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# Interaction of modulated pulses in the nonlinear Schrödinger equation with periodic potential 

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#### Abstract

We consider a cubic nonlinear Schrödinger equation with periodic potential. In a semiclassical scaling the nonlinear interaction of modulated pulses concentrated in one or several Bloch bands is studied. The notion of closed mode systems is introduced which allows for the rigorous derivation of a finite system of amplitude equations describing the macroscopic dynamics of these pulses.


## 1 Introduction and main result

In this work we study the asymptotic behavior for $0<\varepsilon \ll 1$ of the following nonlinear Schrödinger equation (NLS)

$$
\begin{equation*}
\mathrm{i} \varepsilon \partial_{t} u^{\varepsilon}=-\frac{\varepsilon^{2}}{2} \Delta u^{\varepsilon}+V_{\Gamma}\left(\frac{x}{\varepsilon}\right) u^{\varepsilon}+\varepsilon \kappa\left|u^{\varepsilon}\right|^{2} u^{\varepsilon}, \quad x \in \mathbb{R}^{d}, t \in \mathbb{R} \tag{1.1}
\end{equation*}
$$

governing the dynamics of a wave field $u^{\varepsilon}(t, \cdot) \in L^{2}\left(\mathbb{R}^{d}\right)$. Here, $\kappa \in \mathbb{R}$ and the potential $V_{\Gamma}=V_{\Gamma}(y) \in \mathbb{R}$ is assumed to be smooth and periodic with respect to some regular lattice $\Gamma \simeq \mathbb{Z}^{d}$, generated by a given basis $\left\{\zeta_{1}, \ldots, \zeta_{d}\right\}, \zeta_{l} \in \mathbb{R}^{d}$, i.e.

$$
\begin{equation*}
V_{\Gamma}(y+\gamma)=V_{\Gamma}(y), \quad \forall y \in \mathbb{R}^{d}, \gamma \in \Gamma \equiv\left\{\gamma=\sum_{l=1}^{d} \gamma_{l} \zeta_{l} \in \mathbb{R}^{d}: \gamma_{l} \in \mathbb{Z}\right\} \tag{1.2}
\end{equation*}
$$

It is well known that if $\kappa<0$ the solution of (1.1) in general does not exist for all times, i.e. finite-time blow-ups may occur, cf. [25].
The equation (1.1) can be seen as a simplified model of the one considered in [7]. There the main motivation was to study, from a semiclassical point of view, the dynamics of a Bose-Einstein condensate in an optical lattice, described by $V_{\Gamma}$, cf. [8, $9,19]$ for more details. To this end a rescaling of the appearing physical parameters yields an equation similarly to (1.1), but with an additional non-periodic confining potential, which we shall neglect in the following. The parameter $\varepsilon \ll 1$ then describes the microscopic/macroscopic scale ratio. The main assumption for the analysis presented in [7] has been that the initial data $u^{\varepsilon}(0, \cdot)$ is supposed to be of WKB type and in particular it has to be concentrated in a single (isolated) Bloch band $E_{\ell}(k) \in \mathbb{R}$. These energy bands describe the spectral subspaces corresponding to the periodic Hamiltonian operator

$$
\begin{equation*}
H_{\mathrm{per}}^{\varepsilon}:=-\frac{\varepsilon^{2}}{2} \Delta+V_{\Gamma}\left(\frac{x}{\varepsilon}\right), \tag{1.3}
\end{equation*}
$$

cf. Section 2.1 below for more details. In the linear case similar WKB approximations have been established earlier in [3, 14], yielding an approximate macroscopic description (i.e. on time- and length-scales of order one) of the highly oscillatory solution to (1.1). However, the question concerning a generalization of the results in [7], in particular to the case of multiple bands, has been open so far. Here we will answer this question for initial data which correspond to a sum of modulated plane waves.

In order to derive an approximate macroscopic description we shall proceed by a two scale expansion method similar to that in [7]. To this end a detailed understanding of the influence of the nonlinearity is crucial. Indeed we will show that the solution to (1.1) can be approximated (in a suitably scaled Sobolev space) via

$$
\begin{equation*}
u^{\varepsilon}(t, x) \sim \sum_{m=1}^{M} a_{m}(t, x) \chi_{\ell_{m}}\left(\frac{x}{\varepsilon} ; k_{m}\right) \mathrm{e}^{\left.\mathrm{i}\left(k_{m} \cdot x-t E_{\ell_{m}}\left(k_{m}\right)\right) / \varepsilon\right)}+\mathcal{O}(\varepsilon), \tag{1.4}
\end{equation*}
$$

for $M \in \mathbb{N}$, where the set $\left\{\left(k_{m}, \ell_{m}\right): m=1, \ldots, M\right\}$ is assumed to form a closed mode system, see Definition 2.2,t of sufficient high order $\Lambda$. As we shall see $\Lambda$ will depend on the spatial dimension $d$. The amplitudes $a_{m}$ are the (local-in-time) solutions to the nonlinear system of amplitude equations

$$
\begin{equation*}
\mathrm{i} \partial_{t} a_{m}+\mathrm{i} \vartheta_{m} \cdot \nabla_{x} a_{m}=\sum_{\substack{p, q, r=1 ; \\ \Sigma\left(\mu_{p}, \mu_{q}, \mu_{r}\right)=\Sigma\left(\mu_{m}\right)}}^{M} \kappa_{(p, q, r, m)} a_{p} \bar{a}_{q} a_{r}, \quad m=1, \ldots, M \tag{1.5}
\end{equation*}
$$

The above system describes a so-called four-wave interaction, also known (most prominently in laser physics and nonlinear optics) as four-wave mixing, cf. [1, 6, 15]. By (1.4) we allow for nonlinear interactions within the same band but also consider interactions of different bands. In particular, energy or mass transfer between different bands is expected due to the presence of the nonlinearity on the right hand side of (1.5). To our knowledge this phenomenon has never been studied rigorously in the context of Schrödinger type equations.
One should note that the concept of wave mixing is strongly linked to plane waves as considered above. Indeed, if one allows for more general phases (which is possible in the case of a single pulse [7]), a rigorous understanding, even in much simpler cases, is lacking and in particular one can not expect the nonlinear interaction to be maintained on macroscopic time-scales in general. On the other hand we could, without any problems, allow for (smooth) higher order nonlinearities as considered in [7]. This however would result into a much more involved resonance-structure for the nonlinear interactions and in order to keep our presentation simple we restrict ourselves to the cubic case.
Before going into more details, let us briefly mention the following mathematically rigorous works which, besides [7], are most closely related to ours: In [24] the same equation as (1.1) is considered but in a slightly different scaling. Similarly, a nonlinear Schrödinger type model is derived in [5] from a semilinear wave equation with periodic coefficients and in [13] from an underlying oscillator chain model. Concerning nonlinear wave interactions, there exist several results (mostly three-wave mixing) in the context of strictly hyperbolic systems, see, e.g., [17, 18, 20, 22], and in the case of microscopically discrete dynamical systems, cf. [10], [11] and the references given therein.
The paper is organized as follows: In Section 2 we introduce the basic notions needed in the following, i.e. Bloch bands, mode systems and the closure condition. We also discuss several illustrative examples of wave mixing. Then in Section 3 the formal
derivation of the approximate solution is given in detail. In Section 4 the obtained formal asymptotic description is rigorously established and our main theorem is stated for a closed mode system of order $\Lambda=2 N+1$, with $N>d / 2$. In Section 5 it is shown that this condition can be relaxed though by introducing the concept of weak closure. Finally, in Appendix A, we discuss in more detail the underlying Hamiltonian structure of the amplitude equations (1.5).

## 2 Mode systems and resonances

For our work it is essential to study the spectral properties of

$$
H_{\mathrm{per}}=-\frac{1}{2} \Delta_{y}+V_{\Gamma}(y)
$$

since they will basically determine the fast degrees of freedom in our model.

### 2.1 Bloch's spectral problem

In what follows we denote by $Y$ the centered fundamental domain of the lattice $\Gamma$, i.e.

$$
\begin{equation*}
Y:=\left\{\gamma \in \mathbb{R}^{d}: \gamma=\sum_{l=1}^{d} \gamma_{l} \zeta_{l}, \gamma_{l} \in\left[-\frac{1}{2}, \frac{1}{2}\right]\right\} . \tag{2.1}
\end{equation*}
$$

By $\mathcal{Y}$ we denote the $d$-dimensional torus $\mathbb{R}_{/ \Gamma}^{d}$, which is obtained also from $Y$ by identifying opposite faces. Note that writing $H^{s}(\mathcal{Y})$ then automatically includes periodicity conditions. Moreover, $Y^{*}$ denotes the corresponding basic cell of the dual lattice $\Gamma^{*}$. Equipped with periodic boundary conditions $\mathcal{Y}^{*}$ is usually called the Brillouin zone and hence we shall denote it by $\mathcal{B} \equiv \mathcal{Y}^{*}$.
Next, consider the so-called Bloch eigenvalue problem [4], i.e. the following spectral problem on $\mathcal{Y}$ :

$$
\begin{equation*}
H_{\Gamma}(k) \chi_{\ell}(y ; k)=E_{\ell}(k) \chi_{\ell}(y ; k), \quad k \in \mathcal{B}, \ell \in \mathbb{N}, \tag{2.2}
\end{equation*}
$$

where $E_{\ell}(k) \in \mathbb{R}$ and $\chi_{\ell}(y) \equiv \chi_{\ell}(y ; k)$ denote the $\ell$-th eigenvalue and eigenvector of the shifted Hamiltonian operator

$$
H_{\Gamma}(k):=\mathrm{e}^{-\mathrm{i} k \cdot y} H_{\mathrm{per}} \mathrm{e}^{\mathrm{i} k \cdot y}=\frac{1}{2}\left(-\mathrm{i} \nabla_{y}+k\right)^{2}+V_{\Gamma}(y) .
$$

Let us recall some well known facts for this eigenvalue problem, cf. [2, 3, 23, 26]. Since $V_{\Gamma}$ is smooth and periodic, we get that, for every fixed $k \in \mathcal{B}, H_{\Gamma}(k)$ is selfadjoint on $L^{2}(\mathcal{Y})$ with domain $H^{2}(\mathcal{Y})$ and compact resolvent. Hence the spectrum of $H_{\Gamma}(k)$ is given by

$$
\operatorname{spec}\left(H_{\Gamma}(k)\right)=\left\{E_{\ell}(k) ; \ell \in \mathbb{N}\right\} \subset \mathbb{R}
$$

One can order the eigenvalues $E_{\ell}(k)$ according to their magnitude and multiplicity such that

$$
E_{1}(k) \leq \ldots \leq E_{\ell}(k) \leq E_{\ell+1}(k) \leq \ldots
$$

Moreover every $E_{\ell}(k)$ is periodic w.r.t. $\Gamma^{*}$ and it holds that $E_{\ell}(k)=E_{\ell}(-k)$. The set $\left\{E_{\ell}(k) ; k \in \mathcal{B}\right\}$ is called the $\ell$ th energy band. The associated eigenfunctions, the Bloch functions, $\chi_{\ell}(y ; k)$ form (for every fixed $k \in \mathcal{B}$ ) an orthonormal basis in $L^{2}(\mathcal{Y})$. We choose the usual normalization such that

$$
\left\langle\chi_{\ell_{1}}(k), \chi_{\ell_{2}}(k)\right\rangle_{L^{2}(\mathcal{Y})} \equiv \int_{\mathcal{Y}} \overline{\chi_{\ell_{1}}}(y ; k) \chi_{\ell_{2}}(y ; k) \mathrm{d} y=\delta_{\ell_{1}, \ell_{2}}, \quad \ell_{1}, \ell_{2} \in \mathbb{N} .
$$

Concerning the dependence on $k \in \mathcal{B}$, it has been shown, cf. [23], that for any $\ell \in \mathbb{N}$ there exists a closed subset $\mathcal{U} \subset \mathcal{B}$ such that $E_{\ell}(k)$ is analytic in $\mathcal{O}:=\mathcal{B} \backslash \mathcal{U}$. Similarly, the eigenfunctions $\chi_{\ell}$ are found to be analytic and periodic in $k$, for all $k \in \mathcal{O}$. Moreover it holds that

$$
E_{\ell-1}(k)<E_{\ell}(k)<E_{\ell+1}(k), \quad \forall k \in \mathcal{O} .
$$

If this condition is satisfied for all $k \in \mathcal{B}$ then $E_{\ell}(k)$ is said to be an isolated Bloch band. Finally we remark that

$$
\operatorname{meas} \mathcal{U}=\operatorname{meas}\left\{k \in \mathcal{B} \mid E_{\ell_{1}}(k)=E_{\ell_{2}}(k), \ell_{1} \neq \ell_{2}\right\}=0
$$

In this set of measure zero one encounters so-called band crossings. The elements of this set are characterized by the fact that $E_{\ell}(k)$ is only Lipschitz continuous and hence the group velocity $\vartheta:=\nabla_{k} E_{\ell}(k)$ does not exist.

### 2.2 Resonances and closed mode systems

Our goal is to derive an approximate description of our model for $\varepsilon \ll 1$. To this end we shall first introduce several definitions needed to do so.

Definition 2.1. For $k \in \mathcal{B}$ and $\ell \in \mathbb{N}$ we call $\mu=(k, \ell)$ a mode and $\mathcal{M}:=\mathcal{B} \times \mathbb{N}$ the set of all modes.
(i) The graph of all modes is given by

$$
\mathcal{G}=\left\{\left(k, E_{\ell}(k)\right):(k, \ell) \in \mathcal{M}\right\} \subset \mathcal{B} \times \mathbb{R}
$$

(ii) Given a finite set $\mathcal{S}=\left\{\mu_{1}, \ldots, \mu_{M}: M \in \mathbb{N}\right\} \subset \mathcal{M}$ of modes, we call $\mathcal{S}^{\Lambda}$ the (ordered) mode system of size $\Lambda \in \mathbb{N}$ generated by $\mathcal{S}$.
(iii) We further introduce $\Sigma: \mathcal{S}^{\Lambda} \rightarrow \mathcal{B} \times \mathbb{R}$ by

$$
\Sigma\left(\mu_{m_{1}}, \ldots, \mu_{m_{\Lambda}}\right):=\left(\sum_{\lambda=1}^{\Lambda}(-1)^{\lambda+1} k_{m_{\lambda}}, \sum_{\lambda=1}^{\Lambda}(-1)^{\lambda+1} E_{\ell_{m_{\lambda}}}\left(k_{m_{\lambda}}\right)\right)
$$

where $\sum^{*}$ denotes summation modulo $\Gamma^{*}$, and we write $\mathcal{G}_{\mathcal{S}}^{(\Lambda)}:=\Sigma\left(\mathcal{S}^{\Lambda}\right)$ for the corresponding graphs.

In the following the mapping $\Sigma$ will describe the possible nonlinear interaction of modes (in every order of $\varepsilon$ ). Note that $\mathcal{G}_{\mathcal{S}}^{(1)} \subset \mathcal{G}$ and moreover, for any $\Lambda \in \mathbb{N}$, it holds $\mathcal{G}_{\mathcal{S}}^{(\Lambda)} \subset \mathcal{G}_{\mathcal{S}}^{(\Lambda+2)}$. Since we are dealing with a cubic nonlinearity, we shall see that indeed $\Lambda \in \mathbb{N}$ takes only odd values.

Definition 2.2. Given a finite set of modes $\mathcal{S}$ and a subset $\mathcal{T}$ of $\mathcal{M}$.
(i) An element $\left(\mu_{m_{1}}, \ldots, \mu_{m_{\Lambda}}\right) \in \mathcal{S}^{\Lambda}$ is called resonant of order $\Lambda$ to $\mathcal{T}$, if

$$
\Sigma\left(\mu_{m_{1}}, \ldots, \mu_{m_{\Lambda}}\right) \in \mathcal{G}_{\mathcal{T}}^{(1)}
$$

(ii) We say that $\mathcal{S}$ is closed of order $\Lambda$, if the group velocity $\nabla_{k} E_{\ell}(k)$ exists for all $\mu=(k, \ell) \in \mathcal{S}$ and

$$
\begin{equation*}
\mathcal{G}_{\mathcal{S}}^{(\Lambda)} \cap\left(\mathcal{G} \backslash \mathcal{G}_{\mathcal{S}}^{(1)}\right)=\emptyset \tag{2.3}
\end{equation*}
$$

We infer from the above definition that a single mode $\mu$ is always resonant of order $\Lambda=1$ to itself. Note that for any $\Lambda \geq 2$ however, we obtain

$$
\Sigma:\left(\mu_{m_{1}}, \ldots, \mu_{m_{\Lambda}}\right) \mapsto \sigma:=(k, E) \in \mathcal{B} \times \mathbb{R},
$$

which does not necessarily imply $\sigma \in \mathcal{G}_{\mathcal{S}}^{(1)}$. As for the condition (2.3), it is equivalent to saying that for all $\left(\mu_{m_{1}}, \ldots, \mu_{m_{\Lambda}}\right) \in \mathcal{S}^{\Lambda}$ it holds: Either $\Sigma\left(\mu_{m_{1}}, \ldots, \mu_{m_{\Lambda}}\right) \in(\mathcal{B} \times$ $\mathbb{R}) \backslash \mathcal{G}$, or else there exists a $\mu \in \mathcal{S}$ such that $\Sigma\left(\mu_{m_{1}}, \ldots, \mu_{m_{\Lambda}}\right)=\Sigma(\mu) \in \mathcal{G}_{\mathcal{S}}^{(1)}$. The latter obviously means that $\left(\mu_{m_{1}}, \ldots, \mu_{m_{\Lambda}}\right)$ is resonant of order $\Lambda$ to $\mathcal{S}$. We illustrate these concepts in the examples in Section 2.4.
The condition that $\nabla_{k} E_{\ell}(k)$ has to exist implies that we can not deal with mode systems including band crossings, i.e. which include a $\mu=(k, \ell)$ such that $k \in \mathcal{U}$. It is known that the higher the dimension $d$ and the higher the band index $\ell$, the more likely one encounters such crossings. Thus, in terms of practical use one can expect our analysis to be restricted to cases where only a few bands with low energies are taken into account.

Remark 2.3. Also note that condition (2.3) is not equivalent to $\mathcal{G}_{\mathcal{S}}^{(\Lambda)} \cap \mathcal{G} \subset \mathcal{G}_{\mathcal{S}}^{(1)}$, since due to multiple eigenvalues $E_{\ell}(k)$, we may have modes $\mu \in \mathcal{S}$ and $\tilde{\mu} \notin \mathcal{S}$ with $\Sigma(\mu)=\Sigma(\tilde{\mu})$.

### 2.3 The system of amplitude equations

In Section 3 we shall derive a system of amplitude equations describing the macroscopic dynamics of our nonlinear wave interactions. Before going into the details of the derivation let us first discuss the general structure of the obtained system (see also Appendix A). In order to keep our presentation simple, we shall from now on consider only the case of non-degenerated eigenvalues $E_{\ell}(k)$.
Let $\mathcal{S}$ be a given finite set of modes $\left\{\mu_{1}, \ldots, \mu_{M}\right\}$ which is closed of sufficient high order $\Lambda$. Then the equations obtained for the modulating amplitudes $a_{m}(t, x) \in \mathbb{C}$,
for $m \in\{1, \ldots, M\}$, are as in (1.5), i.e.

$$
\mathrm{i} \partial_{t} a_{m}+\mathrm{i} \vartheta_{m} \cdot \nabla_{x} a_{m}=\sum_{\substack{p, q, r=1: \\ \Sigma\left(\mu_{p}, q_{q}, \mu_{r}\right)=\Sigma\left(\mu_{m}\right)}}^{M} \kappa_{(p, q, q, r, m)} a_{p} \bar{a}_{q} a_{r} .
$$

Here $\vartheta_{m} \in \mathbb{R}^{d}$ is the group velocity corresponding to a given mode $\mu_{m}=\left(k_{m}, \ell_{m}\right) \in$ $\mathcal{S}$, i.e. $\vartheta_{m}:=\nabla_{k} E_{\ell_{m}}\left(k_{m}\right)$, and we further denote by

$$
\kappa_{(p, q, r, m)}:=\kappa \int_{\mathcal{Y}} \chi_{\ell_{p}}(y) \bar{\chi}_{\ell_{q}}(y) \chi_{\ell_{r}}(y) \bar{\chi}_{\ell_{m}}(y) \mathrm{d} y
$$

the effective coupling constant $\kappa_{(p, q, r, m)}$, which in general is complex valued.
As in (1.1), the nonlinearity in the amplitude equations is again cubic. It takes into account the sum over all $p, q, r=1, \ldots, M$, for which the resonance condition of order $\Lambda=3$ holds. Explicitly, this is equivalent to the following two conditions

$$
\begin{equation*}
k_{p}-k_{q}+k_{r}=k_{m}, \quad E_{\ell_{p}}\left(k_{p}\right)-E_{\ell_{q}}\left(k_{q}\right)+E_{\ell_{r}}\left(k_{r}\right)=E_{\ell_{m}}\left(k_{m}\right), \tag{2.4}
\end{equation*}
$$

where the first equation has to be understood as a summation modulo $\Gamma^{*}$. The equations (2.4) describe a so-called four-wave interaction, i.e. three modes $\mu_{p}, \mu_{q}, \mu_{r}$ being in resonance with a fourth one $\mu_{m} \in \mathcal{S}$. Clearly, the nonlinear structure on the right-hand side of (1.5) is such that energy transfer between different modes can occur. In other words, even if initially some of the amplitudes $a_{m}(0, \cdot)$ are zero, they will not remain so in general for times $|t| \neq 0$ (see also Subsection 2.4.1 below).
Concerning existence and uniqueness of solutions to the above given amplitude equations, we only need a local-in-time result.

Lemma 2.4. For any initial data $\left(a_{1}(0, \cdot), \ldots, a_{M}(0, \cdot)\right) \in H^{S}\left(\mathbb{R}^{d}\right)^{M}$, with $S>d / 2$, the system (1.5) admits a unique solution

$$
\left(a_{1}, \ldots, a_{M}\right) \in C^{0}\left([0, T) ; H^{S}\left(\mathbb{R}^{d}\right)\right)^{M} \cap C^{1}\left((0, T) ; H^{S-1}\left(\mathbb{R}^{d}\right)\right)^{M}
$$

up to some (finite) time $T>0$.
Proof. Since the left hand side of (1.5) only generates translations by a constant velocity $\vartheta_{m} \in \mathbb{R}^{d}$ and thus conserves the $L^{2}\left(\mathbb{R}^{d}\right)$ norm of each $a_{m}(t, x)$, the assertion of the lemma follows by a standard fixed point argument.

Remark 2.5. Discarding for a moment the nonlinear term in (1.5) we want to point out that the remaining linear transport part is much simpler than in the case of the full WKB type approximation, as discussed in [7] (for a single mode only). In particular we do not run into any problems due to caustics, since the phase functions

$$
\varphi_{m}(t, x)=k_{m} \cdot x-t E_{\ell_{m}}\left(k_{m}\right)
$$

are globally defined. Also note that in the present work the so-called Berry term vanishes, since we do not take into account additional non-periodic potentials, cf. [7, 21].

### 2.4 Examples

In order to illustrate the abstract concepts defined in the former subsections we shall in the following consider several particular examples of nonlinear wave interactions appearing in our model.

### 2.4.1 Four-wave interaction for three pulses

In this example we shall restrict ourselves to a set of modes $\left\{\mu_{m}=\left(k_{m}, \ell_{m}\right): m=\right.$ $1,2,3\}$ which is closed of order $\Lambda=3$, at least, and which satisfies the following resonance conditions

$$
\begin{equation*}
k_{1}-k_{2}+k_{1}=k_{3}, \quad E_{\ell_{1}}\left(k_{1}\right)-E_{\ell_{2}}\left(k_{2}\right)+E_{\ell_{1}}\left(k_{1}\right)=E_{\ell_{3}}\left(k_{3}\right) . \tag{2.5}
\end{equation*}
$$

(Again the first summation has to be understood modulo $\Gamma^{*}$.) This yields the following set of amplitude equations

$$
\left\{\begin{array}{l}
\mathrm{i} \partial_{t} a_{1}+\mathrm{i} \vartheta_{1} \cdot \nabla_{x} a_{1}=W_{1}(a) a_{1}+2 \bar{\kappa}_{(1,2,1,3)} \bar{a}_{1} a_{2} a_{3},  \tag{2.6}\\
\mathrm{i} \partial_{t} a_{2}+\mathrm{i} \vartheta_{2} \cdot \nabla_{x} a_{2}=W_{2}(a) a_{2}+\kappa_{(1,2,1,3)} a_{1}^{2} \bar{a}_{3}, \\
\mathrm{i} \partial_{t} a_{3}+\mathrm{i} \vartheta_{3} \cdot \nabla_{x} a_{3}=W_{3}(a) a_{3}+\kappa_{(1,2,1,3)} a_{1}^{2} \bar{a}_{2},
\end{array}\right.
$$

where we shortly denote

$$
W_{m}(a):=\kappa_{(m, m, m, m)}\left|a_{m}\right|^{2}+2 \sum_{\substack{j=1,2,3: \\ j \neq m}} \kappa_{(m, j, j, m)}\left|a_{j}\right|^{2} \in \mathbb{R} .
$$

We consequently expect the solution of (1.1) to be asymptotically described by

$$
u^{\varepsilon}(t, x) \sim \sum_{m=1}^{3} a_{m}(t, x) \chi_{\ell_{m}}\left(\frac{x}{\varepsilon} ; k_{m}\right) \mathrm{e}^{\mathrm{i}\left(k_{m} \cdot x-E_{\ell_{m}}\left(k_{m}\right) t\right) / \varepsilon}+\mathcal{O}(\varepsilon), \quad \ell_{m} \in \mathbb{N}
$$

It is easily seen that the nonlinearities $W_{m}(a) a_{m}$ can not transfer energy from one band to the other as they are homogeneous in the respective $a_{m}(t, x)$ and that all of the appearing $\kappa_{(m, j, j, m)}, m, j=1,2,3$ are indeed real valued. What is more important though is the fact that the above given amplitude system (2.6) includes an invariant sub-system. Namely, if initially $a_{1}(0, \cdot)=0$, it remains so for all $t \in \mathbb{R}$ and thus the above given system simplifies to

$$
\left\{\begin{array}{l}
\mathrm{i} \partial_{t} a_{2}+\mathrm{i} \vartheta_{2} \cdot \nabla_{x} a_{2}=\left(\kappa_{(2,2,2)}\left|a_{2}\right|^{2}+2 \kappa_{(2,3,3,2)}\left|a_{3}\right|^{2}\right) a_{2}, \\
\mathrm{i} \partial_{t} a_{3}+\mathrm{i} \vartheta_{3} \cdot \nabla_{x} a_{3}=\left(\kappa_{(3,3,3,3)}\left|a_{3}\right|^{2}+2 \kappa_{(2,3,3,2)}\left|a_{2}\right|^{2}\right) a_{3} .
\end{array}\right.
$$

However, such a decoupling does not exist if initially $a_{2}(0, \cdot)=0$, since this amplitude will be generated during the course of time by the remaining two others. Analogously $a_{3}$ is generated by interaction of $a_{1}$ and $a_{2}$.
More formally we can also infer this fact from the closure condition, since any possible combination of $k_{2}$ and $k_{3}$ via the mapping $\Sigma$ (with $\Lambda=3$ ) yields either $k_{2}, k_{3}$, or any other value (like $2 k_{2}-k_{3}$ ) which is not in $\mathcal{G}_{\mathcal{S}}^{(1)}$ by assumption (recall that the system of modes $\mu_{1}, \mu_{2}, \mu_{3}$ is assumed to be closed of order $\Lambda=3$ ). However, if one aims to follow the same argument for, e.g., $k_{1}$ and $k_{2}$, the first equation in (2.5) shows that this sub-system is not closed.

### 2.4.2 The case of several pulses within one Bloch band

This is a particular situation where we keep the band index $\ell \in \mathbb{N}$ fixed and thus only consider the interaction of several pulses within a single Bloch band. Hence the solution to (1.1) takes the asymptotic form

$$
u^{\varepsilon}(t, x) \sim \sum_{m=1}^{M} a_{m}(t, x) \chi_{\ell}\left(\frac{x}{\varepsilon} ; k_{m}\right) \mathrm{e}^{\mathrm{i}\left(k_{m} \cdot x-E_{\ell}\left(k_{m}\right) t\right) / \varepsilon}+\mathcal{O}(\varepsilon), \quad \ell \in \mathbb{N} .
$$

In particular we can expect this description to be correct in cases where the Bloch band $E_{\ell}(k)$ admits only a moderate variation $\Delta_{\ell}:=\operatorname{var}_{k \in \mathcal{B}} E_{\ell}(k)$ and is well separated from the rest of the spectrum of $H_{\Gamma}(k)$ by a sufficiently large gap, i.e.

$$
\min \left\{\left|E_{\ell}(k)-E_{n}(k)\right|: n \in \mathbb{N}, n \neq \ell\right\} \gg \Delta_{\ell}>0
$$

It is easy to show then that a four-wave interaction can always be realized within such a band. Choose $k_{\max }$ and $k_{\min }$ such that $E_{\ell}\left(k_{\min }\right) \leq E_{\ell}(k) \leq E_{\ell}\left(k_{\max }\right)$ for all $k \in \mathcal{B}$. We are looking for a triple $\left(k_{1}, k_{2}, k_{3}\right)$ satisfying (2.5) with $\ell_{j}=\ell, j=1,2,3$. To this end we note that $k_{3}=2 k_{1}-k_{2}$, modulo $\Gamma^{*}$, and define

$$
e\left(k_{1}, k_{2}\right) \equiv 2 E_{\ell}\left(k_{1}\right)-E_{\ell}\left(k_{2}\right)-E_{\ell}\left(2 k_{1}-k_{2}\right) .
$$

Then we have $e\left(k_{\max }, k_{\min }\right)>0>e\left(k_{\min }, k_{\max }\right)$ and by a simple application of the intermediate-value theorem, we easily find $k_{1}, k_{2}$ with $k_{1} \neq k_{2}$ and $e\left(k_{1}, k_{2}\right)=0$.

### 2.4.3 The case of a single pulse decomposed into several bands

This is a second particular case, where we expect an asymptotic description for solutions to (1.1) given by

$$
u^{\varepsilon}(t, x) \sim \sum_{\ell=1}^{L} a_{\ell}(t, x) \chi_{\ell}\left(\frac{x}{\varepsilon} ; k_{0}\right) \mathrm{e}^{\mathrm{i}\left(k_{0} \cdot x-E_{\ell}\left(k_{0}\right) t\right) / \varepsilon}+\mathcal{O}(\varepsilon), \quad k_{0} \in \mathbb{R},
$$

where $k_{0}$ is some given wave vector. One should have the following intuition: Given an initial plane wave of the form

$$
u_{\mathrm{in}}^{\varepsilon}(x)=f(x) \mathrm{e}^{\mathrm{i} k_{0} \cdot x / \varepsilon},
$$

one decomposes the slowly varying macroscopic amplitude $f(x)$ into a sum of terms, each of which is concentrated on a single Bloch band. This is possible due to Bloch's theorem, which ensures that $L^{2}\left(\mathbb{R}^{d}\right)=\bigoplus_{\ell=1}^{\infty} \mathcal{H}_{\ell}$, where $\mathcal{H}_{\ell}$ are the so-called band spaces. Strictly speaking though, one would require countably many terms in the decomposition, which we can take into account here. However for any practical purposes (and if $f(x)$ is sufficiently smooth and rapidly decaying) only the first few Bloch bands need to be taken into account as all higher bands give negligible contributions, cf. [16]. We shall not go into more details here on the precise definition
of $\mathcal{H}_{\ell}$ etc. but refer to [2, 3, 23, 26] for more details on these classical results (see also [16] for a numerical approach).
Concerning the possible generation of resonant modes, it is clear that in this case the first condition in (2.4) is trivially fulfilled, since we are only dealing with a single wave vector $k_{0}$. Thus, it is merely a question on the precise structure of the bands $E_{\ell}\left(k_{0}\right)$ whether one can indeed expect resonances.

## 3 Formal derivation of the approximate solution

In the following we consider a finite set $\mathcal{S}$ of modes which is closed of order $\Lambda=$ $2 N+1$, for some $N \in \mathbb{N}$, to be determined later.
For solutions of (1.1) we seek an asymptotic two-scale expansion of the following form

$$
\begin{equation*}
u_{N}^{\varepsilon}(t, x):=\sum_{n=0}^{N} \varepsilon^{n} v_{n}\left(t, x, \frac{t}{\varepsilon}, \frac{x}{\varepsilon}\right) \tag{3.1}
\end{equation*}
$$

where

$$
\begin{equation*}
v_{n}(t, x, \tau, y):=\sum_{\sigma \in \mathcal{G}_{s}^{(2 n+1)}} A_{n, \sigma}(t, x, y) \mathbf{E}_{\sigma}(\tau, y) \tag{3.2}
\end{equation*}
$$

with $\mathbf{E}_{\sigma}(\tau, y):=\mathrm{e}^{\mathrm{i} \sigma \cdot(y,-\tau)}$, for $\sigma \in \mathcal{G}_{\mathcal{S}}^{(2 n+1)} \equiv \Sigma\left(\mathcal{S}^{2 n+1}\right)$. Note that in the most simple case, where $n=0$, we get that $\sigma=\left(k, E_{\ell}(k)\right)$, with $(k, \ell) \in \mathcal{S}$, which yields $\mathbf{E}_{\sigma}(\tau, y)=\mathrm{e}^{\mathrm{i}\left(k \cdot y-E_{\ell}(k) \tau\right)}$, a simple plane wave. Moreover, if $\sigma_{j} \in \mathcal{G}_{\mathcal{S}}^{\left(2 n_{j}+1\right)}$, for $j=1,2,3$, this yields

$$
\begin{equation*}
\mathbf{E}_{\sigma_{1}} \overline{\mathbf{E}}_{\sigma_{2}} \mathbf{E}_{\sigma_{3}}=\mathbf{E}_{\sigma_{1}-\sigma_{2}+\sigma_{3}} \tag{3.3}
\end{equation*}
$$

with $\sigma_{1}-\sigma_{2}+\sigma_{3} \in \mathcal{G}_{\mathcal{S}}^{\left(2\left(n_{1}+n_{2}+n_{3}\right)+3\right)}$. Of course we also impose that

$$
A_{n, \sigma}(\cdot, \cdot, y+\gamma)=A_{n, \sigma}(\cdot, \cdot, y), \quad \forall y \in \mathbb{R}^{d}, \gamma \in \Gamma
$$

### 3.1 The general strategy

As already said before, we shall only consider the case of simple eigenvalues $E_{\ell}(k)$, for simplicity. Plugging the ansatz (3.1), (3.2) into (1.1) and expanding in powers of $\varepsilon$, we formally obtain

$$
\mathrm{i} \varepsilon \partial_{t} u_{N}^{\varepsilon}-H_{\mathrm{per}}^{\varepsilon} u_{N}^{\varepsilon}-\varepsilon \kappa\left|u_{N}^{\varepsilon}\right|^{2} u_{N}^{\varepsilon}=\sum_{n=0}^{N} \varepsilon^{n} X_{n}+\operatorname{res}\left(u_{N}^{\varepsilon}\right),
$$

with the residual

$$
\begin{equation*}
\operatorname{res}\left(u_{N}^{\varepsilon}\right)=\sum_{n=N+1}^{3 N+1} \varepsilon^{n} X_{n} \tag{3.4}
\end{equation*}
$$

Introducing, for any (general) $\sigma \equiv(k, E) \in \mathcal{B} \times \mathbb{R}$, the operators

$$
\begin{equation*}
L_{0}^{\sigma}:=E-H_{\Gamma}(k), \quad L_{1}^{\sigma}:=\mathrm{i} \partial_{t}+\mathrm{i} k \cdot \nabla_{x}+\operatorname{div}_{x} \nabla_{y}, \quad L_{2}:=\frac{1}{2} \Delta_{x} \tag{3.5}
\end{equation*}
$$

this yields

$$
\begin{equation*}
X_{0}:=\sum_{\sigma \in \mathcal{G}_{s}^{(1)}}\left(L_{0}^{\sigma} A_{0, \sigma}\right) \mathbf{E}_{\sigma} \tag{3.6}
\end{equation*}
$$

and also

$$
\begin{align*}
X_{1}:= & \sum_{\substack{\sigma \in \mathcal{G}_{\mathcal{S}}^{(3)}}}\left(L_{0}^{\sigma} A_{1, \sigma}\right) \mathbf{E}_{\sigma}+\sum_{\sigma \in \mathcal{G}_{\mathcal{S}}^{(1)}}\left(L_{1}^{\sigma} A_{0, \sigma}\right) \mathbf{E}_{\sigma} \\
& -\kappa \sum_{\substack{\sigma_{j} \in G_{\mathcal{S}}^{(1)} ; \\
j=1,2,3}} A_{0, \sigma_{1}} \bar{A}_{0, \sigma_{2}} A_{0, \sigma_{3}} \mathbf{E}_{\sigma_{1}-\sigma_{2}+\sigma_{3}} \tag{3.7}
\end{align*}
$$

where we have used the relation (3.3). In general we get for $n=2, \ldots, 3 N+1$,

$$
\begin{align*}
X_{n}:= & \sum_{\sigma \in \mathcal{G}_{\mathcal{S}}^{(2 n+1)}}\left(L_{0}^{\sigma} A_{n, \sigma}\right) \mathbf{E}_{\sigma}+\sum_{\substack{\sigma \in \mathcal{G}_{\mathcal{S}}^{(2 n-1)}}}\left(L_{1}^{\sigma} A_{n-1, \sigma}\right) \mathbf{E}_{\sigma}+\sum_{\sigma \in \mathcal{G}_{\mathcal{S}}^{(2 n-3)}} L_{2} A_{n-2, \sigma} \mathbf{E}_{\sigma} \\
& -\kappa \sum_{n_{1}+n_{2}+n_{3}=n-1} \sum_{\substack{\sigma_{j} \in G_{\mathcal{S}}^{\left(2 n_{j}+1\right)} \\
j=1,2,3}} A_{n_{1}, \sigma_{1}} \bar{A}_{n_{2}, \sigma_{2}} A_{n_{3}, \sigma_{3}} \mathbf{E}_{\sigma_{1}-\sigma_{2}+\sigma_{3}} . \tag{3.8}
\end{align*}
$$

Here one should note that $A_{n, \sigma}(t, x, y) \equiv 0$, for all $n \geq N+1$, by assumption.
Now we shall subsequently construct $A_{n, \sigma}$ such that $X_{n} \equiv 0$ for $n=0, \ldots, N$. To this end we have to compare equal coefficients of $\mathbf{E}_{\sigma}$. We consequently obtain, from (3.6)-(3.8), equations of the form

$$
\begin{equation*}
L_{0}^{\sigma} A_{n, \sigma}=F_{n, \sigma}, \quad \sigma \in \mathcal{G}_{\mathcal{S}}^{(2 n+1)} \tag{3.9}
\end{equation*}
$$

where the r.h.s. $F_{n, \sigma}$ can be determined from the coefficients $\left(A_{m, \sigma}\right)_{\sigma \in \mathcal{G}_{\mathcal{S}}^{(2 m+1)}}$ for $m=0,1, \ldots, n-1$. More precisely

$$
\begin{align*}
F_{n, \sigma}:= & -L_{1}^{\sigma} A_{n-1, \sigma}-L_{2} A_{n-2, \sigma} \\
& +\kappa \sum_{n_{1}+n_{2}+n_{3}=n-1} \sum_{\substack{\sigma_{j} \in \mathcal{G}_{\left(2 n_{j}+1\right)} \\
\sigma_{1}-\sigma_{2}+\sigma_{3}=\sigma}} A_{n_{1}, \sigma_{1}} \bar{A}_{n_{2}, \sigma_{2}} A_{n_{3}, \sigma_{3}} . \tag{3.10}
\end{align*}
$$

Note that here the summation index has changed in comparison to (3.8). To proceed further we need to distinguish two possible cases:
Case I. On the one hand, for $\sigma \in \mathcal{G}_{\mathcal{S}}^{(2 n+1)} \backslash \mathcal{G}_{\mathcal{S}}^{(1)}$ the closure condition up to order $2 n+1$ implies invertibility of $L_{0}^{\sigma}$, i.e. $\left(L_{0}^{\sigma}\right)^{-1} \in \operatorname{Lin}\left(L^{2}(\mathcal{Y}), H^{2}(\mathcal{Y})\right)$ and we obtain

$$
\begin{equation*}
A_{n, \sigma}(t, x, y)=\left(L_{0}^{\sigma}\right)^{-1} F_{n, \sigma}(t, x, y) \tag{3.11}
\end{equation*}
$$

The corresponding modes are called non-resonant.
Case II. On the other hand, if $\sigma \in \mathcal{G}_{\mathcal{S}}^{(1)}$, then $L_{0}^{\sigma}$ has a nontrivial kernel. In order to distinguish this case more prominently from the one above, we shall from now
on use the notation $\varsigma \equiv \sigma \in \mathcal{G}_{\mathcal{S}}^{(1)}$, which also characterizes the basic resonant modes $\mu \in \mathcal{S}$ via $\varsigma=\Sigma(\mu)$.
Using the orthogonal projections $\mathbb{P}_{\varsigma}$ onto this kernel the necessary and sufficient solvability condition for (3.9) in this case is then given by

$$
\begin{equation*}
\mathbb{P}_{\varsigma} F_{n, \varsigma}(t, x, y)=0 \tag{3.12}
\end{equation*}
$$

which yields (recall that $\operatorname{dim}\left(\operatorname{ran} \mathbb{P}_{\varsigma}\right)=1$, by assumption)

$$
\begin{equation*}
\left\langle\chi_{\varsigma}, F_{n, \varsigma}\right\rangle_{L^{2}(\mathcal{Y})}=0, \quad \text { for } 0 \neq \chi_{\varsigma} \in \operatorname{ran} \mathbb{P}_{\varsigma} . \tag{3.13}
\end{equation*}
$$

Under this condition, we consequently obtain

$$
\begin{equation*}
A_{n, \varsigma}(t, x, y)=a_{n, \varsigma}(t, x) \chi_{\varsigma}(y)+A_{n, \varsigma}^{\perp}(t, x, y) \tag{3.14}
\end{equation*}
$$

with

$$
\begin{equation*}
A_{n, \varsigma}^{\perp}:=\left(L_{0}^{\varsigma}\right)^{-1}\left(1-\mathbb{P}_{\varsigma}\right) F_{n, \varsigma} . \tag{3.15}
\end{equation*}
$$

Note that here $a_{n, \varsigma}$ is still undetermined. However the condition (3.13) provides a partial differential equation for $a_{n-1, \varsigma}$, the so far undetermined part of the previous step.

Remark 3.1. In the case of non-simple eigenvalues $E_{\ell}(k)$, i.e. $\operatorname{dim}\left(\operatorname{ran} \mathbb{P}_{\varsigma}\right)=R>1$, we can simply use a smooth orthonormal basis $\left\{\chi_{\varsigma, r}\right\}_{r=1}^{R}$ of $\operatorname{ran} \mathbb{P}_{\varsigma}$ and generalize the above given formulas (3.12), (3.14) accordingly.

### 3.2 Explicit calculations

In the following we shall determine the approximate solution in more detail, by following the above described strategy.
$n=0$ : We need to solve $X_{0} \equiv 0$ and immediately note that in this case, Case I above is obsolete. Then, in Case II, equation (3.6) implies

$$
L_{0}^{\varsigma} A_{0, \varsigma}=0, \quad \varsigma \in \mathcal{G}_{\mathcal{S}}^{(1)}
$$

and thus (3.14) simplifies to

$$
\begin{equation*}
A_{0, \varsigma}=a_{0, \varsigma}(t, x) \chi_{\varsigma}(y), \tag{3.16}
\end{equation*}
$$

with $a_{0, \varsigma}$ still to be determined.
$n=1$ : In the next step we have to solve $X_{1} \equiv 0$. In Case I, i.e. for $\sigma \in \mathcal{G}_{\mathcal{S}}^{(3)} \backslash \mathcal{G}_{\mathcal{S}}^{(1)}$, the equations (3.7) and (3.11) imply

$$
\begin{align*}
A_{1, \sigma} & =\kappa\left(L_{0}^{\sigma}\right)^{-1}\left(\sum_{\substack{\varsigma_{1}, \varsigma_{2}, \varsigma_{3} \in \mathcal{G}_{\mathcal{S}}^{(1)} \\
\sigma=\varsigma_{1}-\varsigma_{2}+\varsigma_{3}}} A_{0, \varsigma_{1}} \bar{A}_{0, \varsigma_{2}} A_{0, \varsigma_{3}}\right) \\
& =\kappa \sum_{\substack{\varsigma_{1}, \varsigma_{2}, \varsigma_{3} \in \mathcal{G}_{2}^{(1)}, \sigma=\varsigma_{1}-\varsigma_{2}+\varsigma_{3}}}^{a_{0, \varsigma_{1}} \bar{a}_{0, \varsigma_{2}} a_{0, \varsigma_{3}}\left(L_{0}^{\sigma}\right)^{-1}\left(\chi_{\varsigma_{1}} \bar{\chi}_{\varsigma_{2}} \chi_{\varsigma_{3}}\right),} \tag{3.17}
\end{align*}
$$

where for the second equality we simply insert (3.16). We proceed with Case II: To this end the solvability condition (3.12) for $\varsigma \in \mathcal{G}_{\mathcal{S}}^{(1)}$ allows us to determine the so far still unknown $a_{0, \varsigma}$, obtained before. By (3.13), this yields

$$
\begin{equation*}
\int_{\mathcal{Y}} \bar{\chi}_{\varsigma}(y) F_{1, \varsigma}(t, x, y) \mathrm{d} y=0 \tag{3.18}
\end{equation*}
$$

where

$$
\begin{equation*}
F_{1, \varsigma}=L_{1}^{\varsigma} A_{0, \varsigma}-\kappa \sum_{\substack{\varsigma_{1}, \varsigma_{2}, \varsigma_{3} \in \mathcal{G}_{\S}^{(1)}, \varsigma=\varsigma_{1}-\varsigma_{2}+\varsigma_{3}}} A_{0, \varsigma_{1}} \bar{A}_{0, \varsigma_{2}} A_{0, \varsigma_{3}} . \tag{3.19}
\end{equation*}
$$

From the definition of $L_{1}^{\varsigma}$ in (3.5) and using the following basic identity

$$
\left\langle\chi_{\ell},\left(-\mathrm{i} \nabla_{y}+k\right) \chi_{\ell}\right\rangle_{L^{2}(\mathcal{Y})}=\nabla_{k} E_{\ell}(k),
$$

a straightforward calculation shows, cf. the appendix of [7], that (3.18) can be written as

$$
\begin{equation*}
\partial_{t} a_{0, \varsigma}+\vartheta_{\varsigma} \cdot \nabla_{x} a_{0, \varsigma}+\sum_{\substack{\varsigma_{1}, \varsigma_{2}, \varsigma_{\zeta} \in \mathcal{S}_{\mathcal{S}}^{(1)} ; \\ \varsigma=\varsigma_{1}-\varsigma_{2}+\varsigma_{3}}} \mathrm{i} \kappa_{\left(\varsigma_{1}, \varsigma_{2}, \varsigma_{3}, \varsigma\right)} a_{0, \varsigma_{1}} \bar{a}_{0, \varsigma_{2}} a_{0, \varsigma_{3}}=0, \tag{3.20}
\end{equation*}
$$

where we denote

$$
\kappa_{\left(\varsigma_{1}, \varsigma_{2}, \varsigma_{3}, \varsigma\right)}:=\kappa \int_{\mathcal{Y}} \chi_{\varsigma_{1}}(y) \bar{\chi}_{\varsigma_{2}}(y) \chi_{\varsigma_{3}}(y) \bar{\chi}_{\varsigma}(y) \mathrm{d} y .
$$

Since any $\varsigma \in \mathcal{G}_{\mathcal{S}}^{(1)}$ corresponds to a unique $m=1, \ldots, M$, via $\varsigma=\Sigma\left(\mu_{m}\right)=$ $\left(k_{m}, E_{\ell_{m}}\left(k_{m}\right)\right)$, with $\mu_{m}=\left(k_{m}, \ell_{m}\right) \in \mathcal{S}$, we can shortly write

$$
a_{0, \varsigma}(t, x) \equiv a_{m}(t, x), \quad \chi_{\varsigma}(y) \equiv \chi_{\ell_{m}}\left(y ; k_{m}\right), \quad \vartheta_{\varsigma} \equiv \vartheta_{m} .
$$

Hence the amplitude equations (3.20) can be equivalently written in the form (1.5) used before.

In summary we have now fully determined the expressions (3.16) and (3.17), and from (3.14) we finally get that

$$
A_{1, \varsigma}(t, x, y)=a_{1, \varsigma}(t, x) \chi_{\varsigma}(y)+\left(L_{0}^{\varsigma}\right)^{-1}\left(1-\mathbb{P}_{\varsigma}\right) F_{1, \varsigma}(t, x, y),
$$

where again the coefficients $a_{1, \varsigma}$ are still arbitrary and have to be determined by the solvability condition for $n=2$. Since this yields an initial value problem for $a_{1, \varsigma}$ (see below) we are free to choose its value at time $t=0$. For simplicity we shall put $a_{1, \varsigma}(0, \cdot)=0$.
$n \geq 2$ : From here we proceed inductively, as described above, by solving $X_{n} \equiv 0$. The only difference that occurs is that in Case II the coefficients $a_{n, \varsigma}$ do not solve a nonlinear initial value problem, but rather a system of linear, inhomogeneous
transport equations. Indeed, lengthy calculations show that

$$
\begin{align*}
& \partial_{t} a_{n, \varsigma}+\vartheta_{\varsigma} \cdot \nabla_{x} a_{n, \varsigma}+\sum_{\substack{\bar{\zeta}, \varsigma_{1}, \varsigma_{2} \in \mathcal{S}_{\mathcal{S}}^{(1)} ; \\
\bar{\zeta}-\varsigma_{1}+\varsigma_{2}=\varsigma}} 2 \mathrm{i} \kappa_{\left(\tilde{\zeta}, \varsigma_{1}, \varsigma_{2}, \varsigma\right)} a_{n, \bar{\varsigma}} \bar{a}_{0, \varsigma_{1}} a_{0, \varsigma_{2}}  \tag{3.21}\\
& +\sum_{\substack{\varsigma_{1}, \varsigma_{2}, \varsigma_{1} \in \mathcal{G}_{\mathcal{S}}^{(1)}: \\
\varsigma_{1}-\tilde{\zeta}+\varsigma_{2}=\varsigma}} \mathrm{i} \kappa_{\left(\varsigma_{1}, \tilde{\zeta}_{,}, \varsigma_{2}, \varsigma\right)} a_{0, \varsigma_{1}} \bar{a}_{n, \tilde{\zeta}} a_{0, \varsigma_{2}}+\mathrm{i} \Theta_{n}=0,
\end{align*}
$$

with source term

$$
\begin{aligned}
& \Theta_{n}:=-\mathbb{P}_{\varsigma}\left(L_{1}^{\varsigma} A_{n, \varsigma}^{\perp}+L_{2} A_{n-1, \varsigma}\right)
\end{aligned}
$$

$$
\begin{aligned}
& +\kappa \mathbb{P}_{\varsigma}\left(\sum_{\substack{1_{1}, \varsigma_{2} \in \mathcal{G}_{1}^{(1)} \ngtr \sigma: \\
\sigma-\varsigma_{1}+\varsigma_{2}=\varsigma}} 2 A_{n, \sigma} \bar{A}_{0, \varsigma_{1}} A_{0, \varsigma_{2}}+\sum_{\substack{\left.\varsigma_{1}, \varsigma_{2} \in \mathcal{G}_{1}^{(1)}\right) \neq \sigma: \\
\varsigma_{1}-\sigma+\varsigma_{2}=\varsigma}} A_{0, \varsigma_{1}} \bar{A}_{n, \sigma} A_{0, \varsigma_{2}}\right) \\
& +\kappa \mathbb{P}_{\varsigma}\left(\sum_{\substack{n_{1}+n_{2}+n_{3}=n: \\
n_{j} \leq n-1}} \sum_{\substack{\sigma_{j} \in \mathcal{G}_{\mathcal{S}}^{\left(2 n_{j}+1\right)}: \\
\sigma_{1}-\sigma_{2}+\sigma_{3}=\varsigma}} A_{n_{1}, \sigma_{1}} \bar{A}_{n_{2}, \sigma_{2}} A_{n_{3}, \sigma_{3}}\right) .
\end{aligned}
$$

Again we shall put $a_{n, s}(0, \cdot)=0$ for simplicity, since we are free to choose the initial values for (3.21). Finally, we note that, in order to satisfy $X_{N} \equiv 0$ (i.e. in the last step), we do not need to determine the corresponding $a_{N, \varsigma}$ and thus we can impose $a_{N, \varsigma}(t, \cdot) \equiv 0$.

## 4 Justification of the amplitude equations

In order to obtain our main result, Theorem 4.5, we have to justify the above given formal calculations rigorously. In particular we need a nonlinear stability result on our approximation.

### 4.1 Estimates on the approximate solution and on the residual

In the above section we derived an approximate solution $u_{N}^{\varepsilon}$ which admits an asymptotic expansion of the following form

$$
\begin{align*}
u_{N}^{\varepsilon}(t, x)= & \sum_{\varsigma \in G_{\mathcal{S}}^{(1)}} a_{0, \varsigma}(t, x) \chi_{\varsigma}\left(\frac{x}{\varepsilon}\right) \mathbf{E}_{\varsigma}\left(\frac{t}{\varepsilon}, \frac{x}{\varepsilon}\right) \\
& +\sum_{n=1}^{N} \varepsilon^{n} \sum_{\varsigma \in G_{\mathcal{S}}^{(1)}}\left(a_{n, \varsigma}(t, x) \chi_{\varsigma}\left(\frac{x}{\varepsilon}\right)+A_{n, \varsigma}^{\perp}\left(t, x, \frac{x}{\varepsilon}\right)\right) \mathbf{E}_{\varsigma}\left(\frac{t}{\varepsilon}, \frac{x}{\varepsilon}\right)  \tag{4.1}\\
& +\sum_{n=1}^{N} \varepsilon^{n} \sum_{\sigma \in G_{\mathcal{S}}^{(2 n+1)} \backslash G_{\mathcal{S}}^{(1)}} A_{n, \sigma}\left(t, x, \frac{x}{\varepsilon}\right) \mathbf{E}_{\sigma}\left(\frac{t}{\varepsilon}, \frac{x}{\varepsilon}\right)
\end{align*}
$$

where the first two lines on the r.h.s. include only resonant modes while the last one takes into account the generated non-resonant terms. To this end we required the appearing finite set of modes to be closed of order $\Lambda=2 N+1$.
To proceed further let us, for $s \in[0, \infty)$, introduce the scaled Sobolev spaces

$$
H_{\varepsilon}^{s}:=\left\{f^{\varepsilon} \in L^{2}\left(\mathbb{R}^{d}\right) ;\left\|f^{\varepsilon}\right\|_{H_{\varepsilon}^{s}}<\infty\right\}
$$

with

$$
\left\|f^{\varepsilon}\right\|_{H_{\varepsilon}^{s}}^{2}:=\int_{\mathbb{R}^{d}}\left(1+|\varepsilon p|^{2}\right)^{s}|\widehat{f}(p)|^{2} \mathrm{~d} p
$$

Here $\widehat{f}$ denotes the usual Fourier transform of $f$ in $L^{2}\left(\mathbb{R}^{d}\right)$. Note that in $H_{\varepsilon}^{s}$ the following Gagliardo-Nirenberg inequality holds

$$
\begin{equation*}
\forall s>\frac{d}{2} \exists C_{\infty}>0: \quad\|f\|_{L^{\infty}} \leq C_{\infty}\|f\|_{H^{s}} \leq C_{\infty} \varepsilon^{-d / 2}\|f\|_{H_{\varepsilon}^{s}}, \tag{4.2}
\end{equation*}
$$

where the factor $\varepsilon^{-d / 2}$ is easily obtained by scaling. In the following lemma we collect the a-priori estimates on $u_{N}^{\varepsilon}$ needed in order to prove our main result.
Lemma 4.1. For $d, N \in \mathbb{N}$ and $s \in[0, \infty)$ let $K=\max \left\{0, s+\frac{d}{2}-2\right\}$ and $S>$ $N+s+\frac{d}{2}$, or, if $d=1, S \geq N+s+1$. Moreover, assume that $\partial^{\alpha} V_{\Gamma} \in L^{\infty}(\mathcal{Y})$, for all $|\alpha| \leq K$, and let $\left(a_{1}, \ldots, a_{M}\right) \in C^{0}\left([0, T), H^{S}\left(\mathbb{R}^{d}\right)\right)^{M}$ be a solution of the amplitude equations (1.5).
Then, the approximate solution $u_{N}^{\varepsilon}$, given in (4.1), satisfies the following estimates. For each $T_{*} \in(0, T)$ there exist positive constants $C_{a}, C_{b}, C_{r}>0$, such that, for $\varepsilon \in(0,1),|\alpha| \leq s$, and $t \in\left[0, T_{*}\right]$, it holds

$$
\left\|\left(\varepsilon \partial_{x}\right)^{\alpha} u_{N}^{\varepsilon}(t, \cdot)\right\|_{L^{\infty}} \leq C_{a}, \quad\left\|u_{N}^{\varepsilon}(t, \cdot)\right\|_{H_{\varepsilon}^{s}} \leq C_{b}, \quad\left\|\operatorname{res}\left(u_{N}^{\varepsilon}\right)(t, \cdot)\right\|_{H_{\varepsilon}^{s}} \leq C_{r} \varepsilon^{N+1}
$$

Proof. Since $u_{N}^{\varepsilon}$ is given by (4.1), in order to prove the estimates of the lemma we need to establish the regularity of $A_{n, \sigma}$ for $n=0, \ldots, N, \sigma \in \mathcal{G}_{\mathcal{S}}^{(2 n+1)}$, which
have been calculated formally in Section 3. Under the assumptions $\left(a_{1}, \ldots, a_{M}\right) \in$ $C^{0}\left([0, T), H^{S}\left(\mathbb{R}^{d}\right)\right)^{M}$ and $\partial^{\alpha} V_{\Gamma} \in L^{\infty}(\mathcal{Y})$, for $|\alpha| \leq K$, we shall first prove that it holds

$$
\begin{equation*}
A_{n, \sigma} \in C^{0}\left([0, T), H^{S-n}\left(\mathbb{R}^{d}, H^{K+2}(\mathcal{Y})\right)\right) . \tag{4.3}
\end{equation*}
$$

Indeed, by (3.16), $A_{0, \varsigma}(t, x, y)=a_{0, \varsigma}(t, x) \chi_{\varsigma}(y)$, where $a_{0, \varsigma}$ solves (3.20) and $\chi_{\varsigma}$ is given by (2.2), and we get immediately (4.3) for $n=0$ (cf. Lemma 2.4 and Lemma 4.3 below).

From here it follows, on the one hand, that for $\sigma \in \mathcal{G}_{\mathcal{S}}^{(3)} \backslash \mathcal{G}_{\mathcal{S}}^{(1)}$, by (3.9) and (3.17), $F_{1, \sigma}=L_{0}^{\sigma} A_{1, \sigma} \in C^{0}\left([0, T), H^{S}\left(\mathbb{R}^{d}, H^{\tilde{s}}(\mathcal{Y})\right)\right)$ with $\tilde{s}=K+2$ and hence, by (3.10), $A_{1, \sigma} \in C^{0}\left([0, T), H^{S-1}\left(\mathbb{R}^{d}, H^{\tilde{s}}(\mathcal{Y})\right)\right)$. On the other hand, for $\varsigma \in \mathcal{G}_{\mathcal{S}}^{(1)}$ we obtain by (3.19) $F_{1, \varsigma} \in C^{0}\left([0, T), H^{S-1}\left(\mathbb{R}^{d}, H^{\tilde{s}}(\mathcal{Y})\right)\right.$ ), and thus, by (3.15), $A_{1, \varsigma}^{\perp} \in$ $C^{0}\left([0, T), H^{S-1}\left(\mathbb{R}^{d}, H^{\tilde{s}}(\mathcal{Y})\right)\right)$. Moreover, $a_{1, \varsigma} \chi_{\varsigma} \in C^{0}\left([0, T), H^{S-1}\left(\mathbb{R}^{d}, H^{\tilde{s}}(\mathcal{Y})\right)\right)$, since $a_{1, \varsigma}$ solves the linear inhomogeneous transport equation (3.21) (for $n=1$ ) with initial condition $a_{1, \varsigma}(0, \cdot)=0$, coefficients in $C^{0}\left([0, T), H^{S}\left(\mathbb{R}^{d}\right)\right.$ ), and source term $\Theta_{1} \in C^{0}\left([0, T), H^{S-2}\left(\mathbb{R}^{d}\right)\right)$. Hence, by (3.14), we obtain (4.3) for $n=1$. Proceeding inductively, we obtain (4.3) for all $n=0, \ldots, N$ via the above given steps.
Having established (4.3), we aim to show the first estimate of the lemma. To this end we have to guarantee that for each $n=0, \ldots, N, \sigma \in \mathcal{G}_{\mathcal{S}}^{(2 n+1)}$, there exists a $C>0$, such that

$$
\left\|\left(\varepsilon \partial_{x}\right)^{\alpha} A_{n, \sigma}\left(t, \cdot, \frac{\dot{x}}{\varepsilon}\right)\right\|_{L^{\infty}} \leq C
$$

holds true for all $\varepsilon \in(0,1),|\alpha| \leq s$, and $t \in\left[0, T_{*}\right]$. Using the Gagliardo-Nirenberg inequality (4.2), we therefore require

$$
\begin{equation*}
\left(\varepsilon \partial_{x}\right)^{\alpha} A_{n, \sigma}\left(t, \cdot, \frac{\dot{ }}{\varepsilon}\right) \in H^{m}\left(\mathbb{R}^{d}\right) \quad \text { with } m>\frac{d}{2}, \tag{4.4}
\end{equation*}
$$

for all $|\alpha| \leq s$, i.e.

$$
A_{n, \sigma}\left(t, \cdot, \frac{\cdot}{\varepsilon}\right) \in H_{\varepsilon}^{m^{*}}\left(\mathbb{R}^{d}\right) \quad \text { with } m^{*}>s+\frac{d}{2} .
$$

Thus, by (4.3), we need $m^{*}=S-n+\tilde{s}>s+\frac{d}{2}$, that is, either $S-n>|\beta|+\frac{d}{2}$ and $\tilde{s} \geq s-|\beta|$, or $S-n \geq|\beta|$ and $\tilde{s}>s-|\beta|+\frac{d}{2}$, for all $|\beta| \leq s$ and all $n=0, \ldots, N$. It turns out that in order to satisfy either of these conditions, the optimal conditions we have to impose concerning the regularity of $A_{n, \sigma}$, given in (4.3), are $S-n>s+\frac{d}{2}$ and $\tilde{s}=K+2>s+\frac{d}{2}$, for $n=0, \ldots, N$, which directly yields the assumptions on $S$ and $K$ stated in the lemma.
The second estimate of the lemma then follows immediately by (4.4). In order to obtain the third estimate, having in mind the definitions (3.4) and (3.8) with $A_{n, \sigma} \equiv 0$ for $n \geq N+1$ and $a_{N, \varsigma} \equiv 0$ (cf. the note after (3.21)), it is necessary and sufficient to assure additionally that $\Delta a_{N-1, \varsigma} \in H^{s}\left(\mathbb{R}^{d}\right)$, i.e., $a_{0, \varsigma} \in H^{N+s+1}\left(\mathbb{R}^{d}\right)$. This yields the additional condition $S \geq N+s+1$, and concludes the proof.

Remark 4.2. Note that in order to derive the above given a-priori estimates for $u_{N}$, the required regularity imposed on $\chi_{\varsigma}$ is indeed independent of $N$.

We shall also need the following result on the linear time-evolution.
Lemma 4.3. For $\varepsilon \in(0,1)$ denote the unitary propagator corresponding to the linear Hamiltonian $H_{\mathrm{per}}^{\varepsilon}$, defined in (1.3), by

$$
U^{\varepsilon}(t):=\mathrm{e}^{-\mathrm{i} H_{\mathrm{per}}^{\varepsilon} t / \varepsilon} .
$$

For $s \in[0, \infty)$ assume $\partial^{\alpha} V_{\Gamma} \in L^{\infty}(\mathcal{Y})$ for all $|\alpha| \leq \max \{0, s-2\}$. Then there exists a $C_{l}>0$, such that

$$
\begin{equation*}
\left\|U^{\varepsilon}(t) f^{\varepsilon}\right\|_{H_{\varepsilon}^{s}} \leq C_{l}\left\|f^{\varepsilon}\right\|_{H_{\varepsilon}^{s}} \quad \text { for all } t \in \mathbb{R} \tag{4.5}
\end{equation*}
$$

Proof. Recalling the definition of $H_{\text {per }}^{\varepsilon}$, given in (1.3), we clearly have the basic $L^{2}$ estimate,

$$
\left\|f^{\varepsilon}(t)\right\|_{L^{2}} \equiv\left\|U^{\varepsilon}(t) f^{\varepsilon}\right\|_{L^{2}}=\left\|f^{\varepsilon}\right\|_{L^{2}}, \quad \text { for all } t \in \mathbb{R},
$$

since $H_{\mathrm{per}}^{\varepsilon}$ is self-adjoint. Moreover, since $U^{\varepsilon}(t)$ obviously commutes with (its generator) $H_{\text {per }}^{\varepsilon}$, so does any power of the latter and we therefore obtain

$$
\left\|\left(H_{\mathrm{per}}^{\varepsilon}\right)^{s / 2} f^{\varepsilon}(t)\right\|_{L^{2}}=\left\|\left(H_{\mathrm{per}}^{\varepsilon}\right)^{s / 2} f^{\varepsilon}\right\|_{L^{2}}, \quad \text { for all } s, t \in \mathbb{R}
$$

Without loss of generality we assume here that $H_{\text {per }}^{\varepsilon} \geq 1$, otherwise we may add a constant independent of $\varepsilon \in(0,1)$.
For $s=1$ this is nothing but energy conservation. Using the conservation of the $L^{2}$ norm and the condition $V_{\Gamma} \in L^{\infty}(\mathcal{Y})$ we immediately obtain the desired result (4.5) for $s=1$ with $C_{l}=\left(1+4\left\|V_{\Gamma}\right\|_{\infty}\right)^{1 / 2}$.
For general $s \in \mathbb{N} \backslash\{1\}$, we use integration by parts and the assumed regularity of $V_{\Gamma}$ to find $\varepsilon$-independent constants $C_{1}, C_{2}>0$ such that

$$
C_{1}\left\|\left(H_{\mathrm{per}}^{\varepsilon}\right)^{s / 2} f^{\varepsilon}(t)\right\|_{L^{2}} \leq\left\|f^{\varepsilon}(t)\right\|_{H_{\varepsilon}^{s}} \leq C_{2}\left\|\left(H_{\mathrm{per}}^{\varepsilon}\right)^{s / 2} f^{\varepsilon}(t)\right\|_{L^{2}},
$$

i.e. we have equivalence of the norms uniformly in $\varepsilon$. For general $s \geq 0$ the same holds true by interpolation. In summary we obtain boundedness of the unitary group $U^{\varepsilon}(t)$ in all $H_{\varepsilon}^{s}$ and the assertion of the lemma is proved.

### 4.2 Stability of the approximation

First let us recall the following Moser-type lemma, cf. [22, Lemma 8.1], which we shall use in the proof of Theorem 4.5 below.

Lemma 4.4. Let $R>0, s \in[0, \infty)$, and $\mathcal{N}(z)=\kappa|z|^{2} z$ with $\kappa \in \mathbb{R}$. Then there exists a $C_{s}=C_{s}(R, s, d, \kappa)>0$ such that if

$$
\left\|(\varepsilon \partial)^{\alpha} f\right\|_{L^{\infty}} \leq R, \forall|\alpha| \leq s \text { and }\|g\|_{L^{\infty}} \leq R,
$$

then

$$
\|\mathcal{N}(f+g)-\mathcal{N}(f)\|_{H_{\varepsilon}^{s}} \leq C_{s}\|g\|_{H_{\varepsilon}^{s}} .
$$

Proof. The proof can be found in [22] for $s \in \mathbb{N}$ and follows by interpolation for general $s \in[0, \infty)$.

With the above results at hand, we are now able to establish our main result, which justifies rigorously the validity of the amplitude equations (1.5), describing the macroscopic dynamics of $M$ modulated pulses for a closed mode system of order $\Lambda=2 N+1$ (with $N$ depending on $d$ ).

Theorem 4.5. For $d \in \mathbb{N}$ choose $s, S \in[0, \infty)$ and $N \in \mathbb{N}$ such that $N, s>\frac{d}{2}$ and $S>N+s+\frac{d}{2}$, or, if $d=1, S \geq N+s+1$. Let $K=\max \left\{0, s+\frac{d}{2}-2\right\}$ and assume $\partial^{\alpha} V_{\Gamma} \in L^{\infty}(\mathcal{Y})$ for all $|\alpha| \leq K$. Moreover, let the finite system of modes $\mathcal{S}=\left\{\mu_{1}, \ldots, \mu_{M}: M \in \mathbb{N}\right\}$ be closed of order $\Lambda=2 N+1$.
Then for any solution $\left(a_{1}, \ldots, a_{M}\right) \in C^{0}\left([0, T), H^{S}\left(\mathbb{R}^{d}\right)\right)^{M}$ of the amplitude equations (1.5), and any $t_{*} \in(0, T), \beta \in\left(\frac{d}{2}, N\right], c>0$, there exist an $\varepsilon_{0} \in(0,1)$ and a $C>0$, such that for all $\varepsilon \in\left(0, \varepsilon_{0}\right)$, the approximate solution $u_{N}^{\varepsilon}$, constructed above, and any exact solution $u^{\varepsilon}$ of (1.1) with

$$
\left\|u^{\varepsilon}(0, \cdot)-u_{N-1}^{\varepsilon}(0, \cdot)\right\|_{H_{\varepsilon}^{s}} \leq c \varepsilon^{\beta}
$$

satisfy

$$
\left\|u^{\varepsilon}(t, \cdot)-u_{N-1}^{\varepsilon}(t, \cdot)\right\|_{H_{\varepsilon}^{s}} \leq C \varepsilon^{\beta} \quad \text { for all } t \in\left[0, t_{*}\right] .
$$

Remark 4.6. Note, that the approximate solution $u_{N}^{\varepsilon}$ contains terms up to order $\mathcal{O}\left(\varepsilon^{N}\right)$, whereas the above given error estimates include only $u_{N-1}^{\varepsilon}$. The reason for this is that for our proof we need to work with $u_{N}^{\varepsilon}$ but still the obtained error is only of the order $\mathcal{O}\left(\varepsilon^{\beta}\right)$ with $\beta \leq N$. We can therefore eventually neglect the last term in $u_{N}^{\varepsilon}$, or, loosely speaking, we can move it to the right hand side of the above given estimate.

Proof. We write the exact solution of (1.1) in the form

$$
u^{\varepsilon}(t, x)=u_{N}^{\varepsilon}(t, x)+\varepsilon^{\beta} w^{\varepsilon}(t, x)
$$

and denote $\varrho^{\varepsilon}(t):=\left\|w^{\varepsilon}(t)\right\|_{H_{\varepsilon}^{s}}$. Then, clearly, $\varrho^{\varepsilon}(0) \leq c$, and we will show that there exist $C>0, \varepsilon_{0} \in(0,1)$, such that, for all $\varepsilon \in\left(0, \varepsilon_{0}\right]$, it holds $\varrho^{\varepsilon}(t) \leq C$ for $t \in\left[0, t_{*}\right]$. Inserting this ansatz into (1.1), written as

$$
\mathrm{i} \partial_{t} u^{\varepsilon}=-\frac{\varepsilon}{2} \Delta u^{\varepsilon}+\frac{1}{\varepsilon} V_{\Gamma}\left(\frac{x}{\varepsilon}\right) u^{\varepsilon}+\mathcal{N}\left(u^{\varepsilon}\right), \quad \text { where } \mathcal{N}(z)=\kappa|z|^{2} z,
$$

and applying Duhamel's formula we get

$$
\begin{aligned}
w^{\varepsilon}(t)= & U^{\varepsilon}(t) w^{\varepsilon}(0)+\varepsilon^{-\beta} \int_{0}^{t} U^{\varepsilon}(t-\tau)\left(\mathcal{N}\left(u_{N}^{\varepsilon}(\tau)+\varepsilon^{\beta} w^{\varepsilon}(\tau)\right)-\mathcal{N}\left(u_{N}^{\varepsilon}(\tau)\right)\right) \mathrm{d} \tau \\
& -\varepsilon^{-(\beta+1)} \int_{0}^{t} U^{\varepsilon}(t-\tau) \operatorname{res}\left(u_{N}^{\varepsilon}(\tau)\right) \mathrm{d} \tau
\end{aligned}
$$

where $\operatorname{res}\left(u_{N}^{\varepsilon}\right)$ is defined in equation (3.4). Hence, using Lemma 4.1 and Lemma 4.3 to estimate the residual and the linear semi-group, respectively, we obtain

$$
\varrho^{\varepsilon}(t) \leq C_{l} c+C_{l} C_{r} \varepsilon^{N-\beta} t+\varepsilon^{-\beta} \int_{0}^{t} C_{l}\left\|\mathcal{N}\left(u_{N}^{\varepsilon}(\tau)+\varepsilon^{\beta} w^{\varepsilon}(\tau)\right)-\mathcal{N}\left(u_{N}^{\varepsilon}(\tau)\right)\right\|_{H_{\varepsilon}^{s}} \mathrm{~d} \tau
$$

since $\varrho^{\varepsilon}(0) \leq c$, by assumption. Using $N \geq \beta$ and $\varepsilon \in(0,1)$, we consequently obtain

$$
\varrho^{\varepsilon}(t) \leq C_{l}\left(c+C_{r} t_{*}\right)+\varepsilon^{-\beta} \int_{0}^{t} C_{l}\left\|\mathcal{N}\left(u_{N}^{\varepsilon}(\tau)+\varepsilon^{\beta} w^{\varepsilon}(\tau)\right)-\mathcal{N}\left(u_{N}^{\varepsilon}(\tau)\right)\right\|_{H_{\varepsilon}^{s}} \mathrm{~d} \tau
$$

for $t \leq t_{*}$.
Now, we set $C:=C_{l}\left(c+C_{r} t_{*}\right) \mathrm{e}^{C_{l} C_{s} t_{*}}$ and choose a $D>\max \{c, C\}$. Then, since $D>c \geq \varrho^{\varepsilon}(0)$ and $\varrho^{\varepsilon}(t)$ is continuous, there exists, for every $\varepsilon \in(0,1)$, a positive time $t_{D}^{\varepsilon}>0$, such that $\varrho^{\varepsilon}(t) \leq D$ for $t \leq t_{D}^{\varepsilon}$.
The Gagliardo-Nirenberg inequality (4.2) yields, for $s>d / 2$, that

$$
\left\|\varepsilon^{\beta} w^{\varepsilon}(t)\right\|_{L^{\infty}} \leq \varepsilon^{\beta-d / 2} C_{\infty} D \quad \text { for } t \leq t_{D}^{\varepsilon} .
$$

Hence, using $\beta-d / 2>0$ there exists an $\varepsilon_{0} \in(0,1)$, such that

$$
\left\|\varepsilon^{\beta} w^{\varepsilon}(t)\right\|_{L^{\infty}} \leq C_{a}, \quad \text { for } \varepsilon \in\left(0, \varepsilon_{0}\right] \text { and } t \leq t_{D}^{\varepsilon} .
$$

Moreover, by Lemma 4.1 we have $\left\|(\varepsilon \partial)^{\alpha} u_{N}^{\varepsilon}(t)\right\|_{L^{\infty}} \leq C_{a}$ for $|\alpha| \leq s, \varepsilon \in(0,1)$, and $t<T$. Thus, we can apply Lemma 4.4 (with $R=C_{a}$ ) in order to estimate the nonlinear term and obtain

$$
\varrho^{\varepsilon}(t) \leq C_{l}\left(c+C_{r} t_{*}\right)+C_{l} C_{s} \int_{0}^{t} \varrho^{\varepsilon}(\tau) \mathrm{d} \tau, \quad \text { for } \varepsilon \in\left(0, \varepsilon_{0}\right] \text { and } t \leq t_{D}^{\varepsilon} .
$$

Gronwall's lemma then yields

$$
\varrho^{\varepsilon}(t) \leq C_{l}\left(c+C_{r} t_{*}\right) \mathrm{e}^{C_{l} C_{s} t} \leq C, \quad \text { for } \varepsilon \in\left(0, \varepsilon_{0}\right] \text { and } t \leq t_{*} \text {. }
$$

Since $C<D$, we conclude that the assumptions needed in order to apply Lemma 4.4 are fulfilled for all $\varepsilon \in\left(0, \varepsilon_{0}\right]$ and all $t \leq t_{*}$, that is $t_{D}^{\varepsilon} \geq t_{*}$. Thus the above given estimate proves that

$$
\left\|u^{\varepsilon}(t, \cdot)-u_{N}^{\varepsilon}(t, \cdot)\right\|_{H_{\varepsilon}^{s}}=\mathcal{O}\left(\varepsilon^{\beta}\right), \quad \text { for } t \in\left[0, t_{*}\right] .
$$

However, since $N \geq \beta$ and

$$
\left\|u_{N}^{\varepsilon}(t, \cdot)-u_{N-1}^{\varepsilon}(t, \cdot)\right\|_{H_{\varepsilon}^{s}}=\mathcal{O}\left(\varepsilon^{N}\right), \quad \text { for } t \in\left[0, t_{*}\right],
$$

we can finally replace $u_{N}^{\varepsilon}(t, \cdot)$ by $u_{N-1}^{\varepsilon}(t, \cdot)$ in our stability result, which consequently proves the assertion of the theorem.

We want to stress that our stability result is an advancement when compared to the result of [7], in the sense that our asymptotic estimates no longer suffer from a loss of accuracy in powers of $\varepsilon$ (as has been the case in the cited work). We infer in particular, that our nonlinear setting allows for the same kind of stability result as one would expect in the linear case.
Theorem 4.5 also shows that no blow-up can occur in the exact solution $u^{\varepsilon}$ within the interval $t \in[0, T)$, determined by the existence time of (1.5). Hence, if the system (1.5) indeed admits global-in-time solutions we deduce that the solutions $u^{\varepsilon}$, starting close to such modulated pulses, exist for arbitrary long times (but blow-up may occur after the modulational structure is lost on longer time scales.)
Finally note that for $d=1$ the assumption $N>d / 2$ and the required closure of order $\Lambda=2 N+1$ imply that, at least, $N=1$ and thus $\Lambda=3$, whereas for $d=3$ spatial dimensions we need at least $N=2$ and hence $\Lambda=5$. This increase of the required $\Lambda$ for higher spatial dimensions $d$ can be relaxed though, as we shall show in the following section.

## 5 The case of higher-order resonances

In Section 3, we derived the approximate solution $u_{N}^{\varepsilon}$ under the assumption that $\mathcal{S}$ satisfies a closure condition of sufficient high order $\Lambda$, i.e. we required $\Lambda=2 N+1$ and $N>d / 2$. We shall show now that this closure condition can be significantly relaxed if we generalize the ansatz (3.1), (3.2) slightly by allowing for a larger set of modes $\mathcal{G}_{2 N+1}$. To this end we consider the new ansatz

$$
\begin{equation*}
u_{N}^{\varepsilon}(t, x)=\sum_{n=0}^{N} \sum_{\sigma \in \mathcal{G}_{2 n+1}} \varepsilon^{n} A_{n, \sigma}(t, x, y) \mathbf{E}_{\sigma}(\tau, y), \tag{5.1}
\end{equation*}
$$

where the set of modes $\mathcal{G}_{2 N+1}$ is then defined inductively as follows: Again we start from $\mathcal{S}=\left\{\left(k_{1}, \ell_{1}\right), \ldots,\left(k_{M}, \ell_{M}\right)\right\}$, which is now assumed to be only closed of order $\Lambda=3$ and we set $\mathcal{G}_{1}=\mathcal{G}_{\mathcal{S}}^{(1)}$. However, for $n \geq 2$ we will allow the sets $\mathcal{G}_{2 n+1}$ to be larger than $\mathcal{G}_{\mathcal{S}}^{(2 n+1)}$. The reason for this enlargement is motivated by the fact, that

$$
\widetilde{\mathcal{G}}_{\mathcal{S}}^{(2 n+1)}=\mathcal{G}_{\mathcal{S}}^{(2 n+1)} \cap\left(\mathcal{G} \backslash \mathcal{G}_{\mathcal{S}}^{(1)}\right)
$$

may be nonempty for some $n \geq 2$. Then, the corresponding equation (3.9) might not be solvable. To circumvent this problem, we set

$$
\mathcal{G}_{2 n-1}=\mathcal{G}_{\mathcal{S}}^{(2 n-1)} \cup \widetilde{\mathcal{G}}_{\mathcal{S}}^{(2 n+1)}
$$

At stage $n-1$ the corresponding $A_{n-1, \sigma}$ with $\sigma \in \widetilde{\mathcal{G}}_{\mathcal{S}}^{(2 n+1)}$ take the usual product form $a_{n-1, \sigma}(t, x) \chi_{\sigma}(y)$ with $\chi_{\sigma} \in \operatorname{ker} L_{0}^{\sigma} \subset H^{2}(\mathcal{Y})$ and free coefficients $a_{n-1, \sigma}$. Then, in step $n$ the corresponding equations for $\sigma \in \widetilde{\mathcal{G}}_{\mathcal{S}}^{(2 n+1)}$ can be solved as in Case II of Subsection 3.1, i.e. we obtain a linear transport equation for the free coefficients from the solvability condition for the previous step.

To make this more precise, we associate with the sequence $\mathcal{G}_{1}, \mathcal{G}_{3}, \ldots, \mathcal{G}_{2 N+1}$ another sequence $\check{\mathcal{G}}_{3}, \check{\mathcal{G}}_{5}, \ldots, \check{\mathcal{G}}_{2 N+1}$ as follows

$$
\begin{equation*}
\check{\mathcal{G}}_{2 n+1}=\bigcup_{\substack{n_{1}, n_{2}, n_{3}=0, \ldots, n-1 \\ n_{1}+n_{2}+n_{3}=n_{-1}}}\left(\mathcal{G}_{2 n_{1}+1}-\mathcal{G}_{2 n_{2}+1}+\mathcal{G}_{2 n_{3}+1}\right) . \tag{5.2}
\end{equation*}
$$

Hence, $\check{\mathcal{G}}_{2 n+1}$ contains all modes of order $\mathcal{O}\left(\varepsilon^{n}\right)$ that can be generated by the cubic nonlinearity $|u|^{2} u$ from the already given terms of lower order.

Definition 5.1. A mode set $\mathcal{S}=\left\{\left(k_{1}, \ell_{1}\right), \ldots,\left(k_{M}, \ell_{M}\right)\right\}$ is called weakly closed of order $2 N+1$, if there exists a sequence $\mathcal{G}_{1} \subset \mathcal{G}_{3} \subset \cdots \subset \mathcal{G}_{2 N+1} \subset \mathcal{B} \times \mathbb{R}$ such that the following conditions hold:
(i) $\mathcal{G}_{2 n+1}$ is finite for $n=0,1, \ldots, N$;
(ii) $\mathcal{G}_{1}=\mathcal{G}_{\mathcal{S}}^{(1)}=\Sigma(\mathcal{S})$;
(iii) $\check{\mathcal{G}}_{2 n+1} \subset \mathcal{G}_{2 n+1}$ for $n=1, \ldots, N$;
(iv) if $\sigma \in \mathcal{G}_{2 n+1} \cap \mathcal{G}$ with $n \in \mathbb{N}$, then either $\sigma \in \mathcal{G}_{2 n-1}$ or $\sigma \notin \check{\mathcal{G}}_{2 n+1}$;
(v) for each $\sigma \in \mathcal{G}_{2 N+1} \cap \mathcal{G}$ the group velocity $\vartheta_{\sigma}=\nabla_{k} E_{\ell}(k)$ exists,
where $\check{\mathcal{G}}_{3} \subset \cdots \subset \check{\mathcal{G}}_{2 N+1}$ is the sequence associated with $\mathcal{G}_{1}, \ldots, \mathcal{G}_{2 N+1}$ according to (5.2).

Condition (iv) means that any occurring resonant mode must either occur already in an earlier step and hence has a corresponding free coefficient or it appears for the first time but it is not yet generated by the nonlinear interaction. This will become clearer in the following examples.

Example 5.2. We first show that a system which is closed of order $2 N+1$ is also weakly closed of order $2 N+1$. For this we simply set $\mathcal{G}_{2 n+1}=\mathcal{G}_{\mathcal{S}}^{(2 n+1)}$. By construction, the associated sequence is $\check{\mathcal{G}}_{2 n+1}=\mathcal{G}_{2 n+1}$ and the conditions (i) to (iii) hold immediately. Moreover, (iv) follows since $\mathcal{G}_{2 n+1} \cap \mathcal{G}=\mathcal{G}_{1}=\mathcal{G}_{\mathcal{S}}^{(1)}$.

Example 5.3. We now show that a weak closure of order $\Lambda$ is in general a weaker condition than closure of order $\Lambda$ is. For this consider the case $\mathcal{S}=\left\{\left(k_{1}, \ell_{1}\right),\left(k_{2}, \ell_{2}\right)\right\}$ such that $\mathcal{G}_{1}=\mathcal{G}_{\mathcal{S}}^{(1)}=\left\{\sigma_{1}, \sigma_{2}\right\} \subset \mathcal{G}$ with $\sigma_{j}=\left(k_{j}, E_{\ell_{j}}\left(k_{j}\right)\right)$. An easy induction argument gives

$$
\mathcal{G}_{\mathcal{S}}^{(2 n+1)}=\left\{(n+1-j) \sigma_{1}+(j-n) \sigma_{2} ; j=0,1, \ldots, 2 n+1\right\}
$$

Now assume that $\mathcal{G} \cap \mathcal{G}_{\mathcal{S}}^{\left(2 N^{*}+1\right)} \subset\left\{\sigma_{1}, \sigma_{2}, 3 \sigma_{1}-2 \sigma_{2}\right\}$ for some $N^{*} \geq 2$, i.e. we have a closure of order 3 but not of order 5 .
We claim that a weak closure of order $2 N+1$ still holds for any $N$ such that $3 N \leq$ $2 N^{*}$. For $n=0, \ldots, N$ we let

$$
\mathcal{G}_{2 n+1}=\left\{j \sigma_{1}+(1-j) \sigma_{2} ; j=-\beta_{n}, \ldots, \alpha_{n}\right\},
$$

with $\alpha_{n}=[3(n+1) / 2]$ and $\beta_{n}=[3 n / 2]$, where $[\cdot]$ denotes the integer part. By construction of the sequences $\left(\alpha_{n}, \beta_{n}\right)_{n}$ we thus have $\mathcal{G}_{3}=\check{\mathcal{G}}_{3} \cup\left\{3 \sigma_{1}-2 \sigma_{2}\right\}$ and $\mathcal{G}_{2 n+1}=\check{\mathcal{G}}_{2 n+1}$ for $n=2, \ldots, N$. Hence conditions (i) to (iii) of Definition 5.1 are fulfilled. Condition (iv) also holds since $\mathcal{G}_{2 n+1} \cap \mathcal{G} \subset \mathcal{G}_{1} \cup\left\{3 \sigma_{1}-2 \sigma_{2}\right\}$.

Repeating the construction of Section 3.2 analogously we obtain the following result.
Lemma 5.4. Let $\mathcal{S}=\left\{\left(k_{1}, \ell_{1}\right), \ldots,\left(k_{M}, \ell_{M}\right)\right\}$ be a set of modes that is closed of order $\Lambda=3$ and weakly closed of order $\Lambda=2 N+1$ for some $N \in \mathbb{N}$. Then an approximation

$$
\begin{equation*}
u_{N}^{\varepsilon}(t, x)=\sum_{n=0}^{N} \varepsilon^{n} \sum_{\sigma \in \mathcal{G}_{2 n+1}} A_{n, \sigma}(t, x, y) \mathbf{E}_{\sigma} \tag{5.3}
\end{equation*}
$$

can be constructed as in Sections 3.1 and 3.2, where $\mathcal{G}_{1} \subset \cdots \mathcal{G}_{2 N+1}$ is the sequence guaranteed by Definition 5.1.

Hence, via the same steps as in Section 4 one can easily obtain a justification of the amplitude equations (1.5) under these relaxed resonance conditions.

Corollary 5.5. Let $d, N \in \mathbb{N}$ and $\mathcal{S}$ be a mode system that is closed of order $\Lambda=3$ and weakly closed of order $\Lambda=2 N+1$. Then, under the same assumptions as before the statement of Theorem 4.5 holds analogously with the approximate solution given by (5.3).

## A Hamiltonian structure and conservation laws of the amplitude equations

Recalling the amplitude equations in the general form (1.5), i.e.

$$
\begin{equation*}
\mathrm{i} \partial_{t} a_{m}=-\mathrm{i} \vartheta_{m} \cdot \nabla_{x} a_{m}+\sum_{\substack{p, q, r=1: \\ \Sigma\left(\mu_{p}, \mu_{q}, \mu_{r}\right)=\Sigma\left(\mu_{m}\right)}}^{M} \kappa_{(p, q, r, m)} a_{p} \bar{a}_{q} a_{r}, \tag{A.1}
\end{equation*}
$$

for $m=1, \ldots, M$, we want to highlight the Hamiltonian structure of this equation as well as to deduce the conserved quantities which follow from the specific structure of the resonances. For a general method on how to derive the reduced Hamiltonian structure and the first integrals for (A.1) we refer to [11, 12]. Here we just collect the results that are obtained with the general theory developed there.
We first note that (A.1) is a Hamiltonian system generated by

$$
\mathcal{H}^{\mathrm{red}}(a)=\int_{\mathbb{R}^{d}} \sum_{m=1}^{M} \operatorname{Im}\left(\bar{a}_{m} \vartheta_{m} \cdot \nabla a_{m}\right)+\sum_{\substack{p, q, r=1: \\ \Sigma\left(\mu_{p}, q_{q}, \mu_{r}\right)=\Sigma\left(\mu_{m}\right)}}^{M} \frac{\kappa_{(p, q, r, m)}}{2} a_{p} \bar{a}_{q} a_{r} \bar{a}_{m} \mathrm{~d} x,
$$

and the symplectic two-form $i=\sqrt{-1}$. Hence, (A.1) takes the form

$$
\mathrm{i} \partial_{t} a_{m}=\delta_{\bar{a}_{m}} \mathcal{H}^{\mathrm{red}}(a) .
$$

Note that the reduced Hamiltonian $\mathcal{H}^{\text {red }}(a)$ is not obtained as the lowest order expansion of the original Hamiltonian $\mathcal{H}^{\varepsilon}\left(u^{\varepsilon}\right)$ of the full system, which reads

$$
\mathcal{H}^{\varepsilon}\left(u^{\varepsilon}\right)=\int_{\mathbb{R}^{d}} \frac{\varepsilon^{2}}{2}|\nabla u|^{2}+V_{\Gamma}\left(\frac{x}{\varepsilon}\right)\left|u^{\varepsilon}\right|^{2}+\frac{\varepsilon \kappa}{2}|u|^{4} \mathrm{~d} x .
$$

Indeed, inserting the ansatz $u^{\varepsilon}=u_{N}^{\varepsilon}$, as given in (4.1), into $\mathcal{H}^{\varepsilon}\left(u^{\varepsilon}\right)$ we find, as $\varepsilon \rightarrow 0$, that

$$
\mathcal{H}^{\varepsilon}\left(u_{N}^{\varepsilon}\right)=\mathcal{I}(a)+\mathcal{O}(\varepsilon), \quad \text { with } \mathcal{I}(a)=\int_{\mathbb{R}^{d}} \sum_{m=1}^{M} E_{\ell_{m}}\left(k_{m}\right)\left|a_{m}\right|^{2} \mathrm{~d} x .
$$

Even though not a Hamiltonian, $\mathcal{I}(a)$ clearly is a conserved quantity for the amplitude system (A.1).
However, the energy levels $E_{\ell_{m}}\left(k_{m}\right) \equiv \omega_{m}$ do not occur explicitly in the amplitude system (A.1). Hence, any choice of $\widetilde{\omega}_{m} \in \mathbb{R}$ that is compatible with the resonance conditions (2.4) leads to additional conserved quantities. More precisely, if $\widetilde{\omega}_{1}, \ldots, \widetilde{\omega}_{M}$ are chosen such that for any $p, q, r, m \in\{1, \ldots, M\}$ the resonance conditions (2.4) imply the identity

$$
\widetilde{\omega}_{p}-\widetilde{\omega}_{q}+\widetilde{\omega}_{r}=\widetilde{\omega}_{m},
$$

then

$$
\widetilde{\mathcal{I}}(a)=\int_{\mathbb{R}^{d}} \sum_{m=1}^{M} \widetilde{\omega}_{m}\left|a_{m}\right|^{2} \mathrm{~d} x
$$

defines a first integral.
A particular choice is $\widetilde{\omega}_{1}=\cdots=\widetilde{\omega}_{M}=1$, which leads to the trivial fact that the $L^{2}$ norm is preserved. The latter can of course also be deduced from the fact that the $L^{2}$ norm was preserved in the original problem for $u^{\varepsilon}$ or from the fact that the mode system is invariant under the phase shifts $\left(a_{1}, \ldots, a_{M}\right) \mapsto\left(\mathrm{e}^{\mathrm{i} \alpha} a_{1}, \ldots, \mathrm{e}^{\mathrm{i} \alpha} a_{M}\right)$, with $\alpha \in \mathbb{R}$. Similarly, $\widetilde{\mathcal{I}}(a)$ can be understood as a first integral with respect to suitably chosen phase shifts that are compatible with the resonance structure, i.e.

$$
\widetilde{T}_{\alpha}:\left(a_{1}, \ldots, a_{M}\right) \mapsto\left(\mathrm{e}^{\mathrm{i} \widetilde{\omega}_{1} \alpha} a_{1}, \ldots, \mathrm{e}^{\mathrm{i} \widetilde{\omega}_{M} \alpha} a_{M}\right), .
$$

Finally, it should be mentioned that the mode system (A.1) is translation invariant. This provides first integrals associated with the translation operators in the coordinate directions

$$
\mathcal{I}_{\theta}^{\mathrm{trans}}(a)=\int_{\mathbb{R}^{d}} \sum_{m=1}^{M} \operatorname{Im}\left(\bar{a}_{m} \theta \cdot \nabla a_{m}\right) \mathrm{d} x, \quad \theta \in \mathbb{R}^{d}
$$

So far, we are not able to show that these conserved quantities are enough to provide a global existence result. Note in particular, that for the original problem the case $\kappa>0$ leads to an energy that is definite and allows us to conclude global existence for the nonlinear Schrödinger equation (1.1). For the mode system (A.1) the sign of $\kappa$ is no longer helpful, since the term involving derivatives is indefinite. Note also that global existence for the mode system cannot be inferred from global existence of (1.1), since we can not expect the solutions $u^{\varepsilon}$ to remain in the form of modulated pulses for all times.

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