$and $\lambda_{+}$can be chosen real but $p_{j}^{+}$is generally $2 \nu_{+}$-periodic, see [12].

An $L^{2}(\mathbb{R})$-solution is given by

$$
\varphi(x)= \begin{cases}\varphi^{+}(x):=p_{1}^{+}(x) e^{-\lambda_{+} x}, & x>0 \\ \varphi_{-}(x):=p_{2}^{-}(x) e^{\lambda_{-} x}, & x<0\end{cases}
$$

Due to the linearity of (54) the $C^{1}$-matching condition (IFCs) is equivalent to the condition

$$
\begin{equation*}
R_{+}:=\frac{\varphi_{+}^{\prime}(0)}{\varphi_{+}(0)}=\frac{\varphi_{-}^{\prime}(0)}{\varphi_{-}(0)}=: R_{-} . \tag{55}
\end{equation*}
$$

Note that by varying the parameters $\omega, k \in \mathbb{R}$ we get $R_{ \pm}=R_{ \pm}(\omega, k)$.
In [10] the case of real, periodic $W_{ \pm}$, i.e. the one of two periodic conservative materials, was considered and eigenvalues were found by varying $k$ and searching for zeros of $R_{+}(k)-R_{-}(k)$.

In the presence of non-conservative materials the respective potential $W$ is complex. It is easy to see that the single interface of a conservative material (e.g. a classical dielectric) with a real $W$ and non-conservative homogeneous material (e.g. a Drude metal) with a complex $W$ does not support any eigenvalues of (54). Note that this does not contradict the existence of SPPs at such single metal/dielectric interfaces in general because this existence holds for TM-polarisations, see [19]. Without loss of generality we assume that the conservative material is on the half line $x>0$, i.e. $W_{+}(x, \omega)$ is real. The non-conservative material in $x<0$ is homogeneous (described, e.g. by (50)) such that $W_{-}(x, \omega)=W_{-}(\omega)$ is complex and independent of $x$. Hence $\varphi_{-}(x)=c e^{\lambda_{-} x}$ with $^{3} \lambda_{-}=\sqrt{W_{-}(\omega)} \in \mathbb{C} \backslash \mathbb{R}$, such that

$$
R_{-}=\lambda_{-}
$$

and $\varphi_{+}(x)=p_{1}^{+}(x) e^{-\lambda_{+} x}$ with $\lambda_{+} \in \mathbb{R}$ and $p_{1}^{+}(x)$ real and $2 \nu_{+}-$periodic, such that

$$
R_{+}=\frac{p_{1}^{+^{\prime}}(0)}{p_{1}^{+}(0)}-\lambda_{+} \in \mathbb{R}
$$

[^2]Note that $p_{1}^{+}(0)=0$ implies $c=0$ due to (IFCs) such that only the trivial solution $\varphi \equiv 0$ is produced in that case.

Because $R_{-} \in \mathbb{C} \backslash \mathbb{R}$ and $R_{+} \in \mathbb{R}$, condition (55) is not satisfied and no eigenvalue exists.
4.1.2. Three homogeneous layers. Next we consider three homogeneous material layers with interfaces at $x=0$ and $x=d>0$, i.e. a sandwich geometry with two unbounded layers,

$$
\begin{gather*}
\hat{\chi}^{(1)}(x, \omega):= \begin{cases}\eta_{-} & \text {for } x<0, \\
\eta_{*} & \text { for } x \in(0, d), \\
\eta_{+} & \text {for } x>d,\end{cases} \\
\text { i.e. } W(x, \omega)= \begin{cases}\omega^{2}\left(1+\eta_{-}\right)-k^{2}=: W_{-} & \text {for } x<0 \\
\omega^{2}\left(1+\eta_{*}\right)-k^{2}=: W_{*} & \text { for } x \in(0, d), \\
\omega^{2}\left(1+\eta_{+}\right)-k^{2}=: W_{+} & \text {for } x>d,\end{cases} \tag{56}
\end{gather*}
$$

where $\eta_{ \pm}, \eta_{*} \in \mathbb{C}$. A localized solution of (54) is possible only if

$$
\begin{equation*}
0 \notin S:=\bigcup_{ \pm} \sigma\left(-\frac{d^{2}}{d x^{2}}-W_{ \pm}(\omega)\right) \tag{57}
\end{equation*}
$$

Note that if the semi-infinite layers are conservative, then $\eta_{+}, \eta_{-} \in \mathbb{R}$ and (57) is equivalent to

$$
k^{2}>\omega^{2}\left(1+\max \left\{\eta_{+}, \eta_{-}\right\}\right) .
$$

Under assumption (57) we have

$$
\varphi(x)= \begin{cases}A e^{\lambda_{-} x} & \text { for } x<0  \tag{58}\\ B e^{\mu x}+C e^{-\mu x} & \text { for } x \in(0, d) \\ D e^{-\lambda_{+} x} & \text { for } x>d\end{cases}
$$

where $\mu:=\sqrt{-W_{*}}, \lambda_{ \pm}:=\sqrt{-W_{ \pm}}, \operatorname{Re}\left(\lambda_{ \pm}\right)>0$, and where $A, B, C, D \in \mathbb{C}$ are constants to be determined. We can normalize such that $D=1$. Then the $C^{1}$-matching at $x=0$ and $x=d$ is equivalent to

$$
\left\{\begin{array} { l } 
{ A = B + C }  \tag{59}\\
{ A \lambda _ { - } = \mu B - \mu C }
\end{array} \quad \left\{\begin{array}{l}
e^{-\lambda_{+} d}=B e^{\mu d}+C e^{-\mu d} \\
-\lambda_{+} e^{-\lambda_{+} d}=\mu B e^{\mu d}-\mu C e^{-\mu d}
\end{array}\right.\right.
$$

This system has the unique solution

$$
\begin{aligned}
& A=\frac{e^{-\lambda_{+} d}}{2}\left[\left(1-\frac{\lambda_{+}}{\mu}\right) e^{-\mu d}+\left(1+\frac{\lambda_{+}}{\mu}\right) e^{\mu d}\right] \\
& B=\frac{1}{2}\left(1-\frac{\lambda_{+}}{\mu}\right) e^{-\left(\mu+\lambda_{+}\right) d}, \quad C=\frac{1}{2}\left(1+\frac{\lambda_{+}}{\mu}\right) e^{-\left(-\mu+\lambda_{+}\right) d}
\end{aligned}
$$

together with a condition on the parameter $d$ :

$$
\begin{equation*}
\exists m \in \mathbb{Z}: d=\tilde{d}_{m}:=\frac{1}{2 \mu}\left[\log \left(\frac{\left(\mu-\lambda_{-}\right)\left(\mu-\lambda_{+}\right)}{\left(\mu+\lambda_{-}\right)\left(\mu+\lambda_{+}\right)}\right)+2 \pi \mathrm{i} m\right] \in(0, \infty) \tag{60}
\end{equation*}
$$

The term $2 \pi \mathrm{i} m$ appears due the fact that $z=\log (b)+2 \pi \mathrm{i} m$ solves $e^{z}=b$ for any $m \in \mathbb{Z}$. Condition (60) means that under assumption (57) there is a width $d>$ 0 supporting a real eigenvalue if and only if $\tilde{d}_{m}=\tilde{d}_{m}(\omega)$ is positive for some value of $\omega \in \mathbb{R}$ and some $m \in \mathbb{Z}$. If all the layers are homogeneous conservative materials, (60) cannot be satisfied because $\mu>0$ and the argument of the logarithm on the right hand side of $(60)$ lies in $(0,1)$ such that $\operatorname{Re}\left(\tilde{d}_{m}\right)<0$ for all $m \in \mathbb{Z}$.

Next, we consider three sandwich settings, out of which the last two are $\mathcal{P} \mathcal{T}$-symmetric. As the numerical evaluation of $\tilde{d}_{m}$ suggests, both of these apparently lead to the existence of linear eigenvalues $\omega_{0} \in \mathbb{R}$ and hence to real bifurcating nonlinear eigenvalues $\omega$. The first one of these $\mathcal{P} \mathcal{T}$-symmetric cases has been taken from the physics literature $[2,3]$ while the second one corresponds to a hypothetical material.
Case 1: conservative dielectric-Drude metal-conservative dielectric
Note that a Drude metal layer being sandwiched between two conservative dielectrics does not produce a $\mathcal{P} \mathcal{T}$-symmetric potential $W$. Hence, we do not expect a real eigenvalue $\omega$.

We have $\eta_{ \pm}>0$ and $\eta_{*}=-\frac{1}{\omega^{2}+\mathrm{i} \gamma \omega} \in \mathbb{C}$. Our numerical study of $\tilde{d}_{m}$ shows that for any choice of the constants $\eta_{+}, \eta_{-}, \gamma \in \mathbb{R}$ and $m \in \mathbb{Z}$ we cannot find a frequency $\omega$ for which $\tilde{d}_{m}(\omega) \in(0, \infty)$. Figure 1a, b shows as an example the behaviour of the maps $\omega \mapsto \operatorname{Re}\left(\tilde{d}_{0}(\omega)\right)$ and $\omega \mapsto \operatorname{Im}\left(\tilde{d}_{0}(\omega)\right)$ for different choices of $\gamma, \eta_{+}, \eta_{-}$.
Case 2: $\mathcal{P} \mathcal{T}$-symmetric (dielectric-metal-dielectric) setting
For the case of a homogeneous Drude metal sandwiched between two homogeneous non-conservative dielectric layers, we have

$$
\begin{equation*}
\eta_{ \pm}=\eta_{R}^{ \pm}+\mathrm{i} \eta_{I}^{ \pm} \quad \text { and } \quad \eta_{*}=-\frac{1}{\omega^{2}+\mathrm{i} \gamma \omega} \tag{61}
\end{equation*}
$$

This setting leads to a $\mathcal{P} \mathcal{T}$-symmetric $W$ (with respect to $x=d / 2$, i.e. $W\left(\frac{d}{2}-\right.$ $x, \omega)=\overline{W\left(\frac{d}{2}+x, \omega\right)}$ for all $x \in \mathbb{R}$ and $\left.\omega \in \mathbb{R} \backslash\{0\}\right)$ if we choose

$$
\eta_{R}^{+}=\eta_{R}^{-}, \quad \eta_{I}^{+}=-\eta_{I}^{-}, \quad \gamma=0
$$

Hence, one of the dielectric layers is lossy while the other is active and generates energy gain.

Note that in this example the simple transformation $\omega^{\prime}:=\omega^{2}$ leads to a linear dependence on the spectral parameter $\omega^{\prime}$.

This configuration of a conservative metal sandwiched between a couple of well-prepared active and lossy dielectrics was considered, e.g., in [2,3]. As the active material we consider titanium dioxide $\left(\mathrm{TiO}_{2}\right)$, with refractive index $n=3.2+0.2 \mathrm{i}$, and as the metal we choose silver with the bulk plasma frequency $\omega_{p}(\mathrm{Ag})=8,85 \cdot 10^{15} \mathrm{~s}^{-1}$. We use $\omega_{p}(\mathrm{Ag})$ as the rescaling parameter in (51). Recalling the relation $W(x)=\omega^{2} n(x)^{2}-k^{2}$ between the refractive index and the potential $W$, we choose $k=2$ and obtain $\eta_{ \pm}=9.2 \pm 1.28 \mathrm{i}$.

Our numerical tests show that this $\mathcal{P} \mathcal{T}$-symmetric setting leads to $\tilde{d}(\omega)>$ 0 for all $\omega>\omega_{D M D} \approx 2.23$, see Fig. 1c. For the computation of the bifurcation we choose the point $\omega_{0} \approx 3.8275$, for which $d=\tilde{d}_{-1}(3.8275) \approx 1$.
(a)

(c)

(b)

(d)


Figure 1. The graph of $\omega \mapsto \operatorname{Re}\left(\tilde{d}_{m}(\omega)\right)$ (full blue line) and $\omega \mapsto \operatorname{Im}\left(\tilde{d}_{m}(\omega)\right)$ (dashed red line). a, b Case 1 with $k=1$, $\gamma=\mp 0.5, \eta_{+}=5, \eta_{-}=0.05$, and $m=0$. c Case 2 with $k=2, \eta_{ \pm}=9.2 \pm 1.28 \mathrm{i}, \gamma=0$ and $m=1$. d Case 3 with $k=2, \gamma^{+}=-\gamma^{-}=0.5, \eta_{*}=0.2$ and $m=-1$. In $\mathbf{c}$ (resp. d) the chosen value $d=1$, attained at $\omega \approx 3.8275$ (resp. $d=0.5$, attained at $\omega \approx 2.8096$ ), is highlighted in the graph (colour figure online)

As the graph in Fig. 1c suggests, $\frac{d}{d \omega} \operatorname{Re}\left(\tilde{d}_{-1}(\omega)\right) \neq 0$ for all $\omega>\omega_{D M D}$. Hence, for any width $d \in \tilde{d}_{-1}\left(\left(\omega_{D M D}, \infty\right)\right)$ there is only one $\omega_{0}$ for which $d=\tilde{d}_{-1}\left(\omega_{0}\right)$, i.e. for which (60) holds, i.e. there is only one eigenvalue $\omega_{0}$. Because

$$
\sigma_{\mathrm{ess}}\left(-\frac{d^{2}}{d x^{2}}-W\right)=\bigcup_{ \pm}\left\{\lambda-\omega^{2}\left(1+\eta_{ \pm}\right)+k^{2}: \lambda \in(0, \infty)\right\}
$$

and $\eta_{ \pm} \in \mathbb{C} \backslash \mathbb{R}$, we have $\omega_{0} \notin \sigma_{\text {ess }}\left(-\frac{d^{2}}{d x^{2}}-W\right)$. Moreover, the constants $A, B, C, D$ are unique (up to normalization of $\varphi_{0}$ ). Hence, $\omega_{0}$ is an isolated simple eigenvalue.
Case 3: $\mathcal{P} \mathcal{T}$-symmetric (metal-dielectric-metal) setting.

For the case of a homogeneous conservative dielectric sandwiched between two homogeneous Drude metal layers we have

$$
\begin{equation*}
\eta_{ \pm}=-\frac{1}{\omega^{2}+\mathrm{i} \gamma^{ \pm} \omega} \in \mathbb{C} \quad \text { and } \quad \eta_{*}>0 \tag{62}
\end{equation*}
$$

This setting leads to a $\mathcal{P} \mathcal{T}$-symmetric $W$ with respect to $x=d / 2$ if we choose $\omega \in \mathbb{R} \backslash\{0\}$ and $\gamma^{+}=-\gamma^{-}$. Note that this setting (with $\gamma_{+} \neq 0$ ) is hypothetical as materials with $\chi^{(1)}(\omega)=-1 /\left(\omega^{2}+\mathrm{i} \gamma \omega\right)$ and a negative $\gamma$ may not exist.

Also notice that here, unlike the previous example, the dependence of $W$ on $\omega$ is truly nonlinear. We retrieve numerically again a similar plot for the function $\omega \mapsto \tilde{d}(\omega)$, as shown in Fig. 1d: $\tilde{d}(\omega)>0$ for all $\omega>\omega_{M D M}$ with some $\omega_{M D M}>0$. For $k=2, \gamma^{+}=0.5$ and $\eta_{*}=0.2$ we get $\omega_{M D M} \approx 1.83$. For the computation of the bifurcation we choose the point $\omega_{0} \approx 2.8096$ for which $d=\tilde{d}_{-1}\left(\omega_{0}\right) \approx 0.5$. Analogously to Case 2 , such $\omega_{0}$ is an isolated simple eigenvalue.

### 4.2. Bifurcation of a nonlinear eigenvalue

We are now in a position to apply the bifurcation result of Theorem 2.1 in its symmetric version provided by Proposition 3.1 in both settings given by Case 2 and Case 3 and find a bifurcating branch of solutions to the reduced Maxwell's equation (49). In both cases we choose the cubic susceptibility $\chi^{(3)} \equiv 1$.

The numerical computations below were obtained using a centered finite difference discretization of fourth order on an equispaced grid and with the condition that $\varphi(x)=0$ for $x$ outside the computational interval. The linear eigenfunction $\varphi_{0}$ was computed using Matlab's built in eigenvalue solver eigs. The computation of the nonlinear solution $\varphi$ was done via the standard Newton's iteration.

Case 2. In this case the linear susceptibility $\chi^{(1)}$ is chosen as in (56), (61) with parameters as in Sect. 4.1.2 $\left(k=2, \gamma=0\right.$ and $\left.\eta_{ \pm}=9.2 \pm 1.28 \mathrm{i}\right)$. The linear eigenvalue is selected as $\omega_{0}=3.8275$. The last condition to check is (Wt). A numerical approximation produces $\left\langle\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}, \varphi_{0}^{*}\right\rangle \approx 7.602$.

Figure 2a shows the resulting potential $W\left(\cdot, \omega_{0}\right)$. In Fig. 2d the bifurcation diagram is shown, where the actual (numerically computed) branch is plotted with the dashed blue line for $\omega \in\left(3.2, \omega_{0}\right)$, while the first order approximation $\left(\omega_{0}+\varepsilon \nu,\left\|\varepsilon^{1 / 2} \varphi_{0}\right\|_{L^{2}}\right)$ for $\varepsilon \in\left(0,\left(\omega_{0}-3.2\right) /|\nu|\right)$ is plotted in full red. The numerical value of $\nu$ is $\nu \approx-7.4226$. A good agreement is observed between the asymptotic and the numerical curves in the vicinity of $\omega_{0}$. The eigenfunction $\varphi_{0}$ related to the eigenvalue $\omega_{0}$ is plotted in Fig. 2b. Finally, Fig. 2c shows the solution $\varphi$ at $\omega=3.2$, i.e. at the last $\omega$ in the continuation procedure.

Case 3. Here we choose $\chi^{(1)}$ as in (56), (62) with parameters as in Sect. 4.1.2 $\left(k=2, \gamma^{+}=-\gamma^{-}=0.5\right.$, and $\left.\eta_{*}=0.2\right)$. The linear eigenvalue is selected as $\omega_{0}=2.8096$. Again we have to check condition (Wt): a numerical approximation produces $\left\langle\partial_{\omega} W\left(\cdot, \omega_{0}\right) \varphi_{0}, \varphi_{0}^{*}\right\rangle \approx 6.202$.


Figure 2. Equations (49), (56), (61) with the parameters: $d=1, \eta_{ \pm}=9.2 \pm 1.28 \mathrm{i}$, and $k=2$. a The potential $W\left(\cdot, \omega_{0}\right)$ at $\omega_{0} \approx 3.8275$. b The eigenfunction $\varphi_{0}$ of (54) at $\omega_{0}$. c The solution of (49) at $\omega=3.2$. d The bifurcation diagram $(\omega,\|\varphi\|)$ (dashed blue) and the approximation $\left(\omega_{0}+\varepsilon \nu,\left\|\varepsilon^{1 / 2} \varphi_{0}\right\|\right)$ (full red) starting at $\omega_{0}$ (colour figure online)

The resulting bifurcation diagram is shown in Fig. 3d, where again the actual branch is plotted with the dashed blue line, while the first order approximation $\left(\omega_{0}+\varepsilon \nu,\left\|\varepsilon^{1 / 2} \varphi_{0}\right\|\right)$ for $\varepsilon \in\left(0,\left(\omega_{0}-1.5\right) /|\nu|\right)$ is plotted in full red. We get numerically $\nu \approx-2.7233$. Once again, a good agreement is observed between the asymptotic and the numerical curves. The eigenfunction $\varphi_{0}$ related to the eigenvalue $\omega_{0}$ is plotted in Fig. 3b. Figure 3c shows the solution $\varphi$ at $\omega=1.5$, i.e. at the last $\omega$ in the continuation procedure.


Figure 3. Equations (49), (56), (62) with the parameters: $d=0.5, \gamma^{+}=-\gamma^{-}=0.5$, and $\eta_{*}=0.2$. a The potential $W\left(\cdot, \omega_{0}\right)$ at $\omega_{0} \approx 2.8096$. b The eigenfunction $\varphi_{0}$ of (54) at $\omega_{0}$. c The solution of (49) for $\omega=1.5$. d The bifurcation diagram $(\omega,\|\varphi\|)$ (dashed blue) and the approximation $\left(\omega_{0}+\right.$ $\varepsilon \nu,\left\|\varepsilon^{1 / 2} \varphi_{0}\right\|$ ) (full red) starting at $\omega_{0}$ (colour figure online)

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Tomáš Dohnal and Giulio Romani<br>Institut für Mathematik<br>Martin-Luther-Universität Halle-Wittenberg<br>06099 Halle (Saale)<br>Germany<br>e-mail: tomas.dohnal@mathematik.uni-halle.de<br>Giulio Romani<br>e-mail: giulio.romani@mathematik.uni-halle.de

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[^0]:    ${ }^{1}$ Recall that if 0 is a simple isolated eigenvalue of $L\left(\cdot, \omega_{0}\right)$, then it is also a simple isolated eigenvalue of $L\left(\cdot, \omega_{0}\right)^{*}$, cf. [16, Chap.III.6.5-6]

[^1]:    ${ }^{2}$ With a little abuse of notation, henceforth we write $G(\psi):=G(\sigma, \psi)$ and similarly for $R$ when $\sigma$ is assumed to be fixed.

[^2]:    ${ }^{3}$ Henceforth, for any $z \in \mathbb{C}$ we choose the square root $\sqrt{z}$ as the one solution $a$ of $a^{2}=z$ with $\arg (a) \in(-\pi / 2, \pi / 2]$.

