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WELL-POSEDNESS AND PARABOLIC SMOOTHING EFFECT FOR HIGHER ORDER SCHRÖDINGER TYPE EQUATIONS WITH CONSTANT COEFFICIENTS

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Abstract

In this paper, we consider the Cauchy problem of a class of higher order Schrödinger type equations with constant coefficients. By employing the energy inequality, we show the L^2 well-posedness, the parabolic smoothing and a breakdown of the persistence of regularity. We classify this class of equations into three types on the basis of their smoothing property.

1. Introduction

In this paper, we consider the Cauchy problem of the following:

$$(1.1) \quad D_t u(t, x) = D_x^{2m} u(t, x) + \sum_{j=1}^{2m} (a_j D_x^{2m-j} u(t, x) + b_j D_x^{2m-j} \bar{u}(t, x)),$$

$$(1.2) \quad u(0, x) = \varphi(x),$$

where $1 \leq m \in \mathbb{N}$, $\mathcal{M} = \mathbb{R}$ (or \mathbb{T}), $(t, x) \in (-\infty, \infty) \times \mathcal{M}$, $D_t = -i\partial_t$, $D_x = -i\partial_x$ and i is the imaginary unit. The constants $\{a_j\}, \{b_j\} \subset \mathbb{C}$ and the initial data $\varphi(x) : \mathcal{M} \rightarrow \mathbb{C}$ are given and $u(t, x) : (-\infty, \infty) \times \mathcal{M} \rightarrow \mathbb{C}$ is unknown. We are interested in the Cauchy problem of the following higher order nonlinear Schrödinger type equations:

$$(1.3) \quad i\partial_t u(t, x) - \partial_x^{2m} u(t, x) = F(\partial_x^{2m-1} u, \partial_x^{2m-1} \bar{u}, \partial_x^{2m-2} u, \partial_x^{2m-2} \bar{u}, \dots, u, \bar{u}),$$

with (1.2), where F is a polynomial. As important examples, this class of equations includes the nonlinear Schrödinger hierarchy and the derivative nonlinear Schrödinger hierarchy, which are integrable systems appearing in the soliton theory. It is known that the Cauchy problem of (1.3) with (1.2) is locally well-posed on \mathbb{R} in weighted Sobolev spaces (of which functions are also sufficiently smooth). Its proof is based on the Kato type smoothing estimate and the gauge transform [5, 6]. See Section 3 in [11] for this argument. On the one hand, the well-posedness for (1.3) with (1.2) on \mathbb{R} without any weight or on \mathbb{T} remains open. We also refer to [4, 10, 11, 16] for well-posedness results to higher order dispersive equations including the KdV hierarchy. In [2], Chihara studied the well-posedness and the ill-posedness of (1.3) for $m = 1$ with (1.2) on \mathbb{T} . Recently, in [20], the second author has studied a similar problem and shown a non-existence result of solutions of (1.3) for some nonlinearity and $m = 1$ with (1.2) on \mathbb{T} by employing a smoothing for elliptic equations.

In other words, even when we restrict (1.3) to $m = 1$, the well-posedness for (1.3) with (1.2) on \mathbb{T} is remarkably different from that on \mathbb{R} . Therefore, the nonlinearity F must have special structures (expected to include the case where (1.3) is an integrable system) when the Cauchy problem of (1.3) with (1.2) is (locally) well-posed on \mathbb{T} . In proofs of [2, 20], the so called “energy inequality” of (1.1) with variable coefficients $\{a_j(t, x)\}$ and $\{b_j(t, x)\}$ plays an important role. Our plan is to extend this result to $m \geq 2$. However, the energy inequality for higher m is much complicated. Therefore, we assume $\{a_j\}$ and $\{b_j\}$ are constants to make the problem simple in the present paper and will study the variable coefficients case in the forthcoming paper. λ defined below is used to classify (1.1) into three types.

DEFINITION 1. We write $\sum_{k=1}^0 c_k = 0$ for any sequence $\{c_k\}$. $\gamma = \{\gamma_j\}_{j=1}^{m-1}$ and $\lambda = \{\lambda_j\}_{j=1}^{2m-1}$ are defined as

$$\begin{aligned} \gamma_j &= b_{2j} - \sum_{k=1}^{j-1} \bar{a}_{2(j-k)} \gamma_k, \quad 1 \leq j \leq m-1, \\ \begin{cases} \lambda_{2j} = 2 \operatorname{Im} a_{2j} - 2 \sum_{k=1}^{j-1} \operatorname{Im} \bar{b}_{2(j-k)} \gamma_k, & 1 \leq j \leq m-1, \\ \lambda_{2j-1} = 2 \operatorname{Im} a_{2j-1} + 2 \sum_{k=1}^{j-1} \operatorname{Im} \bar{b}_{2(j-k)-1} \gamma_k, & 1 \leq j \leq m. \end{cases} \end{aligned}$$

Our main result is the following. We write $P^+ f(x) := \mathcal{F}^{-1}(\chi(\xi \geq 1) \mathcal{F} f)(x)$ and $P^- f(x) := \mathcal{F}^{-1}(\chi(\xi \leq -1) \mathcal{F} f)(x)$, where \mathcal{F} is the Fourier transform and χ is the definition function.

Theorem 1.1.

(Dispersive type, L^2 well-posedness) Assume that $\lambda_j = 0$ for $1 \leq j \leq 2m-1$. Then, for any $\varphi \in L^2(\mathcal{M})$, there exists a unique solution $u(t, x)$ of (1.1)–(1.2) such that $u(t, x) \in C((-\infty, \infty); L^2(\mathcal{M}))$.

(Parabolic type) Assume that there exists $j^* \in \mathbb{N}$ such that $\lambda_j = 0$ for $1 \leq j < 2j^*$ and $\lambda_{2j^*} > 0$ (resp. $\lambda_{2j^*} < 0$). (i) For any $\varphi \in L^2(\mathcal{M})$, there exist a unique solution $u(t, x)$ of (1.1)–(1.2) on $[0, \infty)$ (resp. $(-\infty, 0]$) such that $u(t, x) \in C([0, \infty); L^2(\mathcal{M})) \cap C^\infty((0, \infty) \times \mathcal{M})$ (resp. $C((-\infty, 0]; L^2(\mathcal{M})) \cap C^\infty((-\infty, 0) \times \mathcal{M})$). (ii) For any $\varphi \in L^2(\mathcal{M}) \setminus C^\infty(\mathcal{M})$ and $\delta > 0$, no solution u of (1.1)–(1.2) exists on $(-\delta, 0]$ (resp. $[0, \delta)$) such that $u(t, x) \in C((-\delta, 0]; L^2(\mathcal{M}))$ (resp. $C([0, \delta); L^2(\mathcal{M}))$).

(Elliptic type) Assume that there exists $j^* \in \mathbb{N}$ such that $\lambda_j = 0$ for $1 \leq j < 2j^* - 1$ and $\lambda_{2j^*-1} > 0$ (resp. $\lambda_{2j^*-1} < 0$). (i) Let $\varphi \in L^2(\mathcal{M})$ satisfy $P^+ \varphi \notin H^{1/2}(\mathcal{M})$. Then, for any $\delta > 0$, there exist no solution $u(t, x)$ of (1.1)–(1.2) on $[-\delta, 0]$ (resp. $[0, \delta]$) satisfying $u \in C([-\delta, 0]; L^2(\mathcal{M}))$ (resp. $u \in C([0, \delta]; L^2(\mathcal{M}))$). Moreover, the same result as above holds even if we replace P^+ , $[-\delta, 0]$ and $[0, \delta]$ with P^- , $[0, \delta]$ and $[-\delta, 0]$, respectively. (ii) Let $\varphi \in L^2(\mathcal{M}) \setminus C^\infty(\mathcal{M})$. Then, for any $\delta > 0$, there exist no solution $u(t, x)$ of (1.1)–(1.2) on $[-\delta, \delta]$ satisfying $u \in C([-\delta, \delta]; L^2(\mathcal{M}))$.

REMARK 1.1. Put $v(t) = \langle \partial_x \rangle^{-s} u(t)$. Then v satisfies (1.1) if u is the solution of (1.1) and $u(t) \in L^2(\mathcal{M}) \Leftrightarrow v(t) \in H^s(\mathcal{M})$. Therefore, Theorem 1.1 holds even if we replace $L^2(\mathcal{M})$ with $H^s(\mathcal{M})$ and $H^{1/2}(\mathcal{M})$ with $H^{s+1/2}(\mathcal{M})$ for any $s \in \mathbb{R}$.

REMARK 1.2. In “Dispersive type,” the persistence of regularity holds on both $(-\infty, 0]$ and $[0, \infty)$. In “Parabolic type,” the equations have the parabolic smoothing effect on either $(-\infty, 0]$ or $[0, \infty)$, which means the persistence of regularity breaks down on either $[0, \infty)$ or $(-\infty, 0]$. Non-existence results in “Parabolic type” and “Elliptic type” is by the breakdown of the persistence of regularity.

REMARK 1.3. We give some examples of $\{a_j\}$ and $\{b_j\}$.

- When $m = 1$, “Parabolic type” does not occur. In fact, the equation $D_t u = D_x^2 u + a_1 D_x u + a_2 u + b_1 D_x \bar{u} + b_2 \bar{u}$ is “Dispersive type” if $\text{Im } a_1 = 0$ and it is “Elliptic type” otherwise.
- Let $H(u)$ be a quadratic form defined by

$$H(u) := \frac{1}{2} \int |\partial_x^m u|^2 + \sum_{j=1}^{2m} (c_j u \partial_x^{2m-j} u + d_j u \partial_x^{2m-j} \bar{u} + e_j \bar{u} \partial_x^{2m-j} \bar{u}) dx$$

for given $\{c_j\}, \{d_j\}, \{e_j\} \subset \mathbb{C}$. Then, it is easy to check that $H(u)$ is the Hamiltonian of the equation (1.1) if and only if $\text{Im } a_j = b_{2n-1} = \text{Re } d_{2n-1} = \text{Im } d_{2n} = 0$ and $c_{2n} = \bar{e}_{2n}$ for $1 \leq j \leq 2m-1$ and $1 \leq n \leq m$. In particular, we can write $c_{2j} = (-1)^{m-j} \bar{b}_{2j}/2$, $d_{2j-1} = i(-1)^{m-j} a_{2j-1}$ and $d_{2j} = (-1)^{m-j} a_{2j}$ for $1 \leq j \leq m$ (without loss of generality we can assume $c_{2j-1} = e_{2j-1} = 0$ since c_{2j-1} - and e_{2j-1} -terms always vanish by the integration by parts). In this case, we see from Definition 1 and Remark 2.2 that $\lambda_j = 0$ for $1 \leq j \leq 2m-1$, which implies that Hamiltonian equations are “Dispersive type.”

- By using the equation (1.1), we have

$$\begin{aligned} \frac{d}{dt} \|u(t)\|^2 &= 2 \text{Re } i \langle D_t u, u \rangle \\ &= -2 \sum_{j=1}^{2m} (\text{Im } a_j) \langle D_x^{2m-j} u, u \rangle - 2 \sum_{n=1}^m \text{Im } b_{2n} \langle D_x^{2(m-n)} \bar{u}, u \rangle. \end{aligned}$$

Therefore, when $\text{Im } a_j = b_{2n} = 0$ for $1 \leq j \leq 2m-1$ and $1 \leq n \leq m$, the solution of the equation (1.1) conserves the mass, i.e., $\|u(t)\|$. We see from the scaling argument that this condition is also necessary. In this case, the equation (1.1) is “Dispersive type.” Indeed, it is easy to see $\gamma_j = 0$ for $1 \leq j \leq m-1$ by Definition 1. It then follows that $\lambda_k = 2 \text{Im } a_k = 0$ for $1 \leq k \leq 2m-1$.

- When $m = 2$, we have

$$\lambda_1 = 2 \text{Im } a_1, \quad \lambda_2 = 2 \text{Im } a_2, \quad \lambda_3 = 2 \text{Im } a_3 + 2 \text{Im } \bar{b}_1 b_2.$$

So, equations $D_t u = D_x^4 u + i D_x u$ and $D_t u = D_x^4 u + D_x^3 \bar{u} - i D_x^2 \bar{u}$ are “Elliptic type.” On the other hand, $D_t u = D_x^4 u + i D_x u + D_x^3 \bar{u} - i D_x^2 \bar{u}$ is “Dispersive type” although this equation does not have the Hamiltonian.

We recall several results for equations related to (1.1). There is a large literature on the well-posedness for the Cauchy problem of Schrödinger type equations, especially $m = 1$. In [13], Mizohata showed that if the Cauchy problem

$$\begin{aligned}\partial_t u &= \sum_{j=1}^n (i\partial_j^2 + c_j(x)\partial_j)u + f(t, x), \quad (t, x) \in \mathbb{R} \times \mathbb{R}^n, \\ u(0, x) &= \varphi(x)\end{aligned}$$

is L^2 well-posed, then the condition

$$\sup_{(t, \omega, x) \in \mathbb{R} \times S^{n-1} \times \mathbb{R}^n} \left| \operatorname{Im} \int_0^t \sum_{j=1}^n c_j(x + s\omega) \omega_j ds \right| < \infty$$

holds. In particular, this condition is also sufficient for the L^2 well-posedness when $n = 1$. See [1, 2, 3, 8, 9, 14, 17, 18] (and references therein) for related results. For $m = 2$, in [15], Mizuhara studied L^2 well-posedness for the Cauchy problem:

$$(1.4) \quad (D_t - D_x^4 - c_1(x)D_x^3 - c_2(x)D_x^2 - c_3(x)D_x - c_4(x))u = f(t, x)$$

$$(1.5) \quad u(0, x) = \varphi(x),$$

where $(t, x) \in \mathbb{R} \times \mathcal{M}$. To be precise, he also studied another equation of the KdV type. When $\mathcal{M} = \mathbb{T}$, he deduced the necessary and sufficient conditions for the L^2 well-posedness for (1.4)–(1.5). On the other hand, when $\mathcal{M} = \mathbb{R}$, he showed some conditions for the L^2 well-posedness. Indeed, his sufficient condition for the L^2 well-posedness is also necessary under the additional assumption. In [19], Tarama removed Mizuhara's additional assumption, so he obtained the necessary and sufficient conditions for the L^2 well-posedness for (1.4)–(1.5) on \mathbb{R} .

Since the coefficients are constants, by the Fourier transform, (1.1) can be rewritten into the following:

$$(1.6) \quad D_t \widehat{u}(t, \xi) = \xi^{2m} \widehat{u}(t, \xi) + \sum_{j=1}^{2m} (a_j \xi^{2m-j} \widehat{u}(t, \xi) + b_j \xi^{2m-j} \widehat{u}(t, -\xi)).$$

Here, we fix $\xi \in \mathbb{R}$ (or \mathbb{Z}) and put

$$U_\xi(t) = \begin{pmatrix} \widehat{u}(t, \xi) \\ \widehat{\bar{u}}(t, -\xi) \end{pmatrix}, \quad X_0 = \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix}, \quad X_j = \begin{pmatrix} a_j & b_j \\ (-1)^{j+1} \overline{b_j} & (-1)^{j+1} \overline{a_j} \end{pmatrix},$$

for $1 \leq j \leq 2m$. Then, by (1.6) with (1.2), it follows that

$$(1.7) \quad D_t U_\xi(t) = \sum_{j=0}^{2m} \xi^{2m-j} X_j U_\xi(t), \quad U_\xi(0) = {}^t(\widehat{\varphi}(\xi), \widehat{\bar{\varphi}}(-\xi)),$$

which is a system of linear ordinary differential equations. We can easily obtain the unique solution

$$(1.8) \quad U_\xi(t) = U_\xi(0) \exp it \sum_{j=0}^{2m} \xi^{2m-j} X_j$$

on $t \in (-\infty, \infty)$ for each $\xi \in \mathbb{R}$ (or \mathbb{Z}). Therefore, our interest in Theorem 1.1 is essentially on the regularity of the solution. Here, note that $X_j X_k = X_k X_j$ holds for any $0 \leq j, k \leq 2m$ if and only if $b_j = 0$ holds for any $1 \leq j \leq 2m$. If we assume this assumption, (1.7) is not a system but a single ordinary differential equation and

$$(1.9) \quad \widehat{u}(t, \xi) = \widehat{\varphi}(\xi) \exp it \left(\xi^{2m} + \sum_{j=1}^{2m} \xi^{2m-j} a_j \right)$$

for each $\xi \in \mathbb{R}$ (or \mathbb{Z}). Since $\gamma_j = 0$ and $\lambda_j = 2\operatorname{Im} a_j$, it follows that

$$|\widehat{u}(t, \xi)| = |\widehat{\varphi}(\xi)| \prod_{j=1}^{2m} \exp \frac{-t \xi^{2m-j} \lambda_j}{2},$$

by which we obtain Theorem 1.1 easily. On the other hand, it seems difficult to obtain Theorem 1.1 by (1.8) for general $\{b_j\}$ since $X_j X_k \neq X_k X_j$ for some j, k . To avoid this difficulty, we employ the energy estimate. In particular, we modify the energy, adding correction terms so as to cancel out derivative losses. See e.g. [7, 12] for this argument. However, some of derivative losses cannot be eliminated by correction terms, and they may essentially affect the well-posedness of the Cauchy problem of (1.1)–(1.2) as stated in Theorem 1.1. Propositions 2.1 and 2.2 are main estimates in this paper. The first term of the left-hand side of (2.1) is the main part of the energy. The second term is the correction term. For “Dispersive type,” the third and the fourth terms vanish. Thus, we easily obtain the L^2 a priori estimate. For “Parabolic type,” the third term includes $\lambda_{2j^*} \|\partial_x^{m-j^*} u\|^2$. The parabolic smoothing is caused by the term. For “Elliptic type,” the fourth term includes $\lambda_{2j^*-1} \langle D_x^{2(m-j^*)+1} u, u \rangle$. We want to show the parabolic smoothing by making use of the term. However, the sign of the term is not definite. That is unfavorable in our argument. Therefore, we compute the energy inequalities of $P^+ u$ and $P^- u$ instead of u and obtain Lemma 2.3. Note that the sign of all terms except the correction terms in (2.4) and (2.5) are definite. Though (2.4) is the energy inequality for $\|P^+ u\|$, it includes $\lambda_j^- \|\partial_x^{m-j/2} P^- u\|^2$. This is because (1.1) is essentially coupled system of $P^+ u$ and $P^- u$ as (1.6). The term $\lambda_j^- \|\partial_x^{m-j/2} P^- u\|^2$ cannot be estimated by $\|u\|$. This is the main difficulty in the proof of “Elliptic type” in Theorem 1.1. The key idea is to eliminate these terms in two steps where $1 \leq j \leq 2j^* + 1$ and $2(j^* + 1) \leq j \leq 2m - 1$. First we analyse a property of $\{\lambda_j^-\}$ so that $\lambda_{2k} = 0$ for $1 \leq k \leq j^* - 1$ implies that the first $2j^* + 1$ of $\{\lambda_j^-\}$ vanish (see Lemma 2.5). In order to cancel out the rest of unfavorable terms $\lambda_j^- \|\partial_x^{m-j/2} P^- u\|^2$ for $2(j^* + 1) \leq j \leq 2m - 1$, we use an additional correction term F_k^- and obtain (2.2) (see also (2.3)). Here, F_k^- originates from the energy inequality for $\|\partial_x^{-(k+2)/2} P^- u\|$, and F_k^- does not yield a bad effect thanks to the first step.

The rest of this paper is organized as follows. In Section 2, we state main estimates which are energy estimates for u and $P^\pm u$, and give proofs of them. In Section 3, we show Theorem 1.1. In particular, we show a smoothing for “Elliptic type” (Proposition 3.2) from the energy estimate for $P^\pm u$.

Here, we set some notation. Let $\langle \cdot, \cdot \rangle := \langle \cdot, \cdot \rangle_{L^2}$ and $\|\cdot\| := \|\cdot\|_{L^2}$. We also use the same symbol for $\langle \cdot \rangle := (1 + |\cdot|^2)^{1/2}$. P_0 and $P_{\neq 0}$ are defined by $P_0 f(x) := \mathcal{F}^{-1}(\chi(|\xi| < 1) \mathcal{F} f)(x)$ and $P_{\neq 0} f(x) := \mathcal{F}^{-1}(\chi(|\xi| \geq 1) \mathcal{F} f)(x)$. We define the Riesz and Bessel potentials by $|\partial_x|^s f := \mathcal{F}^{-1}(|\xi|^s \mathcal{F} f)(x)$ and $\langle \partial_x \rangle^s f := \mathcal{F}^{-1}(\langle \xi \rangle^s \mathcal{F} f)(x)$.

2. the energy estimates

Our purpose in this section is to show Propositions 2.1 and 2.2. Proposition 2.1 below is used to show “Dispersive type” and “Parabolic type” in Theorem 1.1.

Proposition 2.1. *Let u satisfy (1.1). Then, there exists $C = C(\{a_j\}, \{b_j\}) > 0$ such that*

$$(2.1) \quad \left| \frac{d}{dt} \left(\|u\|^2 + \sum_{j=1}^{m-1} \operatorname{Re} \gamma_j \langle D_x^{-2j} P_{\neq 0} \bar{u}, P_{\neq 0} u \rangle \right) + \sum_{j=1}^{m-1} \lambda_{2j} \| |\partial_x|^{m-j} u \|^2 + \sum_{j=1}^m \lambda_{2j-1} \langle D_x^{2(m-j)+1} u, u \rangle \right| \leq C \|u\|^2.$$

DEFINITION 2. $\alpha = \{\alpha_j\}_{j=1}^{2m-1}$, $\lambda^+ = \{\lambda_j^+\}_{j=1}^{2m-1}$, $\lambda^- = \{\lambda_j^-\}_{j=1}^{2m-1}$ are defined as

$$\begin{aligned} \alpha_j &= b_j - \frac{1}{2} \sum_{k=1}^{j-1} (1 + (-1)^{j-k}) \bar{a}_{j-k} \alpha_k, \quad 1 \leq j \leq 2m-1, \\ \lambda_j^+ &= 2 \operatorname{Im} a_j + \sum_{k=1}^{j-1} (-1)^{j-k+1} \operatorname{Im} \bar{b}_{j-k} \alpha_k, \quad 1 \leq j \leq 2m-1, \\ \lambda_j^- &= - \sum_{k=1}^{j-1} \operatorname{Im} \bar{b}_{j-k} \alpha_k, \quad 1 \leq j \leq 2m-1. \end{aligned}$$

Let $1 \leq j^* \leq m-2$. Assume that $\lambda_{2j^*-1}^+ \neq 0$. $\beta^+ = \{\beta_k^+\}_{k=1}^{2(m-j^*-1)}$ and $\beta^- = \{\beta_k^-\}_{k=1}^{2(m-j^*-1)}$ are inductively defined as

$$\begin{aligned} \lambda_{2j^*+k+1}^- &= \sum_{j=1}^k (-1)^{k-j} \lambda_{2j^*+k-j-1}^+ \beta_j^+, \quad 1 \leq k \leq 2(m-j^*-1), \\ \lambda_{2j^*+k+1}^- &= \sum_{j=1}^k (-1)^k \lambda_{2j^*+k-j-1}^+ \beta_j^-, \quad 1 \leq k \leq 2(m-j^*-1). \end{aligned}$$

Note that $\lambda_{2j^*+k-j-1}^+ = \lambda_{2j^*-1}^+ \neq 0$ when $j = k$. So, β_k^+ and β_k^- are well-defined.

REMARK 2.1. It is easy to see that $\gamma_j = \alpha_{2j}$ for $1 \leq j \leq m-1$. Thus, we have

$$\lambda_{2j} = \lambda_{2j}^+ + \lambda_{2j}^-, \quad \lambda_{2k-1} = \lambda_{2k-1}^+ - \lambda_{2k-1}^-$$

for $1 \leq j \leq m-1$ and $1 \leq k \leq m$.

Proposition 2.2 below is used to show “Elliptic type” in Theorem 1.1.

Proposition 2.2. *Let u satisfy (1.1). Assume that there exists $j^* \in \mathbb{N}$ such that $\lambda_j = 0$ for $1 \leq j \leq 2(j^*-1)$ and $\lambda_{2j^*-1} \neq 0$. Put*

$$\begin{aligned} F_k^-(u) &= \| |\partial_x|^{-(k+2)/2} P^- u \|^2 + \sum_{j=1}^{2m-1} \operatorname{Re} \alpha_j \langle D_x^{-j} |\partial_x|^{-k-2} \overline{P^+ u}, P^- u \rangle, \\ F_k^+(u) &= \| |\partial_x|^{-(k+2)/2} P^+ u \|^2 + \sum_{j=1}^{2m-1} \operatorname{Re} \alpha_j \langle D_x^{-j} |\partial_x|^{-k-2} \overline{P^- u}, P^+ u \rangle. \end{aligned}$$

Then, there exists $C = C(\{a_j\}, \{b_j\}) > 0$ such that

$$(2.2) \quad \left| \frac{d}{dt} \left(\|P^+ u\|^2 + \sum_{j=1}^{2m-1} \operatorname{Re} \alpha_j \langle D_x^{-j} \overline{P^- u}, P^+ u \rangle + \sum_{k=1}^{2(m-j^*-1)} \beta_k^+ F_k^-(u) \right) \right. \\ \left. + \lambda_{2j^*-1}^+ \|\partial_x\|^{m-j^*+1/2} P^+ u \|^2 \right| \leq C \|u\|^2 + C \|\partial_x\|^{m-j^*} P^+ u \|^2,$$

and

$$(2.3) \quad \left| \frac{d}{dt} \left(\|P^- u\|^2 + \sum_{j=1}^{2m-1} \operatorname{Re} \alpha_j \langle D_x^{-j} \overline{P^+ u}, P^- u \rangle + \sum_{k=1}^{2(m-j^*-1)} \beta_k^- F_k^+(u) \right) \right. \\ \left. - \lambda_{2j^*-1}^+ \|\partial_x\|^{m-j^*+1/2} P^- u \|^2 \right| \leq C \|u\|^2 + C \|\partial_x\|^{m-j^*} P^- u \|^2.$$

To prove Propositions 2.1 and 2.2, we use the following lemma.

Lemma 2.3. *Let u satisfy (1.1). Then, there exists $C = C(\{a_j\}, \{b_j\}) > 0$ such that*

$$(2.4) \quad \left| \frac{d}{dt} \left(\|P^+ u\|^2 + \sum_{j=1}^{2m-1} \operatorname{Re} \alpha_j \langle D_x^{-j} \overline{P^- u}, P^+ u \rangle \right) \right. \\ \left. + \sum_{j=1}^{2m-1} (\lambda_j^+ \|\partial_x\|^{m-j/2} P^+ u \|^2 + \lambda_j^- \|\partial_x\|^{m-j/2} P^- u \|^2) \right| \leq C \|u\|^2$$

and

$$(2.5) \quad \left| \frac{d}{dt} \left(\|P^- u\|^2 + \sum_{j=1}^{2m-1} \operatorname{Re} \alpha_j \langle D_x^{-j} \overline{P^+ u}, P^- u \rangle \right) \right. \\ \left. + \sum_{j=1}^{2m-1} (-1)^j (\lambda_j^+ \|\partial_x\|^{m-j/2} P^- u \|^2 + \lambda_j^- \|\partial_x\|^{m-j/2} P^+ u \|^2) \right| \leq C \|u\|^2.$$

Proof of Lemma 2.3. First, we show (2.4). For simplicity, we set $v^+ := P^+ u$ and $v^- := P^- u$. Note that $P^+ \bar{u} = \overline{P^- u} = \bar{v}^-$ and $P^- \bar{u} = \overline{P^+ u} = \bar{v}^+$. Then, v^+ and v^- satisfy

$$(2.6) \quad D_t v^+ = D_x^{2m} v^+ + \sum_{k=1}^{2m} (a_k D_x^{2m-k} v^+ + b_k D_x^{2m-k} \bar{v}^-)$$

and

$$(2.7) \quad D_t \bar{v}^- = -D_x^{2m} \bar{v}^- - \sum_{k=1}^{2m} (-1)^k (\bar{a}_k D_x^{2m-k} \bar{v}^- + \bar{b}_k D_x^{2m-k} v^+).$$

By (2.6), we have

$$\begin{aligned} \frac{d}{dt} \|v^+\|^2 &= 2 \operatorname{Re} \langle \partial_t v^+, v^+ \rangle = -2 \operatorname{Im} \langle D_t v^+, v^+ \rangle \\ &= -2 \sum_{j=1}^{2m} (\operatorname{Im} a_j \langle D_x^{2m-j} v^+, v^+ \rangle + \operatorname{Im} b_j \langle D_x^{2m-j} \bar{v}^-, v^+ \rangle) \\ &= -2 \sum_{j=1}^{2m} (\operatorname{Im} a_j \|\partial_x\|^{m-j/2} P^+ u \|^2 + \operatorname{Im} b_j \langle D_x^{2m-j} \bar{v}^-, v^+ \rangle). \end{aligned}$$

Here, we consider the time derivative of correction terms to cancel out the second term. Fix $1 \leq j \leq 2m-1$. We see from (2.6) and (2.7) that

$$\begin{aligned}
 \frac{d}{dt} \operatorname{Re} \alpha_j \langle D_x^{-j} \overline{P^- u}, v^+ \rangle &= -\operatorname{Im} \alpha_j \langle D_x^{-j} D_t \overline{v^-}, v^+ \rangle + \operatorname{Im} \alpha_j \langle D_x^{-j} \overline{v^-}, D_t v^+ \rangle \\
 &= \operatorname{Im} \alpha_j \langle D_x^{-j} (D_x^{2m} \overline{v^-}), v^+ \rangle + \operatorname{Im} \alpha_j \langle D_x^{-j} \overline{v^-}, D_x^{2m} v^+ \rangle \\
 &\quad + \sum_{k=1}^{2m} ((-1)^k \operatorname{Im} \alpha_j \bar{a}_k \langle D_x^{2m-k-j} \overline{v^-}, v^+ \rangle + (-1)^k \operatorname{Im} \alpha_j \bar{b}_k \langle D_x^{2m-k-j} v^+, v^+ \rangle) \\
 &\quad + \operatorname{Im} \alpha_j \bar{a}_k \langle D_x^{2m-k-j} \overline{v^-}, v^+ \rangle + \operatorname{Im} \alpha_j \bar{b}_k \langle D_x^{2m-k-j} \overline{v^-}, \overline{v^-} \rangle) \\
 &=: A_1^j + B_1^j + \sum_{k=1}^{2m} (A_{2,k}^j + A_{3,k}^j + B_{2,k}^j + B_{3,k}^j).
 \end{aligned}$$

Observe that

$$\begin{aligned}
 A_1^j + B_1^j &= 2 \operatorname{Im} \alpha_j \langle D_x^{2m-j} \overline{v^-}, v^+ \rangle, \\
 A_{2,k}^j + B_{2,k}^j &= (1 + (-1)^k) \operatorname{Im} \alpha_j \bar{a}_k \langle D_x^{2m-k-j} \overline{v^-}, v^+ \rangle, \\
 A_{3,k}^j &= (-1)^k \operatorname{Im} \alpha_j \bar{b}_k |||\partial_x|^{m-(k+j)/2} P^+ u|||^2, \\
 B_{3,k}^j &= \operatorname{Im} \alpha_j \bar{b}_k |||\partial_x|^{m-(k+j)/2} P^- u|||^2.
 \end{aligned}$$

We collect coefficients of derivative losses with rearranging the summation order. Note that for any sequence $\{c_{j,k}\}$, it holds that

$$(2.8) \quad \sum_{j=1}^p \sum_{k=1}^{p-j} c_{j,k} = \sum_{j=1}^{p-1} \sum_{k=0}^{p-1-j} c_{j,k+1} = \sum_{j=1}^{p-1} \sum_{k=1}^j c_{k,j-k+1}.$$

It is easy to see that

$$\left| \sum_{j=1}^{2m} \sum_{k=2m-j}^{2m} (A_{2,k}^j + A_{3,k}^j + B_{2,k}^j + B_{3,k}^j) \right| \lesssim \|u\|^2.$$

Then, by (2.8), we have

$$\begin{aligned}
 \sum_{j=1}^{2m-1} \sum_{k=1}^{2m-1-j} (A_{2,k}^j + B_{2,k}^j) &= \sum_{j=1}^{2(m-1)} \sum_{k=1}^j (A_{2,j-k+1}^k + B_{2,j-k+1}^k) \\
 &= \sum_{j=1}^{2(m-1)} \sum_{k=1}^j (1 + (-1)^{j-k+1}) \operatorname{Im} \alpha_k \bar{a}_{j-k+1} \langle D_x^{2m-1-j} \overline{v^-}, v^+ \rangle.
 \end{aligned}$$

Similarly, we obtain

$$\begin{aligned}
 \sum_{j=1}^{2m-1} \sum_{k=1}^{2m-1-j} A_{3,k}^j &= \sum_{j=1}^{2(m-1)} \sum_{k=1}^j (-1)^{j-k+1} \operatorname{Im} \alpha_k \bar{b}_{j-k+1} |||\partial_x|^{m-(j+1)/2} P^+ u|||^2, \\
 \sum_{j=1}^{2m-1} \sum_{k=1}^{2m-1-j} B_{3,k}^j &= \sum_{j=1}^{2(m-1)} \sum_{k=1}^j \operatorname{Im} \alpha_k \bar{b}_{j-k+1} |||\partial_x|^{m-(j+1)/2} P^- u|||^2.
 \end{aligned}$$

This concludes the proof of (2.4). For the proof of (2.5), we set $v^+ := P^- u$ and $v^- := P^+ u$.

Then, they satisfy (2.6) and (2.7). Therefore, the exactly same proof works. \square

Now we are ready to prove Proposition 2.1. Though we can prove it directly without using Lemma 2.3, we give the proof of it by the lemma.

Proof of Proposition 2.1. Note that $\langle P^+ f, P^- g \rangle = \langle P^- f, P^+ g \rangle = 0$ for any functions f, g . This implies that $\langle P_{\neq 0} \bar{f}, P_{\neq 0} g \rangle = \langle \overline{P^- f}, P^+ g \rangle + \langle \overline{P^+ f}, P^- g \rangle$. So, collecting (2.4) and (2.5), we obtain

$$\left| \frac{d}{dt} \left(\|P_{\neq 0} u\|^2 + \sum_{j=1}^{m-1} \operatorname{Re} \alpha_{2j} \langle D_x^{-2j} P_{\neq 0} \bar{u}, P_{\neq 0} u \rangle \right) + \sum_{j=1}^{m-1} \lambda_{2j} \| |\partial_x|^{m-j} P_{\neq 0} u \|^2 + \sum_{j=1}^m \lambda_{2j-1} \langle D_x^{2(m-j)+1} P_{\neq 0} u, P_{\neq 0} u \rangle \right| \leq C \|u\|^2.$$

We also note that $\gamma_k = \alpha_{2k}$. Finally, it is easy to see that

$$\left| \frac{d}{dt} \|P_0 u\|^2 + \sum_{j=1}^{m-1} \lambda_{2j} \| |\partial_x|^{m-j} P_0 u \|^2 + \sum_{j=1}^m \lambda_{2j-1} \langle D_x^{2(m-j)+1} P_0 u, P_0 u \rangle \right| \leq C \|u\|^2.$$

Therefore, we have (2.1). \square

The terms $\lambda_j^- \| |\partial_x|^{m-j/2} P^- u \|^2$ (resp. $\lambda_j^- \| |\partial_x|^{m-j/2} P^+ u \|^2$) in (2.4) (resp. (2.5)) with $1 \leq j \leq 2j^* - 1$ in Lemma 2.3 are unfavorable in our argument to prove Proposition 2.2. Indeed, Proposition 2.2 is used to show “Elliptic type” in Theorem 1.1 when $\lambda_{2j^*-1} \neq 0$ under the assumption $\lambda_j = 0$ for $1 \leq j \leq 2(j^* - 1)$. So, we analyse the coefficients λ^- below in order to ensure the condition $\lambda_j = 0$ for $1 \leq j \leq 2(j^* - 1)$ implies $\lambda_j^+ = \lambda_n^- = 0$ for $1 \leq j \leq 2(j^* - 1)$ and $1 \leq n \leq 2j^* - 1$.

Lemma 2.4. *It holds that*

$$\lambda_{j+1}^- = -\frac{1}{2} \sum_{l=1}^{j-1} (1 + (-1)^l) (\operatorname{Re} a_l) \lambda_{j+1-l}^- + \frac{1}{2} \sum_{l=1}^{j-1} \sum_{k=1}^{j-l} (1 + (-1)^l) (\operatorname{Im} a_l) \operatorname{Re} \bar{b}_{j-l-k+1} \alpha_k$$

for $1 \leq j \leq 2(m-1)$.

Proof. By the definitions of λ_j^- and α_k , we have

$$\lambda_{j+1}^- = - \sum_{l=1}^j \operatorname{Im} b_l \bar{b}_{j-l+1} + \frac{1}{2} \sum_{l=1}^j \sum_{k=1}^{l-1} (1 + (-1)^{l-k}) \operatorname{Im} \bar{b}_{j-l+1} \bar{a}_{l-k} \alpha_k =: A + B.$$

It is easy to see that $A = 0$. Observe that

$$(2.9) \quad \sum_{l=1}^p \sum_{k=1}^{l-1} c_{l-k} d_l e_k = \sum_{l=1}^{p-1} \sum_{k=1}^{p-l} c_l d_{l+k} e_k$$

for any sequences $\{c_j\}$, $\{d_j\}$ and $\{e_j\}$. This implies that

$$B = \frac{1}{2} \sum_{l=1}^{j-1} \sum_{k=1}^{j-l} (1 + (-1)^l) \operatorname{Im} \bar{b}_{j-l-k+1} \bar{a}_l \alpha_k$$

$$= \frac{1}{2} \sum_{l=1}^{j-1} \sum_{k=1}^{j-l} (1 + (-1)^l) ((\operatorname{Re} a_l) \operatorname{Im} \bar{b}_{j-l-k+1} \alpha_k - (\operatorname{Im} a_l) \operatorname{Re} \bar{b}_{j-l-k+1} \alpha_k).$$

Here we used the fact that $\operatorname{Im} cd = (\operatorname{Re} c) \operatorname{Im} d + (\operatorname{Im} c) \operatorname{Re} d$ for any $c, d \in \mathbb{C}$. This completes the proof. \square

Lemma 2.5. *Assume that there exists $j^* \in \mathbb{N}$ such that $\lambda_{2j} = 0$ for $1 \leq j \leq j^*$. Then, it holds that $\operatorname{Im} a_{2j} = \lambda_{2j}^+ = 0$ for $1 \leq j \leq j^*$ and $\lambda_j^- = 0$ for $1 \leq j \leq 2j^* + 3$.*

Proof. First note that $\lambda_1^- = \lambda_2^- = \lambda_3^- = 0$ even without the hypothesis. Indeed, it is clear that $\lambda_1^- = 0$. We also have $\lambda_2^- = -\operatorname{Im} \bar{b}_1 \alpha_1 = 0$ and $\lambda_3^- = -\operatorname{Im} \bar{b}_2 \alpha_1 - \operatorname{Im} \bar{b}_1 \alpha_2 = 0$ since $\alpha_1 = b_1$ and $\alpha_2 = b_2$. Assume that there exists $j^* \in \mathbb{N}$ such that $\lambda_{2j} = 0$ for $1 \leq j \leq j^*$. The rest of proof proceeds by the induction on j . We prove the following claim: it holds that

$$\operatorname{Im} a_{2j} = \lambda_{2j}^+ = \lambda_{2j+2}^- = \lambda_{2j+3}^- = \sum_{k=1}^{j-1} \operatorname{Im} \bar{b}_{2(j-k)} \alpha_{2k} = 0$$

for $1 \leq j \leq j^*$. It is easy to see that the claim above with $j = 1$ follows. Indeed, by the definition of α_j , we obtain $\alpha_3 = b_3 - \bar{a}_2 b_1$ and $\alpha_4 = b_4 - \bar{a}_2 b_2$, which implies that

$$\begin{aligned} \lambda_4^- &= -\operatorname{Im} \bar{b}_3 b_1 - \operatorname{Im} |b_2|^2 - \operatorname{Im} \bar{b}_1 b_3 + \operatorname{Im} |b_1|^2 \bar{a}_2 = 0, \\ \lambda_5^- &= -\operatorname{Im} \bar{b}_4 b_1 - \operatorname{Im} \bar{b}_3 b_2 - \operatorname{Im} \bar{b}_2 (b_3 - \bar{a}_2 b_1) - \operatorname{Im} \bar{b}_1 (b_4 - \bar{a}_2 b_2) = 0 \end{aligned}$$

since $\lambda_2 = 2 \operatorname{Im} a_2 = 0$. We also have $\lambda_2^+ = \operatorname{Im} \bar{b}_2 \alpha_2 = 0$ easily. Next, we assume that the claim above holds for $j(\leq j^* - 1)$. By the hypothesis, it holds that $\lambda_{2j+2}^- = \lambda_{2j+3}^- = 0$. Thus, by Remark 2.1, we have $\lambda_{2j+2}^+ = 0$. We claim that $M := \sum_{l=1}^j \operatorname{Im} \bar{b}_{2(j-l+1)} \gamma_l = 0$. Indeed, we see from the definition of γ_l that

$$M = \sum_{l=1}^j \operatorname{Im} \bar{b}_{2(j-l+1)} b_{2l} - \sum_{l=1}^j \sum_{k=1}^{l-1} \operatorname{Im} \bar{b}_{2(j-l+1)} \bar{a}_{2(l-k)} \gamma_k =: A + B.$$

It is easy to see that $A = 0$. We have

$$\begin{aligned} B &= - \sum_{l=1}^{j-1} \sum_{k=1}^{j-l} \operatorname{Im} \bar{b}_{2(j-l-k+1)} \bar{a}_{2l} \gamma_k \\ &= - \sum_{l=1}^{j-1} (\operatorname{Re} a_{2l}) \sum_{k=1}^{j-l} \operatorname{Im} \bar{b}_{2(j-l-k+1)} \gamma_k + \sum_{l=1}^{j-1} (\operatorname{Im} a_{2l}) \sum_{k=1}^{j-l} \operatorname{Re} \bar{b}_{2(j-l-k+1)} \gamma_k = 0 \end{aligned}$$

by (2.9) and the hypothesis. This shows that $\operatorname{Im} a_{2j+2} = 0$ by the definition of λ_{2j+2} . By Lemma 2.4, we obtain $\lambda_{2j+4}^- = \lambda_{2j+5}^- = 0$, which completes the proof. \square

REMARK 2.2. From the proof of the above lemma, we also see that

$$\lambda_{2j^*+2} = 2 \operatorname{Im} a_{2j^*+2}, \quad \lambda_{2j^*+4} = 2 \operatorname{Im} a_{2j^*+4}$$

when $\lambda_{2j} = 0$ for $1 \leq j \leq j^*$.

Now, we prove Proposition 2.2.

Proof of Proposition 2.2. We give only the proof of (2.2) since we can show (2.3) in the same manner. When $j^* = 1$, we see from the definition that $\lambda_n^- = 0$ for $n = 1, 2, 3$. When $j^* \geq 2$, Lemma 2.5 implies that $\lambda_j^+ = \lambda_j^- = 0$ for $1 \leq j \leq 2(j^* - 1)$ and $\lambda_{2j^*-1}^- = \lambda_{2j^*}^- = \lambda_{2j^*+1}^- = 0$. This together with Remark 2.1 implies $\lambda_{2j^*-1}^+ \neq 0$. By (2.4), interpolation inequalities and the Young inequality, we have

$$\begin{aligned} & \left| \frac{d}{dt} \left(\|P^+ u\|^2 + \sum_{j=1}^{2m-1} \operatorname{Re} \alpha_j \langle D_x^{-j} \overline{P^- u}, P^+ u \rangle \right) + \sum_{j=2j^*+2}^{2m-1} \lambda_j^- \| |\partial_x|^{m-j/2} P^- u \|^2 \right. \\ & \left. + \lambda_{2j^*-1}^+ \| |\partial_x|^{m-j^*+1/2} P^+ u \|^2 \right| \lesssim \|u\|^2 + \| |\partial_x|^{m-j^*} P^+ u \|^2 \end{aligned}$$

Thus, we only need to show

$$\begin{aligned} (2.10) \quad & \left| \frac{d}{dt} \sum_{k=1}^{2(m-j^*-1)} \beta_k^+ F_k^-(u) - \sum_{j=2j^*+2}^{2m-1} \lambda_j^- \| |\partial_x|^{m-j/2} P^- u \|^2 \right| \\ & \lesssim \|u\|^2 + \| |\partial_x|^{m-j^*} P^+ u \|^2. \end{aligned}$$

Put $v = |\partial_x|^{-(k+2)/2} P_{\neq 0} u$. Since v satisfies (1.1), by (2.5), we have

$$\begin{aligned} & \left| \frac{d}{dt} \left(\|P^- v\|^2 + \sum_{j=1}^{2m-1} \operatorname{Re} \alpha_j \langle D_x^{-j} \overline{P^+ v}, P^- v \rangle \right) + \sum_{j=2j^*-1}^{2m-1} (-1)^j \lambda_j^+ \| |\partial_x|^{m-j/2} P^- v \|^2 \right| \\ & \lesssim \|v\|^2 + \| |\partial_x|^{m-j^*} P^+ v \|^2. \end{aligned}$$

Thus, we obtain

$$\begin{aligned} & \left| \sum_{k=1}^{2(m-j^*-1)} \beta_k^+ \left(\frac{d}{dt} F_k^-(u) + \sum_{j=2j^*-1}^{2m-k-3} (-1)^j \lambda_j^+ \| |\partial_x|^{m-(j+k+2)/2} P^- u \|^2 \right) \right| \\ & \lesssim \|u\|^2 + \| |\partial_x|^{m-j^*} P^+ u \|^2. \end{aligned}$$

By (2.8), we have

$$\begin{aligned} & \sum_{k=1}^{2(m-j^*-1)} \sum_{j=2j^*-1}^{2m-k-3} (-1)^j \beta_k^+ \lambda_j^+ \| |\partial_x|^{m-(j+k+2)/2} P^- u \|^2 \\ & = \sum_{k=1}^{2(m-j^*-1)} \sum_{j=1}^k (-1)^{k-j+1} \beta_j^+ \lambda_{2j^*+k-j-1}^+ \| |\partial_x|^{m-j^*-(k+1)/2} P^- u \|^2. \end{aligned}$$

Therefore, by the definition of β_k^+ , we conclude (2.10). \square

3. Proof of main theorem

In this section, we show Theorem 1.1.

DEFINITION 3. For $f \in L^2(\mathcal{M})$ and $N > 0$, we define

$$E(f; N) := \|f\|^2 + N \|\partial_x^{-m} P_{\neq 0} f\|^2 + \sum_{j=1}^{m-1} \operatorname{Re} \gamma_j \langle D_x^{-2j} P_{\neq 0} \bar{f}, P_{\neq 0} f \rangle.$$

We choose N sufficiently large so that Lemma 3.1 holds. If there is no confusion, we write $E(f) := E(f; N)$.

Lemma 3.1. *Let $N > 0$ be sufficiently large. Then, for any $f \in L^2(\mathcal{M})$ it holds that*

$$\frac{1}{2}E(f) \leq \|f\|^2 + N\|\partial_x^{-m}P_{\neq 0}f\|^2 \leq 2E(f).$$

Proof. The Gagliardo-Nirenberg inequality and the Young inequality show that

$$(3.1) \quad \sum_{j=1}^{m-1} |\operatorname{Re} \gamma_j \langle D_x^{-2j} P_{\neq 0} \bar{f}, P_{\neq 0} f \rangle| \leq \frac{1}{2}\|f\|^2 + C\|\partial_x^{-m}P_{\neq 0}f\|^2.$$

So, it suffices to choose $N = 2C$. \square

We prove the first part of Theorem 1.1.

Proof of “Dispersive type” in Theorem 1.1. We consider our problem only on $[0, \infty)$ since the result on $(-\infty, 0]$ follows from the same argument. Let $T > 0$, which can be arbitrary large. We first show the a priori estimate $\sup_{t \in [0, T]} \|u(t)\| \leq C\|\varphi\|$. We assume that u satisfies (1.1) and (1.2). Then, it is easy to see that $\frac{d}{dt}\|\partial_x^{-m}P_{\neq 0}u\|^2 \leq 2|\langle D_t \partial_x^{-2m}P_{\neq 0}u, P_{\neq 0}u \rangle| \leq C\|u\|^2$. This together with (2.1), Lemma 3.1 and $\lambda_j = 0$ for $1 \leq j \leq 2m - 1$ implies that $\frac{d}{dt}E(u(t)) \leq CE(u(t))$ on $[0, T]$. Thus, by the Gronwall inequality and Lemma 3.1, we obtain the a priori estimate. Next, we show the existence. Let $\varphi_n = \mathcal{F}^{-1}\chi_{(|\xi| < n)}\mathcal{F}\varphi$ for $n \in \mathbb{N}$. Then, we have the solution u_n of (1.1) with $u_n(0) = \varphi_n$ by (1.8). Moreover, $u_n \in C([0, T]; L^2(\mathcal{M}))$ since $|\sum_{j=0}^{2m} \xi^{2m-j} X_j| \leq C(\{a_j\}, \{b_j\}, n)$ for $|\xi| < n$. Since $\{\varphi_n\}$ is a Cauchy sequence in $L^2(\mathcal{M})$, by the a priori estimate, we conclude $\{u_n\}$ is also a Cauchy sequence in $C([0, T]; L^2(\mathcal{M}))$. Thus, we obtain the solution $u \in C([0, T]; L^2(\mathcal{M}))$ of (1.1)–(1.2) as the limit of u_n . Finally, the uniqueness easily follows from the a priori estimate. \square

Proof of “Parabolic type” in Theorem 1.1. We use the argument from the proof of Theorem 1.2 in [21]. We consider only the case $\lambda_{2j^*} > 0$ since the other case follows from the same argument. Let $T > 0$, which can be arbitrary large. By the Gagliardo-Nirenberg inequality and the Young inequality, we have

$$\left| \sum_{j=j^*+1}^{m-1} \lambda_{2j} \|\partial_x^{m-j} u\|^2 + \sum_{j=j^*+1}^m \lambda_{2j-1} \langle D_x^{2(m-j)+1} u, u \rangle \right| \leq \frac{1}{2} \lambda_{2j^*} \|\partial_x^{m-j^*} u\|^2 + C\|u\|^2.$$

Recall that $\lambda_j = 0$ for $1 \leq j \leq 2j^* - 1$. Therefore, in the same manner as the proof of “Dispersive type,” we obtain the a priori estimate:

$$\sup_{t \in [0, T]} \left(\|u(t)\|^2 + \frac{\lambda_{2j^*}}{2} \int_0^t \|\partial_x^{m-j^*} u(\tau)\|^2 d\tau \right) \leq C\|\varphi\|^2.$$

It then follows that we have the unique existence of the solution $u \in C([0, T]; L^2(\mathcal{M})) \cap L^2([0, T]; H^{m-j^*}(\mathcal{M}))$, which implies that $u(t) \in H^{m-j^*}(\mathcal{M})$ for a.e. $t \in [0, T]$. Let $0 < \varepsilon < T$. Then there exists $t_0 \in (0, \varepsilon/2)$ such that $u(t_0) \in H^{m-j^*}(\mathcal{M})$. Since $\langle \partial_x \rangle^{m-j^*} u$ satisfies (1.1)–(1.2) with initial data $\varphi := \langle \partial_x \rangle^{m-j^*} u(t_0) \in L^2(\mathcal{M})$, applying the same argument as above, we conclude $\langle \partial_x \rangle^{m-j^*} u \in C([t_0, T]; L^2(\mathcal{M})) \cap L^2([t_0, T]; H^{m-j^*}(\mathcal{M}))$. That is, $u \in C([t_0, T]; H^{m-j^*}(\mathcal{M})) \cap L^2([t_0, T]; H^{2(m-j^*)}(\mathcal{M}))$. We can choose t_1 so that $\varepsilon/2 < t_1 <$

$\varepsilon/2 + \varepsilon/4$ and $u(t_1) \in H^{2(m-j^*)}(\mathcal{M})$. Again, applying the same argument as above with the initial data $\varphi := \langle \partial_x \rangle^{2(m-j^*)} u(t_0) \in L^2(\mathcal{M})$, we conclude $u \in C([t_1, T]; H^{2(m-j^*)}(\mathcal{M})) \cap L^2([t_1, T]; H^{3(m-j^*)}(\mathcal{M}))$. By repeating this process, we conclude $u \in C([\varepsilon, T]; H^{k(m-j^*)}(\mathcal{M}))$ for any $k \in \mathbb{N}$, which implies $u \in C^\ell([\varepsilon, T]; H^{k(m-j^*)-2m\ell}(\mathcal{M}))$ for any $k, \ell \in \mathbb{N}$ by (1.1). By the Sobolev embedding, we obtain $u \in C^\infty([\varepsilon, T] \times \mathcal{M})$. Since we can take $\varepsilon > 0$ arbitrary small and $T > 0$ arbitrary large, we conclude $u \in C^\infty((0, \infty) \times \mathcal{M})$. Finally, we show the nonexistence result by contradiction. Assume that there exists a solution $u \in C((-\delta, 0]; L^2(\mathcal{M}))$ of (1.1)–(1.2) with $\varphi \in L^2(\mathcal{M}) \setminus C^\infty(\mathcal{M})$. We take t_0 such that $-\delta < t_0 < 0$. Then, as we proved above, we have $u \in C^\infty((t_0, 0] \times \mathcal{M})$, which contradicts to the assumption $\varphi = u(0) \notin C^\infty(\mathcal{M})$. \square

The following proposition is the main tool to show the result for “Elliptic type” in Theorem 1.1.

Proposition 3.2 (A smoothing for “Elliptic type”). *Let $u \in C([t_0, t_1]; L^2(\mathcal{M}))$ satisfy (1.1). Assume that there exists $j^* \in \mathbb{N}$ such that $\lambda_j = 0$ for $1 \leq j < 2j^* - 1$ and $\lambda_{2j^*-1} > 0$ (resp. < 0). Then, it follows that*

$$(3.2) \quad P^+u \text{ (resp. } P^-u) \in C((t_0, t_1]; H^{1/2}(\mathcal{M})) \quad (\text{forward smoothing}),$$

$$(3.3) \quad P^-u \text{ (resp. } P^+u) \in C([t_0, t_1]; H^{1/2}(\mathcal{M})) \quad (\text{backward smoothing}).$$

In particular, it holds that $u \in C^\infty((t_0, t_1) \times \mathcal{M})$.

Proof. We consider only the case $\lambda_{2j^*-1}^+ > 0$ since the same proof works for the case $\lambda_{2j^*-1}^+ < 0$. For simplicity, set

$$G^+(u) := \sum_{j=1}^{2m-1} \operatorname{Re} \alpha_j \langle D_x^{-j} \overline{P^-} u, P^+ u \rangle + \sum_{k=1}^{2(m-j^*-1)} \beta_k^+ F_k^-(u),$$

where F_k^- is defined in Proposition 2.2 and $\{\alpha_j\}$ and $\{\beta_k\}$ are defined in Definition 2. Set $M := \sup_{t \in [t_0, t_1]} \|u(t)\|$. Note that $\sup_{t \in [t_0, t_1]} (|G^+(u(t))| + |G^+(|\partial_x|^{1/2} u(t))|) \leq CM$ and $G^+ (|\partial_x|^{1/2} u(t))$ is continuous on $[t_0, t_1]$ by the presence of D_x^{-j} in the definition of $G^+(u)$ above. By the Gagliardo-Nirenberg inequality and the Young inequality, we have

$$\| |\partial_x|^{m-j^*} Qu \|^2 \leq \delta \| |\partial_x|^{m-j^*+1/2} Qu \|^2 + C\delta^{-1} \|u\|^2$$

for $\delta > 0$, $Q = P^+$ or P^- . Take $\delta > 0$ sufficiently small. Then, this together with (2.2) and (2.3) yields

$$\lambda_{2j^*-1}^+ \int_{t_0}^{t_1} \| |\partial_x|^{m-j^*+1/2} Qu(\tau) \|^2 d\tau \leq C(M)(1 + |t_1 - t_0|),$$

for $Q = P^+$ or P^- . By the interpolation, we also have

$$\begin{aligned}
 (3.4) \quad & \int_{t_0}^{t_1} \| |\partial_x|^s u(\tau) \|^2 d\tau \\
 &= \int_{t_0}^{t_1} (\| |\partial_x|^s P^- u(\tau) \|^2 + \| |\partial_x|^s P_0 u(\tau) \|^2 + \| |\partial_x|^s P^+ u(\tau) \|^2) d\tau \\
 &\leq C(M, \lambda_{2j^*-1}^+)(1 + |t_1 - t_0|)
 \end{aligned}$$

for $0 \leq s \leq m + j^* - 1/2$. It then follows that $\| |\partial_x|^{m-j^*+1/2} u(t) \| < \infty$ for a.e. $t \in [t_0, t_1]$. Then, for any $\varepsilon > 0$ there exists $t_* \in (t_0, t_0 + \varepsilon)$ such that $\| |\partial_x|^{m-j^*+1/2} u(t_*) \| < \infty$. Note that (2.2) holds even if we replace u with $|\partial_x|^{1/2} u$ since $|\partial_x|^{1/2} u$ satisfies (1.1). Thus,

$$(3.5) \quad \left| \frac{d}{dt} (\| |\partial_x|^{1/2} P^+ u \|^2 + G^+(|\partial_x|^{1/2} u)) + \lambda_{2j^*-1}^+ \| |\partial_x|^{m-j^*+1} P^+ u \|^2 \right| \\ \leq C \| |\partial_x|^{1/2} u \|^2 + C \| |\partial_x|^{m-j^*+1/2} P^+ u \|^2,$$

By the Gagliardo-Nirenberg inequality and the Young inequality, we have

$$\| |\partial_x|^{m-j^*+1/2} P^+ u \|^2 \leq \delta \| |\partial_x|^{m-j^*+1} P^+ u \|^2 + C \delta^{-1} \| |\partial_x|^{1/2} u \|^2$$

for $\delta > 0$. Taking $\delta > 0$ sufficiently small and integrating (3.5) on $[t_*, t] \subset [t_0, t_1]$ with (3.4), we obtain

$$(3.6) \quad \| |\partial_x|^{1/2} P^+ u(t) \|^2 + \frac{\lambda_{2j^*-1}^+}{2} \int_{t_*}^t \| |\partial_x|^{m-j^*+1} P^+ u(\tau) \|^2 d\tau \\ \leq C(M, \lambda_{2j^*-1}^+, |t_1 - t_0|) + \| |\partial_x|^{1/2} P^+ u(t_*) \|^2 < \infty$$

since $u(t_*) \in H^{m-j^*+1/2}(\mathcal{M})$. Therefore, by (3.5) again, it follows that for any $t_* \leq t' \leq t \leq t_1$

$$\left| \| |\partial_x|^{1/2} P^+ u(t) \|^2 - \| |\partial_x|^{1/2} P^+ u(t') \|^2 \right| \\ \leq \left| \| |\partial_x|^{1/2} P^+ u(\tau) \|^2 + G^+(|\partial_x|^{1/2} u) \right|_{\tau=t'}^{t-t} + \lambda_{2j^*-1}^+ \int_{t'}^t \| |\partial_x|^{m-j^*+1} P^+ u(\tau) \|^2 d\tau \\ + \left| \left[G^+(|\partial_x|^{1/2} u) \right]_{\tau=t'}^{t-t} \right| + \lambda_{2j^*-1}^+ \int_{t'}^t \| |\partial_x|^{m-j^*+1} P^+ u(\tau) \|^2 d\tau \\ \leq C \int_{t'}^t \| |\partial_x|^{1/2} u(\tau) \|^2 d\tau + C \int_{t'}^t \| |\partial_x|^{m-j^*+1/2} P^+ u(\tau) \|^2 d\tau \\ + \lambda_{2j^*-1}^+ \int_{t'}^t \| |\partial_x|^{m-j^*+1} P^+ u(\tau) \|^2 d\tau + \left| \left[G^+(|\partial_x|^{1/2} u) \right]_{\tau=t'}^{t-t} \right|.$$

(3.4), (3.6) and the dominated convergence theorem imply that the right-hand side goes to 0 as $|t - t'| \rightarrow 0$, which shows that $\| |\partial_x|^{1/2} P^+ u(t) \|^2$ is continuous on $[t_*, t_1]$. The fact $P^+ u \in C([t_0, t_1]; L^2(\mathcal{M}))$ with $P^+ u \in L^\infty([t_*, t_1]; H^{1/2}(\mathcal{M}))$ yields $P^+ u \in C_w([t_*, t_1]; H^{1/2}(\mathcal{M}))$. Combining the continuity of $\| |\partial_x|^{1/2} P^+ u(t) \|^2$ and the weak continuity of $P^+ u(t)$ in $H^{1/2}(\mathcal{M})$, we obtain $P^+ u \in C([t_*, t_1]; H^{1/2}(\mathcal{M}))$. Since we can take $\varepsilon > 0$ arbitrary small, we get $P^+ u \in C((t_0, t_1]; H^{1/2}(\mathcal{M}))$. We also obtain $P^- u \in C([t_0, t_1]; H^{1/2}(\mathcal{M}))$ in the same manner. Therefore, $u = P^- u + P_0 u + P^+ u \in C((t_0, t_1); H^{1/2}(\mathcal{M}))$. By repeating this process, we also obtain $u \in C((t_0, t_1); H^{k/2}(\mathcal{M}))$ for any $k \in \mathbb{N}$, which yields $u \in C^\infty((t_0, t_1) \times \mathcal{M})$ since u satisfies (1.1). \square

Proof of “Elliptic type” in Theorem 1.1. We use the argument from the proof of Theorem 1.2 in [20]. We consider only the case $\lambda_{2j^*-1} > 0$ since the case $\lambda_{2j^*-1} < 0$ follows from the same argument. Let $\varphi \in L^2(\mathcal{M})$ satisfy $P^+ \varphi \notin H^{1/2}(\mathcal{M})$. We prove Theorem 1.1 by contradiction. We assume that there exists $u \in C([-\delta, 0]; L^2(\mathcal{M}))$ satisfying (1.1)–(1.2) on $[-\delta, 0]$. Then, we have $P^+ u \in C((-\delta, 0]; H^{1/2}(\mathcal{M}))$ by Proposition 3.2. However, it contradicts to $P^+ \varphi \notin H^{1/2}(\mathcal{M})$. This proof works even if we replace P^+ and $[-\delta, 0]$ with P^- and $[0, \delta]$, respectively. Similarly, we can show that for any $\delta > 0$ there exist no so-

lution $u(t, x)$ of (1.1)–(1.2) with $u(0, x) = \varphi(x) \in L^2(\mathcal{M}) \setminus C^\infty(\mathcal{M})$ on $[-\delta, \delta]$ satisfying $u \in C([-\delta, \delta]; L^2(\mathcal{M}))$. \square

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