Weakly hyperbolic equations with non-analytic coefficients and lower order terms

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Abstract In this paper we consider weakly hyperbolic equations of higher orders in arbitrary dimensions with time-dependent coefficients and lower order terms. We prove the Gevrey well-posedness of the Cauchy problem under C^k -regularity of coefficients of the principal part and natural Levi conditions on lower order terms which may be only continuous. In the case of analytic coefficients in the principal part we establish the C^{∞} well-posedness. The proofs are based on using the quasi-symmetriser for the corresponding companion system and inductions on the order of equation and on the frequency regions. The main novelty compared to the existing literature is the possibility to include lower order terms to the equation (which have been untreatable until now in these problems) as well as considering any space dimensions. We also give results on the ultradistributional and distributional well-posedness of the problem, and we look at new effects for equations with discontinuous lower order terms.

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1 Introduction

In this paper we study the well-posedness of the weakly hyperbolic Cauchy problem

$$\begin{cases} D_t^m u + \sum_{j=0}^{m-1} A_{m-j}(t, D_x) D_t^j u = 0, & (t, x) \in [0, T] \times \mathbb{R}^n, \\ D_t^{h-1} u(0, x) = g_h(x), & h = 1, \dots, m, \end{cases}$$
(1)

where each $A_{m-j}(t, D_x)$ is a differential operator of order m - j with continuous coefficients depending only on t. Later we will also relax the continuity assumption replacing it by the boundedness. As usual, $D_t = \frac{1}{i}\partial_t$ and $D_x = \frac{1}{i}\partial_x$. Let $A_{(m-j)}$ denote the principal part of the operator A_{m-j} and let $\lambda_l(t, \xi), l = 1, ..., m$, be the real-valued roots of the characteristic polynomial which we write as

$$\tau^{m} + \sum_{j=0}^{m-1} A_{(m-j)}(t,\xi)\tau^{j} = \tau^{m} + \sum_{j=0}^{m-1} \sum_{|\gamma|=m-j} a_{m-j,\gamma}(t)\xi^{\gamma}\tau^{j}.$$
 (2)

This means that

$$\tau^{m} + \sum_{j=0}^{m-1} \sum_{|\gamma|=m-j} a_{m-j,\gamma}(t) \xi^{\gamma} \tau^{j} = \prod_{l=1}^{m} (\tau - \lambda_{l}(t,\xi)).$$
(3)

If (1) is strictly hyperbolic and its coefficients are in the Hölder class in $t, A_{m-j}(\cdot, \xi) \in C^{\alpha}([0, T]), 0 < \alpha < 1$, for all j and ξ , it was shown in [11, Case 3; Remark 8] by the authors that the Cauchy problem (1) is well-posed in Gevrey classes $G^{s}(\mathbb{R}^{n})$ provided that $1 \leq s < 1 + \frac{\alpha}{1-\alpha}$. If $\alpha = 1$, it is sufficient to assume the Lipschitz continuity of coefficients to get the well-posedness in G^{s} for all $s \geq 1$. Earlier, for certain second order equations the same Gevrey index was obtained by Colombini et al. [3], who also showed that it is sharp. We also refer to [11, Remark 16] for the Gevrey-Beurling ultradistributional well-posedness of (1) for $1 \leq s \leq 1 + \frac{\alpha}{1-\alpha}$. In this paper we will deal with more regular coefficients in the weakly hyperbolic case.

The well-posedness of the weakly hyperbolic equations has been a challenging problem for a long time. For example, even for the second order Cauchy problem in one space dimension,

$$\partial_t^2 u - a(t, x) \partial_x^2 u = 0, \quad u(0, x) = g_1(x), \quad \partial_t u(0, x) = g_2(x),$$
 (4)

up until now there is no characterisation of smooth functions $a(t, x) \ge 0$ for which (4) would be C^{∞} well-posed. On one hand, there are sufficient conditions. For example, Oleinik has shown in [16] that (4) is C^{∞} well-posed provided there is a constant C > 0 such that $Ca(t, x) + \partial_t a(t, x) \ge 0$. In the case of a(t, x) = a(t) depending only on t, when the problem becomes

$$\partial_t^2 u - a(t)\partial_x^2 u = 0, \quad u(0,x) = g_1(x), \quad \partial_t u(0,x) = g_2(x),$$
 (5)

the Oleinik's condition is satisfied for $a(t) \ge 0$ with $a'(t) \ge 0$. On the other hand, in the celebrated paper [8], Colombini and Spagnolo constructed a C^{∞} function $a(t) \ge 0$ such that (5) is not C^{∞} well-posed. The situation becomes even more complicated if one adds mixed terms to (5), even depending only on *t* and analytic. For example, the Cauchy problem for the equation

$$\partial_t^2 u - 2t \partial_t \partial_x u + t^2 \partial_x^2 u = 0$$

is Gevrey G^s well-posed for s < 2 while it is ill-posed for any s > 2. For other positive and negative results for second order equations with time-dependent coefficients we refer to seminal papers of Colombini et al. [3,5], and to Nishitani [15] for the necessary and sufficient conditions for the C^{∞} well-posedness of (4) with analytic $a(t, x) \ge 0$ in one dimension.

A reasonable substitute for the C^{∞} well-posedness in the weakly hyperbolic setting is the well-posedness in the space $G^{\infty} = \bigcup_{s>1} G^s$. Thus, Colombini et al., proved in [4] that for every C^{∞} function $a(t) \ge 0$, the Cauchy problem (5) is G^{∞} wellposed. More precisely, they showed that if a(t) is in C^k , it is well posed in G^s with $s \le 1 + k/2$, and if a(t) is analytic, it is C^{∞} well-posed.

From another direction, there are also general results for (1). For example, it was shown by Bronshtein in [2] that, in particular, the Cauchy problem (1) with C^{∞} coefficients is G^s well-posed provided that $1 \le s < 1 + \frac{1}{m-1}$. In some cases, this can be improved. For example, for constant multiplicities, see paper [6] by Colombini and Kinoshita in one-dimension (see also D'Ancona and Kinoshita [9]), and the authors' paper [11] for further improvements of Gevrey indices and all dimensions, with a survey of literature therein.

In this paper our interest in analysing the Cauchy problem (1) is motivated by

- (A) allowing any space dimension $n \ge 1$;
- (B) considering the effect of lower order terms or, rather, the properties of the lower order terms which do not influence the results on the Gevrey well-posedness (we will look at new effects for both continuous and discontinuous lower order terms); the inclusion of lower order terms in this setting has been untreatable by previous methods;
- (C) providing well-posedness results in spaces of distributions and ultradistributions.

Our main reference here is the paper [13] of Kinoshita and Spagnolo who have studied the Cauchy problem (1) for operators with homogeneous symbols in one dimension, $x \in \mathbb{R}$. Under the condition

$$\exists M > 0: \quad \lambda_i(t,\xi)^2 + \lambda_j(t,\xi)^2 \le M(\lambda_i(t,\xi) - \lambda_j(t,\xi))^2,$$

for $1 \le i, j \le m, t \in [0,T]$, for all ξ , (6)

on the roots $\lambda_i(t,\xi)$ they have obtained the following well-posedness result:

Theorem 1 ([13]) Assume that n = 1 and that the differential operator is homogeneous, i.e. $A_{m-j}(t,\xi) = A_{(m-j)}(t,\xi) = a_{m-j}(t)\xi^{m-j}$ for all j = 0, ..., m-1. If $a_{m-j} \in C^{\infty}([0,T])$ and the characteristic roots are real and satisfy (6), then the Cauchy problem (1) is well-posed in any Gevrey space. More precisely, if $a_j \in C^k([0, T])$ for some $k \ge 2$ then we have G^s -well-posedness for

$$1 \le s < 1 + \frac{k}{2(m-1)}$$

The proof is based on the construction of a quasi-symmetriser $Q_{\varepsilon}^{(m)}$ which thanks to the condition (6) is nearly diagonal. Previously, equations of second and third order with analytic coefficients, still with n = 1 and without low order terms, have been analysed by Colombini and Orrú [7]. They have shown the C^{∞} well-posedness of (1) under assumption (6). Moreover, if all the coefficients $a_{m-j}(t)$ vanish at t = 0, they showed that the *condition* (6) *is also necessary*. So, for us it will be natural to adopt (6) for our analysis.

Let us briefly discuss the difficulties of aims (A)-(C) above. For the dimensional extension (A), even under condition (6) on the characteristic roots, for space dependent coefficients such an extension is impossible, see e.g. Bernardi and Bove [1], for examples of second order operators with polynomial coefficients for which the C^{∞} well-posedness fails for any $n \ge 2$. It is interesting to note that for these examples the usual Ivrii-Petkov conditions on lower order terms are also satisfied. As we will show, the C^{∞} (and other) well-posedness holds in our case in any dimension $n \ge 1$ since the coefficients depend only on time. In part (B), the proof of the well-posedness for equations with lower order terms highlights several interesting and somewhat surprising phenomena. For example, if the coefficients of the principal part are analytic and the lower order terms are only bounded (in particular, they may be discontinuous, or may exhibit more irregular oscillating behaviour), but the Cauchy data is Gevrey, we still obtain the solution in Gevrey spaces. Indeed, the Levi conditions in this paper control the zeros of the lower order terms but not their regularity. Finally, aim (C) is motivated by an interesting and challenging problem for weakly hyperbolic equations: analysing the propagation of singularities. For this, in order to be able to use also non-Gevrey techniques, we need to have first well-posedness in some bigger space. This will be achieved for the Cauchy problem (1) in the spaces of Beurling Gevrey ultradistributions. A subtle point of this construction is that we will have to use the Beurling Gevrey ultradistributions and not the usual Roumieu Gevrey ultradistributional class. In the case of the analytic principal part we will obtain the well-posedness in the usual space of distributions.

In particular, in this paper we extend Theorem 1 to weakly hyperbolic equations with non-homogeneous symbols and in any space dimension $n \ge 1$, and find suitable assumptions on the lower order terms for the Gevrey well-posedness. Already from the beginning we deviate from [13] by using pseudo-differential techniques to reduce the equation to the system. This will allow us to treat all the dimensions $n \ge 1$. However, the main challenge in the present paper is the analysis of the lower order terms. In fact, in most (if not all) of the literature on the application of the quasi-symmetriser to weakly hyperbolic equations the considered equations are always assumed to have homogeneous symbols. It is our intension to show that the quasi-symmetriser can be effectively used to control parts of the energy corresponding to the lower order terms. It is interesting to see the appearing Levi conditions expressing the dependence of

the lower order terms on the principal part of the operator. Such control becomes possible by exploiting the Sylvester form of the system corresponding to Eq. (1), and the structure of the quasi-symmetriser.

An interesting effect that we observe is that the results remain true assuming just the continuity of the lower order terms in time. For example, we will have the C^{∞} well-posedness for equations with analytic coefficients in the principal part and only continuous lower order terms. Moreover, we give a variant of our results with only assuming the boundedness of lower order terms in time (instead of continuity).

In this paper we formulate the conditions on the lower order terms in terms of the symbols A_{m-j+1} . Note that in (1), the operator $A_{m-j+1}(t, D_x)$ is the coefficient in front of the derivative D_t^{j-1} . We assume that there is some constant C > 0 such that we have

$$|(A_{m-j+1} - A_{(m-j+1)})(t,\xi)| \le C \sum_{i=1}^{m} \left| \sum_{\substack{1 \le l_1 < \dots < l_{m-j} \le m \\ l_h \neq i \ \forall h}} \lambda_{l_1}(t,\xi) \cdots \lambda_{l_{m-j}}(t,\xi) \right|,$$
(7)

for all $t \in [0, T]$, j = 1, ..., m, and for ξ away from 0 (i.e., for $|\xi| \ge R$ for some R > 0). Note that in terms of the coefficients of the original equation, using (2) we have

$$(A_{m-j+1} - A_{(m-j+1)})(t,\xi) = \sum_{|\gamma| \le m-j} a_{m-j+1,\gamma}(t)\xi^{\gamma}.$$
(8)

For j = m, the condition (7) is the condition on the low order terms coming from the coefficient in front of D_t^{m-1} , in which case $A_1 - A_{(1)}$ is independent of ξ , and assumption (7) should read as

$$|(A_1 - A_{(1)})(t,\xi)| \le C, \quad t \in [0,T],$$

which will be automatically satisfied due to the boundedness of A_1 in *t*. In Sect. 2 we will give examples of the condition (7). We will also show in treating the case m = 3 that from the point of view of the desired energy inequality for (1) the assumption (7) is rather natural.

For a better understanding of the right hand side of the condition (7), let us simplify it in the case when all the characteristic roots are nonnegative, i.e. when $\lambda_l \ge 0$ for all l = 1, ..., m. In this situation, using (3) we see that the right hand side of (7) can be replaced by the coefficient of τ^j in (3) and, therefore, condition (7) becomes

$$|(A_{m-j+1} - A_{(m-j+1)})(t,\xi)| \le C|A_{(m-j)}(t,\xi)|$$

for all $t \in [0, T]$, j = 1, ..., m, and for ξ away from 0.

In the Appendix in Sect. 6 we will show that also in general, the Levi conditions (7) can be expressed entirely in terms of the coefficients of the Eq. (1).

We are now ready to formulate the well-posedness results. Part (i) of Theorem 2 is the extension of Theorem 1 in [13] to any space-dimension and to equations with low order terms. In the sequel $\mathcal{D}'_{(s)}(\mathbb{R}^n)(\mathcal{E}'_{(s)}(\mathbb{R}^n))$ denotes the space of Gevrey Beurling (compactly supported) ultradistributions. For the relevant details on these spaces of ultradistributions and their characterisations, with their appearance in the analysis of weakly hyperbolic equations, we refer to our paper [11], where these have been applied to the low (Hölder) regularity constant multiplicities case.

Theorem 2 Let $n \ge 1$. If the coefficients satisfy $A_j(\cdot, \xi) \in C([0, T])$ and $A_{(j)}(\cdot, \xi) \in C^{\infty}([0, T])$ for all ξ and j = 1, ..., m, the characteristic roots are real and satisfy (6), and the low order terms satisfy (7), then the Cauchy problem (1) is well-posed in any Gevrey space. More precisely, for $A_j(\cdot, \xi) \in C([0, T])$, we have:

(i) if $A_{(j)}(\cdot, \xi) \in C^k([0, T])$ for some $k \ge 2$ and $g_j \in G^s(\mathbb{R}^n)$ for j = 1, ..., m, then there exists a unique solution $u \in C^m([0, T]; G^s(\mathbb{R}^n))$ provided that

$$1 \le s < 1 + \frac{k}{2(m-1)};$$

(ii) if $A_{(j)}(\cdot, \xi) \in C^k([0, T])$ for some $k \ge 2$ and $g_j \in \mathcal{E}'_{(s)}(\mathbb{R}^n)$ for j = 1, ..., m, then there exists a unique solution $u \in C^m([0, T]; \mathcal{D}'_{(s)}(\mathbb{R}^n))$ provided that

$$1 \le s \le 1 + \frac{k}{2(m-1)}$$

In the case of analytic coefficients, we have C^{∞} and distributional well-posedness.

Theorem 3 If $A_j(\cdot, \xi) \in C([0, T])$ and the coefficients $A_{(j)}(\cdot, \xi)$ are analytic on [0, T] for all ξ and j = 1, ..., m, the characteristic roots are real and satisfy (6), and the lower order terms fulfil the conditions (7), then the Cauchy problem (1) is C^{∞} and distributionally well-posed.

By $W^{\infty,m}$ we denote the Sobolev space of functions having *m* derivatives in L^{∞} . In the case of discontinuous but bounded lower order terms we have the following:

Theorem 4 (i) Assume the conditions of Theorem 2, with $A_j(\cdot, \xi) \in C([0, T])$ replaced by $A_j(\cdot, \xi) \in L^{\infty}([0, T]), j = 1, ..., m$. Then the statement remains true provided that we replace the conclusion

$$u \in C^m([0, T]; G^s(\mathbb{R}^n))$$

by

$$u \in C^{m-1}([0, T]; G^{s}(\mathbb{R}^{n})) \cap W^{\infty, m}([0, T]; G^{s}(\mathbb{R}^{n})).$$

(ii) Assume the conditions of Theorem 3 with $A_j(\cdot,\xi) \in C([0,T])$ replaced by $A_j(\cdot,\xi) \in L^{\infty}([0,T]), j = 1, ..., m$. Then the C^{∞} well-posedness remains true provided that we replace the conclusion

$$u \in C^m([0, T]; C^{\infty}(\mathbb{R}^n))$$

by

$$u \in C^{m-1}([0, T]; C^{\infty}(\mathbb{R}^n)) \cap W^{\infty, m}([0, T]; C^{\infty}(\mathbb{R}^n)).$$

Similar conclusions remain true in the ultradistributional/distributional settings as well.

In Remark 2 we show that in certain cases the results can be improved, but the Gevrey index may change.

We refer to Remark 1 for a brief discussion of the strictly hyperbolic case. In this case, even in the situation of the lower regularity of coefficients (C^1) , one can analyse the global behaviour of solutions with respect to time (see [14]). The cases of constant coefficients and systems with controlled oscillations have been treated in [17,18], respectively.

Finally, we describe the contents of the sections in more details.

Section 2 collects some motivating examples of applications of our results. In Sect. 3 we recall the required facts about the quasi-symmetriser and in Sect. 4 we use it to derive the energy estimate for the solutions of the hyperbolic system in Sylvester form corresponding to the Cauchy problem (1). The estimate on the part of the energy corresponding to lower order terms is given in Sect. 5. In Sect. 6 we prove Theorems 2, 3 and 4 and we end the paper with a final remark on the Levi conditions (7).

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2 Examples

Let us first give an example of the Levi conditions (7) for the equations of third order, m = 3. In this case (7) become

$$\begin{aligned} |A_3 - A_{(3)}|^2 &\leq C(\lambda_1^2 \lambda_2^2 + \lambda_2^2 \lambda_3^2 + \lambda_3^2 \lambda_1^2), \\ |A_2 - A_{(2)}|^2 &\leq C((\lambda_1 + \lambda_2)^2 + (\lambda_2 + \lambda_3)^2 + (\lambda_3 + \lambda_1)^2), \\ |A_1 - A_{(1)}|^2 &\leq C, \end{aligned}$$
(9)

for some C > 0. It is convenient in certain applications, whenever possible, to write conditions (6) and (7) in terms of the coefficients of the equation. Such analysis for (1) has been recently carried out by Jannelli and Taglialatela [12]. In Example 3 below we will give an example of the meaning of conditions (9).

For the future technicality, similarly to (9), we may also use an equivalent formulation of (7) as

$$|(A_{m-j+1} - A_{(m-j+1)})(t,\xi)|^{2} \leq C \sum_{i=1}^{m} \left| \sum_{\substack{1 \le l_{1} < \dots < l_{m-j} \le m \\ l_{h} \neq i \ \forall h}} \lambda_{l_{1}}(t,\xi) \cdots \lambda_{l_{m-j}}(t,\xi) \right|^{2}.$$
 (10)

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Condition (6) can be often reformulated in terms of the discriminant of (1) defined by $\Delta(t,\xi) = \prod_{i < j} (\lambda_i(t,\xi) - \lambda_j(t,\xi))^2$. Thus, for m = 2, n = 1, and the equation

$$\partial_t^2 u + a_1(t)\partial_t \partial_x u + a_2(t)\partial_x^2 u = 0,$$

condition (6) is equivalent to the existence of c > 0 such that $\Delta(t) \ge ca_1(t)^2$, where $\Delta(t, \xi) = \Delta(t)\xi$ and $\Delta(t) = a_1^2(t) - 4a_2(t) \ge 0$ is the condition of the hyperbolicity. For m = 3, n = 1, and the equation

$$\partial_t^3 u + a_1(t)\partial_x \partial_t^2 u + a_2(t)\partial_x^2 \partial_t u + a_3(t)\partial_x^3 u = 0,$$

following [13], we have $\Delta(t,\xi) = \Delta(t)\xi$, with $\Delta(t) = -4a_2^3 - 27a_3^2 + a_1^2a_2^2 - 4a_1^3a_3 + 18a_1a_2a_3 \ge 0$, and (6) is equivalent to $\Delta(t) \ge c(a_1(t)a_2(t) - 9a_3(t))^2$.

Since the hyperbolic equations above have homogeneous symbols, the coefficients are real. We refer to Colombini–Orrú [7] and Kinoshita–Spagnolo [13] for more examples of equations without lower order terms in one dimension n = 1.

We now give more examples, which correspond to the new possibility, ensured by Theorems 2 and 3, to consider equations with lower order terms and equations in higher dimensions $n \ge 1$. However, we note that in the case of second order equations much stronger results can be obtained, see Remark 3, so Examples 1 and 2 serve primarily at demonstrating the meaning of conditions (7) only.

Example 1

As a first example we consider the second order equation

$$D_t^2 u + a_{2,2}(t) D_x^2 u + a_{2,1}(t) D_x u + a_{1,0}(t) D_t u + a_{2,0}(t) u = 0,$$

for $t \in [0, T]$ and $x \in \mathbb{R}$. Assume $a_{2,2}(t)$ is real and $a_{2,2}(t) \le 0$. The condition (6) is trivially satisfied by the roots

$$\lambda_1(t,\xi) = -\sqrt{-a_{2,2}(t)}|\xi|,$$

$$\lambda_2(t,\xi) = +\sqrt{-a_{2,2}(t)}|\xi|.$$

The well-posedness results of Sect. 1 are obtained under the conditions (7) on the lower order terms. In this case (7) means that the coefficient $a_{1,0}(t)$ is bounded on [0, T] and that there exists a constant c > 0 such that $|a_{2,1}(t)\xi + a_{2,0}(t)|^2 \le -ca_{2,2}(t)\xi^2$ for all

 $t \in [0, T]$ and for ξ away from 0. Note that this last condition holds if $|a_{2,1}(t)|^2 + |a_{2,0}(t)|^2 \leq -c'a_{2,2}(t)$ for some c' > 0 on the *t*-interval [0, T].

Take now the general second order equation

$$D_t^2 u + a_{1,1}(t) D_x D_t u + a_{2,2}(t) D_x^2 u + a_{2,1}(t) D_x u + a_{1,0}(t) D_t u + a_{2,0}(t) u = 0.$$

As observed above condition (6) coincides with the bound from below

$$\Delta(t) = a_{1,1}^2(t) - 4a_{2,2}(t) \ge c_0 a_{1,1}^2(t),$$

valid for some $c_0 > 0$ on [0, T] ([13, (15)]). Here $a_{1,1}, a_{2,2}$ are assumed real. The conditions (7) on the lower order terms are of the type $|a_{1,0}(t)| \le c_1$ and $|a_{2,1}(t)\xi + a_{2,0}(t)|^2 \le c_2(a_{1,1}^2(t) - 2a_{2,2}(t))\xi^2$ for all $t \in [0, T]$ and for ξ away from 0.

Example 2

The equation

$$D_t^2 u + \sum_{j=1}^n a_{1,j}(t) D_{x_j} D_t u + a_2(t) \sum_{j=1}^n D_{x_j}^2 u + \sum_{j=1}^n b_j(t) D_{x_j} u + b(t) D_t u + d(t) u = 0$$

is an *n*-dimensional version of the previous example, with real $a_{1,j}$ and a_2 . The condition (6) is trivially satisfied when $a_2(t) \le 0$. The conditions (7) on the lower order terms are as follows:

$$|b(t)| \le c,$$

$$\left|\sum_{j=1}^{n} b_j(t)\xi_j + d(t)\right|^2 \le c \left[\left(\sum_{j=1}^{n} a_{1,j}(t)\xi_j\right)^2 - 2a_2(t)|\xi|^2\right],$$

for $t \in [0, T]$ and ξ away from 0.

Example 3

We finally give an example of a higher order equation. Let

$$D_t^3 u - (a+b+c)D_x D_t^2 u + (ab+ac+bc)D_x^2 D_t u - abc D_x^3 u + \sum_{l<3} a_{3,l}(t)D_x^l u + \sum_{l<2} a_{2,l}(t)D_x^l D_t u + a_{1,0}(t)D_t^2 u = 0,$$

where a(t), b(t) and c(t) are real-valued functions with b and c bounded above and from below by a (for instance, $a(t)/4 \le b(t) \le a(t)/2$ and $a(t)/16 \le c(t) \le a(t)/8$ for all $t \in [0, T]$). It follows that condition (6) on the roots $\lambda_1(t, \xi) = a(t)\xi$, $\lambda_2(t, \xi) = b(t)\xi$ and $\lambda_3(t, \xi) = c(t)\xi$ is fulfilled on [0, T] for all $\xi \in \mathbb{R}$. The Levi conditions (7) on the lower order terms are of the following type:

$$\begin{aligned} |a_{3,2}(t)\xi^{2} + a_{3,1}(t)\xi + a_{3,0}(t)|^{2} &\leq c \, a^{4}(t)\xi^{4}, \\ |a_{2,1}(t)\xi + a_{2,0}(t)|^{2} &\leq c \, a^{2}(t)\xi^{2}, \\ |a_{1,0}(t)|^{2} &\leq c, \end{aligned}$$

for $t \in [0, T]$ and ξ away from 0.

3 The quasi-symmetriser

We begin by recalling a few facts concerning the quasi-symmetriser. For more details see [10,13]. Note that for $m \times m$ matrices A_1 and A_2 the notation $A_1 \leq A_2$ means $(A_1v, v) \leq (A_2v, v)$ for all $v \in \mathbb{C}^m$ with (\cdot, \cdot) the scalar product in \mathbb{C}^m . Let $A(\lambda)$ be the $m \times m$ Sylvester matrix with real eigenvalues λ_l , i.e.,

$$A(\lambda) = \begin{pmatrix} 0 & 1 & 0 & \dots & 0 \\ 0 & 0 & 1 & \dots & 0 \\ \dots & \dots & \dots & \dots & 1 \\ -\sigma_m^{(m)}(\lambda) & -\sigma_{m-1}^{(m)}(\lambda) & \dots & \dots & -\sigma_1^{(m)}(\lambda) \end{pmatrix},$$

where

$$\sigma_h^{(m)}(\lambda) = (-1)^h \sum_{1 \le i_1 < \dots < i_h \le m} \lambda_{i_1} \dots \lambda_{i_h}$$

for all $1 \leq h \leq m$. In the sequel we make use of the following notations: \mathcal{P}_m for the class of permutations of $\{1, \ldots, m\}, \lambda_{\rho} = (\lambda_{\rho_1}, \ldots, \lambda_{\rho_m})$ with $\lambda \in \mathbb{R}^m$ and $\rho \in \mathcal{P}_m, \pi_i \lambda = (\lambda_1, \ldots, \lambda_{i-1}, \lambda_{i+1}, \ldots, \lambda_m)$ and $\lambda' = \pi_m \lambda = (\lambda_1, \ldots, \lambda_{m-1})$. Following Sect. 4 in [13] we have that the quasi-symmetriser is the Hermitian matrix

$$Q_{\varepsilon}^{(m)}(\lambda) = \sum_{\rho \in \mathcal{P}_m} P_{\varepsilon}^{(m)}(\lambda_{\rho})^* P_{\varepsilon}^{(m)}(\lambda_{\rho}),$$

where $\varepsilon \in (0, 1]$, $P_{\varepsilon}^{(m)}(\lambda) = H_{\varepsilon}^{(m)}P^{(m)}(\lambda)$, $H_{\varepsilon}^{(m)} = \text{diag}\{\varepsilon^{m-1}, \ldots, \varepsilon, 1\}$ and the matrix $P^{(m)}(\lambda)$ is defined inductively by $P^{(1)}(\lambda) = 1$ and

$$P^{(m)}(\lambda) = \begin{pmatrix} & & & 0 \\ P^{(m-1)}(\lambda') & & \vdots \\ \sigma_{m-1}^{(m-1)}(\lambda') & \dots & \sigma_1^{(m-1)}(\lambda') & 1 \end{pmatrix}.$$

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Note that $P^{(m)}(\lambda)$ is depending only on λ' . Finally, let $W_i^{(m)}(\lambda)$ denote the row vector

$$(\sigma_{m-1}^{(m-1)}(\pi_i\lambda), \ldots, \sigma_1^{(m-1)}(\pi_i\lambda), 1), \quad 1 \le i \le m,$$

and let $\mathcal{W}^{(m)}(\lambda)$ be the matrix with row vectors $W_i^{(m)}$. The following proposition collects the main properties of the quasi-symmetriser $Q_{\varepsilon}^{(m)}(\lambda)$. For a detailed proof we refer the reader to Propositions 1 and 2 in [13] and to Proposition 1 in [10].

Proposition 1 (i) The quasi-symmetriser $Q_{\varepsilon}^{(m)}(\lambda)$ can be written as

$$Q_0^{(m)}(\lambda) + \varepsilon^2 Q_1^{(m)}(\lambda) + \dots + \varepsilon^{2(m-1)} Q_{m-1}^{(m)}(\lambda)$$

where the matrices $Q_i^{(m)}(\lambda)$, i = 1, ..., m - 1, are nonnegative and Hermitian with entries being symmetric polynomials in $\lambda_1, ..., \lambda_m$.

(ii) There exists a function $C_m(\lambda)$ bounded for bounded $|\lambda|$ such that

$$C_m(\lambda)^{-1}\varepsilon^{2(m-1)}I \leq Q_{\varepsilon}^{(m)}(\lambda) \leq C_m(\lambda)I.$$

(iii) We have

$$-C_m(\lambda)\varepsilon Q_{\varepsilon}^{(m)}(\lambda) \leq Q_{\varepsilon}^{(m)}(\lambda)A(\lambda) - A(\lambda)^* Q_{\varepsilon}^{(m)}(\lambda) \leq C_m(\lambda)\varepsilon Q_{\varepsilon}^{(m)}(\lambda).$$

(iv) For any $(m-1) \times (m-1)$ matrix T let T^{\sharp} denote the $m \times m$ matrix

$$\begin{pmatrix} T & 0 \\ 0 & 0 \end{pmatrix}.$$

Then, $Q_{\varepsilon}^{(m)}(\lambda) = Q_0^{(m)}(\lambda) + \varepsilon^2 \sum_{i=1}^m Q_{\varepsilon}^{(m-1)}(\pi_i \lambda)^{\sharp}$. (v) We have

$$Q_0^{(m)}(\lambda) = (m-1)! \mathcal{W}^{(m)}(\lambda)^* \mathcal{W}^{(m)}(\lambda).$$

(vi) We have

$$\det \mathcal{Q}_0^{(m)}(\lambda) = (m-1)! \prod_{1 \le i < j \le m} (\lambda_i - \lambda_j)^2.$$

(vii) There exists a constant C_m such that

$$q_{0,11}^{(m)}(\lambda)\cdots q_{0,mm}^{(m)}(\lambda) \leq C_m \prod_{1\leq i< j\leq m} (\lambda_i^2+\lambda_j^2).$$

We finally recall that a family $\{Q_{\alpha}\}$ of nonnegative Hermitian matrices is called *nearly* diagonal if there exists a positive constant c_0 such that

$$Q_{\alpha} \geq c_0 \operatorname{diag} Q_{\alpha}$$

for all α , with diag $Q_{\alpha} = \text{diag}\{q_{\alpha,11}, \ldots, q_{\alpha,mm}\}$. The following linear algebra result is proven in [13, Lemma 1].

Lemma 1 Let $\{Q_{\alpha}\}$ be a family of nonnegative Hermitian $m \times m$ matrices such that det $Q_{\alpha} > 0$ and

$$\det Q_{\alpha} \ge c \, q_{\alpha, 11} q_{\alpha, 22} \cdots q_{\alpha, mm}$$

for a certain constant c > 0 independent of α . Then,

$$Q_{\alpha} \ge c m^{1-m} \operatorname{diag} Q_{\alpha}$$

for all α , i.e., the family $\{Q_{\alpha}\}$ is nearly diagonal.

Lemma 1 is employed to prove that the family $Q_{\varepsilon}^{(m)}(\lambda)$ of quasi-symmetrisers defined above is nearly diagonal when λ belongs to a suitable set. The following statement is proven in [13, Proposition 3].

Proposition 2 For any M > 0 define the set

$$\mathcal{S}_M = \{\lambda \in \mathbb{R}^m : \lambda_i^2 + \lambda_j^2 \le M(\lambda_i - \lambda_j)^2, \quad 1 \le i < j \le m\}.$$

Then the family of matrices $\{Q_{\varepsilon}^{(m)}(\lambda) : 0 < \varepsilon \leq 1, \lambda \in S_M\}$ is nearly diagonal.

We conclude this section with a result on nearly diagonal matrices depending on 3 parameters (i.e. ε , t, ξ) which will be crucial in the next section. Note that this is a straightforward extension of Lemma 2 in [13] valid for 2 parameter (i.e. ε , t) dependent matrices.

Lemma 2 Let $\{Q_{\varepsilon}^{(m)}(t,\xi): 0 < \varepsilon \leq 1, 0 \leq t \leq T, \xi \in \mathbb{R}^n\}$ be a nearly diagonal family of coercive Hermitian matrices of class C^k in $t, k \geq 1$. Then, there exists a constant $C_T > 0$ such that for any continuous function $V : [0, T] \times \mathbb{R}^n \to \mathbb{C}^m$ we have

$$\int_{0}^{T} \frac{|(\partial_{t} Q_{\varepsilon}^{(m)}(t,\xi) V(t,\xi), V(t,\xi))|}{(Q_{\varepsilon}^{(m)}(t,\xi) V(t,\xi), V(t,\xi))^{1-1/k} |V(t,\xi)|^{2/k}} dt \leq C_{T} \|Q_{\varepsilon}^{(m)}(\cdot,\xi)\|_{C^{k}([0,T])}^{1/k}$$

for all $\xi \in \mathbb{R}^n$.

4 Reduction to a first order system and energy estimate

We now go back to the Cauchy problem (1) and perform a reduction of the *m*-order equation to a first order system as in [19]. Let $\langle D_x \rangle$ be the pseudo-differential operator with symbol $\langle \xi \rangle = (1 + |\xi|^2)^{\frac{1}{2}}$. The transformation

$$u_j = D_t^{j-1} \langle D_x \rangle^{m-j} u,$$

with j = 1, ..., m, makes the Cauchy problem (1) equivalent to the following system

$$D_t \begin{pmatrix} u_1 \\ \cdot \\ \cdot \\ u_m \end{pmatrix} = \begin{pmatrix} 0 & \langle D_x \rangle & 0 & \dots & 0 \\ 0 & 0 & \langle D_x \rangle & \dots & 0 \\ \dots & \dots & \dots & \dots & \langle D_x \rangle \\ b_1 & b_2 & \dots & \dots & b_m \end{pmatrix} \begin{pmatrix} u_1 \\ \cdot \\ \cdot \\ u_m \end{pmatrix},$$
(11)

where

$$b_j = -A_{m-j+1}(t, D_x) \langle D_x \rangle^{j-m},$$

with initial condition

$$u_j|_{t=0} = \langle D_x \rangle^{m-j} g_j, \quad j = 1, \dots, m.$$
 (12)

The matrix in (11) can be written as $\mathbb{A}_1 + B$ with

$$\mathbb{A}_{1} = \begin{pmatrix} 0 & \langle D_{x} \rangle & 0 & \dots & 0 \\ 0 & 0 & \langle D_{x} \rangle & \dots & 0 \\ \dots & \dots & \dots & \dots & \langle D_{x} \rangle \\ b_{(1)} & b_{(2)} & \dots & \dots & b_{(m)} \end{pmatrix},$$

where $b_{(j)} = -A_{(m-j+1)}(t, D_x) \langle D_x \rangle^{j-m}$ is the principal part of the operator $b_j = -A_{m-j+1}(t, D_x) \langle D_x \rangle^{j-m}$ and

$$B = \begin{pmatrix} 0 & 0 & 0 & \dots & 0 \\ 0 & 0 & 0 & \dots & 0 \\ \dots & \dots & \dots & \dots & 0 \\ b_1 - b_{(1)} & b_2 - b_{(2)} & \dots & \dots & b_m - b_{(m)} \end{pmatrix}$$

By Fourier transforming both sides of (11) in x we obtain the system

$$D_t V = \mathbb{A}_1(t,\xi)V + B(t,\xi)V,$$

$$V|_{t=0}(\xi) = V_0(\xi),$$
(13)

where V is the *m*-column with entries $v_j = \hat{u}_j$, V_0 is the *m*-column with entries $v_{0,j} = \langle \xi \rangle^{m-j} \hat{g}_j$ and

$$\mathbb{A}_{1}(t,\xi) = \begin{pmatrix} 0 & \langle \xi \rangle & 0 & \dots & 0 \\ 0 & 0 & \langle \xi \rangle & \dots & 0 \\ \dots & \dots & \dots & \dots & \langle \xi \rangle \\ b_{(1)}(t,\xi) & b_{(2)}(t,\xi) & \dots & \dots & b_{(m)}(t,\xi) \end{pmatrix},$$

$$b_{(j)}(t,\xi) = -A_{(m-j+1)}(t,\xi)\langle\xi\rangle^{j-m},$$
(14)

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$$B(t,\xi) = \begin{pmatrix} 0 & 0 & 0 & \dots & 0 \\ 0 & 0 & 0 & \dots & 0 \\ \dots & \dots & \dots & \dots & 0 \\ (b_1 - b_{(1)})(t,\xi) & \dots & \dots & \dots & (b_m - b_{(m)})(t,\xi) \end{pmatrix},$$
$$(b_j - b_{(j)})(t,\xi) = -(A_{m-j+1} - A_{(m-j+1)})(t,\xi)\langle\xi\rangle^{j-m}.$$

From now on we will concentrate on the system (13) and on the matrix

$$A(t,\xi) := \langle \xi \rangle^{-1} \mathbb{A}_1(t,\xi)$$

for which we will construct a quasi-symmetriser. Note that the eigenvalues of the matrix \mathbb{A}_1 are exactly the roots $\lambda_l(t,\xi)$, $l = 1, \ldots, m$. It is clear that the condition (6) holds for the eigenvalues $\langle \xi \rangle^{-1} \lambda_l(t,\xi)$ of the 0-order matrix $A(t,\xi)$ as well.Let us define the energy

$$E_{\varepsilon}(t,\xi) = (Q_{\varepsilon}^{(m)}(t,\xi)V(t,\xi), V(t,\xi)).$$

We have

$$\begin{aligned} \partial_t E_{\varepsilon}(t,\xi) &= (\partial_t Q_{\varepsilon}^{(m)} V, V) + i(Q_{\varepsilon}^{(m)} D_t V, V) - i(Q_{\varepsilon}^{(m)} V, D_t V) \\ &= (\partial_t Q_{\varepsilon}^{(m)} V, V) + i(Q_{\varepsilon}^{(m)} (\mathbb{A}_1 V + BV), V) - i(Q_{\varepsilon}^{(m)} V, \mathbb{A}_1 V + BV) \\ &= (\partial_t Q_{\varepsilon}^{(m)} V, V) + i\langle\xi\rangle ((Q_{\varepsilon}^{(m)} A - A^* Q_{\varepsilon}^{(m)}) V, V) \\ &+ i((Q_{\varepsilon}^{(m)} B - B^* Q_{\varepsilon}^{(m)}) V, V). \end{aligned}$$

It follows that

$$\partial_{t} E_{\varepsilon}(t,\xi) \leq \frac{|(\partial_{t} Q_{\varepsilon}^{(m)} V, V)| E_{\varepsilon}}{(Q_{\varepsilon}^{(m)}(t,\xi) V(t,\xi), V(t,\xi))} + |\langle \xi \rangle ((Q_{\varepsilon}^{(m)} A - A^{*} Q_{\varepsilon}^{(m)}) V, V)| + |((Q_{\varepsilon}^{(m)} B - B^{*} Q_{\varepsilon}^{(m)}) V, V)|.$$
(15)

We recall that from Proposition 1, $Q_{\varepsilon}^{(m)}(t, \xi)$ is a family of smooth nonnegative Hermitian matrices such that

$$Q_{\varepsilon}^{(m)}(t,\xi) = Q_0^{(m)}(t,\xi) + \varepsilon^2 Q_1^{(m)}(t,\xi) + \dots + \varepsilon^{2(m-1)} Q_{m-1}^{(m)}(t,\xi).$$
(16)

In addition there exists a constant $C_m > 0$ such that for all $t \in [0, T], \xi \in \mathbb{R}^n$ and $\varepsilon \in (0, 1]$ the following estimates hold uniformly in V:

$$C_m^{-1}\varepsilon^{2(m-1)}|V|^2 \le (Q_\varepsilon^{(m)}(t,\xi)V,V) \le C_m|V|^2,$$
(17)

$$|((Q_{\varepsilon}^{(m)}A - A^*Q_{\varepsilon}^{(m)})(t,\xi)V,V)| \le C_m \varepsilon(Q_{\varepsilon}^{(m)}(t,\xi)V,V).$$
(18)

Finally, condition (6) and Proposition 2 ensure that the family

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$$\{Q_{\varepsilon}^{(m)}(t,\xi): \varepsilon \in (0,1], t \in [0,T], \xi \in \mathbb{R}^n\}$$

is nearly diagonal.

In the sequel we assume that the coefficients a_j in the Eq. (1) are of class C^k , or in other words that the matrix $A(t, \xi)$ has entries of class C^k in $t \in [0, T]$. It follows by construction that the quasi-symmetriser has the same regularity property. We now estimate the three terms of the right hand side of (15).

4.1 First term

We write $\frac{|(\partial_t Q_{\varepsilon}^{(m)}V, V)|}{(Q_{\varepsilon}^{(m)}(t,\xi)V(t,\xi), V(t,\xi))}$ as

$$\frac{|(\partial_t Q_{\varepsilon}^{(m)}V, V)|}{(Q_{\varepsilon}^{(m)}(t,\xi)V(t,\xi), V(t,\xi))^{1-1/k}(Q_{\varepsilon}^{(m)}(t,\xi)V(t,\xi), V(t,\xi))^{1/k}}$$

From (17) we have

$$\begin{aligned} \frac{|(\partial_t Q_{\varepsilon}^{(m)}V, V)|}{(Q_{\varepsilon}^{(m)}(t, \xi)V(t, \xi), V(t, \xi))} \\ &\leq \frac{|(\partial_t Q_{\varepsilon}^{(m)}V, V)|}{(Q_{\varepsilon}^{(m)}(t, \xi)V(t, \xi), V(t, \xi))^{1-1/k}(C_m^{-1}\varepsilon^{2(m-1)}|V|^2)^{1/k}} \\ &\leq C_m^{1/k}\varepsilon^{-2(m-1)/k}\frac{|(\partial_t Q_{\varepsilon}^{(m)}V, V)|}{(Q_{\varepsilon}^{(m)}(t, \xi)V(t, \xi), V(t, \xi))^{1-1/k}|V|^{2/k}}\end{aligned}$$

An application of Lemma 2 yields the estimate

$$\int_{0}^{T} \frac{|(\partial_{t} Q_{\varepsilon}^{(m)} V, V)|}{(Q_{\varepsilon}^{(m)}(t, \xi) V(t, \xi), V(t, \xi))} dt$$
$$\leq C_{m}^{1/k} \varepsilon^{-2(m-1)/k} C_{T} \sup_{\xi \in \mathbb{R}^{n}} \|Q_{\varepsilon}(\cdot, \xi)\|_{C^{k}([0,T])}^{1/k}$$
$$\leq C_{1} \varepsilon^{-2(m-1)/k},$$

for all $\varepsilon \in (0, 1]$. Setting $\frac{|(\partial_t Q_{\varepsilon}^{(m)}V, V)|}{(Q_{\varepsilon}^{(m)}(t, \xi)V(t, \xi), V(t, \xi))} = K_{\varepsilon}(t, \xi)$ we conclude that

$$\frac{|(\partial_t Q_{\varepsilon}^{(m)} V, V)| E_{\varepsilon}}{(Q_{\varepsilon}^{(m)}(t, \xi) V(t, \xi), V(t, \xi))} = K_{\varepsilon}(t, \xi) E_{\varepsilon},$$
(19)

with

$$\int_{0}^{T} K_{\varepsilon}(t,\xi) dt \le C_1 \varepsilon^{-2(m-1)/k}.$$
(20)

4.2 Second term

From the property (18) we have that

$$|\langle \xi \rangle ((Q_{\varepsilon}^{(m)}A - A^*Q_{\varepsilon}^{(m)})V, V)| \le C_m \varepsilon \langle \xi \rangle (Q_{\varepsilon}^{(m)}(t,\xi)V, V) \le C_2 \varepsilon \langle \xi \rangle E_{\varepsilon}.$$

4.3 Third term

We now concentrate on

$$((Q_{\varepsilon}^{(m)}B - B^*Q_{\varepsilon}^{(m)})V, V),$$

which is the main task in this paper. By Proposition 1(iv) and the definition of the matrix $B(t, \xi)$ we have that

$$((Q_{\varepsilon}^{(m)}B - B^*Q_{\varepsilon}^{(m)})V, V) = ((Q_0^{(m)}B - B^*Q_0^{(m)})V, V) + \varepsilon^2 \sum_{i=1}^m ((Q_{\varepsilon}^{(m-1)}(\pi_i\lambda)^{\sharp}B - B^*Q_{\varepsilon}^{(m-1)}(\pi_i\lambda)^{\sharp})V, V),$$

where we notice that $(Q_{\varepsilon}^{(m-1)}(\pi_i\lambda)^{\sharp}B - B^*Q_{\varepsilon}^{(m-1)}(\pi_i\lambda)^{\sharp}) = 0$ due to the structure of zeros in *B* and in $Q_{\varepsilon}^{(m-1)}(\pi_i\lambda)^{\sharp}$. Hence

$$((Q_{\varepsilon}^{(m)}B - B^*Q_{\varepsilon}^{(m)})V, V) = ((Q_0^{(m)}B - B^*Q_0^{(m)})V, V).$$

Note that from Proposition 1(i) we have that $(Q_0^{(m)}V, V) \leq E_{\varepsilon}$. In the next section we will show that the conditions on *B* corresponding to (7) imply that

$$|((Q_0^{(m)}B - B^*Q_0^{(m)})V, V)| \le C_3(Q_0^{(m)}V, V) \le C_3E_{\varepsilon},$$
(21)

for some constant $C_3 > 0$ independent of $t \in [0, T], \xi \in \mathbb{R}^n$ and $V \in \mathbb{C}^m$.

Remark 1 Note that condition (6) is trivially satisfied when the roots are distinct, i.e. in the strictly hyperbolic case. It follows that the family $\{Q_{\varepsilon}(\lambda)\}$ of quasi-symmetrisers is nearly diagonal and, therefore, there exists a constant $c_0 > 0$ such that $Q_0^{(m)} \ge c_0 \text{diag } Q_0^{(m)}$. This means that

$$(Q_0^{(m)}V, V) \ge c_0 \sum_{i=1}^m q_{0,ii} |V_i|^2$$

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holds for all $V \in \mathbb{C}^m$. From the hypothesis of strict hyperbolicity it easily follows that

$$\inf_{t\in[0,T],|\xi|\geq R,i=1,...,m}q_{0,ii}(t,\xi)>0,$$

for any R > 0. This bound from below implies

$$(Q_0^{(m)}(t,\xi)V,V) \ge c_0'|V|^2$$
(22)

for $t \in [0, T]$ and $|\xi| \ge R$ and hence the estimate

$$|((Q_0^{(m)}B - B^*Q_0^{(m)})V, V)| \le C_3(Q_0V, V)$$

holds trivially in the strictly hyperbolic case for any lower order term *B* (for our purposes it will not be restrictive to assume $|\xi| \ge R$). Concluding, when the roots λ_i are distinct the Gevrey and ultradistributional well-posedness results in Theorems 2 and 3 can be stated without additional conditions on the lower order terms. Strictly hyperbolic equations under low regularity (Hölder C^{α} , $0 < \alpha < 1$) of the coefficients have been analysed by the authors in [11, Case 3], to which we refer for general statements on the Gevrey and ultradistributional well-posedness in this setting.

5 Estimates for the lower order terms

We begin by rewriting $((Q_0^{(m)}B - B^*Q_0^{(m)})V, V)$ in terms of the matrix $\mathcal{W} = \mathcal{W}^{(m)}$. From Proposition 1(v) we have

$$((Q_0^{(m)}B - B^*Q_0^{(m)})V, V) = (m-1)!((WBV, WV) - (WV, WBV))$$

= 2i(m-1)!\(WBV, WV).

It follows that

$$|((Q_0^{(m)}B - B^*Q_0^{(m)})V, V)| \le 2(m-1)!|WBV||WV|.$$

Since

$$(Q_0V, V) = (m-1)!|WV|^2$$

we have that if

$$|\mathcal{W}BV| \le C|\mathcal{W}V| \tag{23}$$

for some constant C > 0 independent of t, ξ and V, then the condition (21) will hold. It is our task to show that the condition (7) on the matrix B of the lower order terms implies the estimate (23). Before dealing with the general case of $Bm \times m$ -matrix, let us consider the instructive case m = 3, which will illustrate the general argument in a simplified setting. In the sequel, for f and g real-valued functions (in the variable y) we write $f(y) \prec g(y)$ if there exists a constant C > 0 such that $f(y) \le Cg(y)$ for all y. More precisely, we will set $y = (t, \xi)$ or $y = (t, \xi, V)$.

5.1 The case m = 3

By definition of the row vectors $W_i^{(3)}$, i = 1, 2, 3, we have that

$$\mathcal{W} = \begin{pmatrix} \lambda_2 \lambda_3 & -\lambda_2 & -\lambda_3 & 1 \\ \lambda_3 \lambda_1 & -\lambda_3 & -\lambda_1 & 1 \\ \lambda_1 \lambda_2 & -\lambda_1 & -\lambda_2 & 1 \end{pmatrix},$$

where λ_i , i = 1, 2, 3, are the 0-order normalised roots. Hence, by definition of the matrix *B* setting $b_j - b_{(j)} = B_j$, j = 1, 2, 3, we get

$$\mathcal{W}BV = \begin{pmatrix} B_1V_1 + B_2V_2 + B_3V_3\\ B_1V_1 + B_2V_2 + B_3V_3\\ B_1V_1 + B_2V_2 + B_3V_3 \end{pmatrix}$$

and

$$\mathcal{W}V = \begin{pmatrix} (\lambda_2\lambda_3)V_1 - (\lambda_2 + \lambda_3)V_2 + V_3\\ (\lambda_3\lambda_1)V_1 - (\lambda_3 + \lambda_1)V_2 + V_3\\ (\lambda_1\lambda_2)V_1 - (\lambda_1 + \lambda_2)V_2 + V_3 \end{pmatrix}$$

Thus, instead of working on proving (23) we can work on the equivalent inequality

$$|B_{1}V_{1} + B_{2}V_{2} + B_{3}V_{3}|^{2} \prec |(\lambda_{2}\lambda_{3})V_{1} - (\lambda_{2} + \lambda_{3})V_{2} + V_{3}|^{2} + |(\lambda_{3}\lambda_{1})V_{1} - (\lambda_{3} + \lambda_{1})V_{2} + V_{3}|^{2} + |(\lambda_{1}\lambda_{2})V_{1} - (\lambda_{1} + \lambda_{2})V_{2} + V_{3}|^{2}.$$
(24)

In terms of the coefficients of the matrix B the Levi conditions (9) on the lower order terms can be written as

$$|B_1|^2 \prec \lambda_1^2 \lambda_2^2 + \lambda_2^2 \lambda_3^2 + \lambda_3^2 \lambda_1^2, |B_2|^2 \prec (\lambda_1 + \lambda_2)^2 + (\lambda_2 + \lambda_3)^2 + (\lambda_3 + \lambda_1)^2, |B_3|^2 \prec c.$$
(25)

Under these conditions we now want to prove that (24) holds for all vectors V. We note here that actually for the right hand side of (24) by the triangle inequality we have the upper bound

$$\begin{split} &|(\lambda_{2}\lambda_{3})V_{1} - (\lambda_{2} + \lambda_{3})V_{2} + V_{3}|^{2} + |(\lambda_{3}\lambda_{1})V_{1} - (\lambda_{3} + \lambda_{1})V_{2} + V_{3}|^{2} \\ &+ |(\lambda_{1}\lambda_{2})V_{1} - (\lambda_{1} + \lambda_{2})V_{2} + V_{3}|^{2} \\ &\prec (\lambda_{1}^{2}\lambda_{2}^{2} + \lambda_{2}^{2}\lambda_{3}^{2} + \lambda_{3}^{2}\lambda_{1}^{2})|V_{1}|^{2} + ((\lambda_{1} + \lambda_{2})^{2} + (\lambda_{2} + \lambda_{3})^{2} + (\lambda_{3} + \lambda_{1})^{2})|V_{2}|^{2} + |V_{3}|^{2}, \end{split}$$

in which the right hand side of (25) appears naturally.

Our strategy is to proceed by 3 steps making use of the following partition of \mathbb{R}^3 :

$$\mathbb{R}^3 = \Sigma_1^{\delta_1} \cup \left(\left(\Sigma_1^{\delta_1} \right)^c \cap \Sigma_2^{\delta_2} \right) \cup \left(\left(\Sigma_1^{\delta_1} \right)^c \cap \left(\Sigma_2^{\delta_2} \right)^c \right),$$

where

$$\begin{split} \Sigma_1^{\delta_1} &:= \{ V \in \mathbb{R}^3 : |V_3|^2 + ((\lambda_1 + \lambda_2)^2 + (\lambda_2 + \lambda_3)^2 + (\lambda_3 + \lambda_1)^2) |V_2|^2 \\ &\leq \delta_1 (\lambda_1^2 \lambda_2^2 + \lambda_2^2 \lambda_3^2 + \lambda_3^2 \lambda_1^2) |V_1|^2 \}, \end{split}$$

and

$$\Sigma_2^{\delta_2} := \{ V \in \mathbb{R}^3 : |V_3|^2 \le \delta_2 ((\lambda_1 + \lambda_2)^2 + (\lambda_2 + \lambda_3)^2 + (\lambda_3 + \lambda_1)^2) |V_2|^2 \}.$$

Estimate on $\Sigma_1^{\delta_1}$.

Making use of the conditions (25) we have that

$$|B_{1}V_{1}+B_{2}V_{2}+B_{3}V_{3}|^{2} \prec |B_{1}|^{2}|V_{1}|^{2}+|B_{2}|^{2}|V_{2}|^{2}+|B_{3}|^{2}|V_{3}|^{2} \prec (\lambda_{1}^{2}\lambda_{2}^{2}+\lambda_{2}^{2}\lambda_{3}^{2}+\lambda_{3}^{2}\lambda_{1}^{2})|V_{1}|^{2}+((\lambda_{1}+\lambda_{2})^{2}+(\lambda_{2}+\lambda_{3})^{2}+(\lambda_{3}+\lambda_{1})^{2})|V_{2}|^{2}+|V_{3}|^{2} \prec (\lambda_{1}^{2}\lambda_{2}^{2}+\lambda_{2}^{2}\lambda_{3}^{2}+\lambda_{3}^{2}\lambda_{1}^{2})|V_{1}|^{2}$$
(26)

on $\Sigma_1^{\delta_1}$. Note¹ that we have

$$\begin{aligned} |(\lambda_{2}\lambda_{3})V_{1} - (\lambda_{2} + \lambda_{3})V_{2} + V_{3}|^{2} + |(\lambda_{3}\lambda_{1})V_{1} - (\lambda_{3} + \lambda_{1})V_{2} + V_{3}|^{2} \\ & \succ |(\lambda_{2}\lambda_{3} - \lambda_{3}\lambda_{1})V_{1} - (\lambda_{2} - \lambda_{1})V_{2}|^{2} \\ & \succ (\lambda_{2} - \lambda_{1})^{2}|\lambda_{3}V_{1} - V_{2}|^{2} \succ (\lambda_{1}^{2} + \lambda_{2}^{2})|\lambda_{3}V_{1} - V_{2}|^{2}, \end{aligned}$$

where also in the last line we make use of the condition (6) on the roots λ_i . Hence, by applying this to different combinations of terms, we get

$$\begin{aligned} |(\lambda_{2}\lambda_{3})V_{1} - (\lambda_{2} + \lambda_{3})V_{2} + V_{3}|^{2} \\ + |(\lambda_{3}\lambda_{1})V_{1} - (\lambda_{3} + \lambda_{1})V_{2} + V_{3}|^{2} + |(\lambda_{1}\lambda_{2})V_{1} - (\lambda_{1} + \lambda_{2})V_{2} + V_{3}|^{2} \\ > (\lambda_{1}^{2} + \lambda_{2}^{2})|\lambda_{3}V_{1} - V_{2}|^{2} + (\lambda_{2}^{2} + \lambda_{3}^{2})|\lambda_{1}V_{1} - V_{2}|^{2} + (\lambda_{3}^{2} + \lambda_{1}^{2})|\lambda_{2}V_{1} - V_{2}|^{2} \\ > \lambda_{1}^{2}(|\lambda_{3}V_{1} - V_{2}|^{2} + |\lambda_{2}V_{1} - V_{2}|^{2}) + \lambda_{2}^{2}(|\lambda_{3}V_{1} - V_{2}|^{2} + |\lambda_{1}V_{1} - V_{2}|^{2}) \\ + \lambda_{3}^{2}(|\lambda_{2}V_{1} - V_{2}|^{2} + |\lambda_{1}V_{1} - V_{2}|^{2}) \\ > (\lambda_{1}^{2}(\lambda_{3} - \lambda_{2})^{2} + \lambda_{2}^{2}(\lambda_{3} - \lambda_{1})^{2} + \lambda_{3}^{2}(\lambda_{2} - \lambda_{1})^{2})|V_{1}|^{2} \\ > (\lambda_{1}^{2}\lambda_{2}^{2} + \lambda_{2}^{2}\lambda_{3}^{2} + \lambda_{3}^{2}\lambda_{1}^{2})|V_{1}|^{2}. \end{aligned}$$

$$(27)$$

¹ Using the formula $|z_1|^2 + |z_2|^2 \ge \frac{1}{2}|z_1 - z_2|^2, z_1, z_2 \in \mathbb{C}.$

From the bound from below (27) and the estimate (26) one has that the inequality (24) holds true in the region $\Sigma_1^{\delta_1}$ for all $\delta_1 > 0$. **Estimate on** $(\Sigma_1^{\delta_1})^c \cap \Sigma_2^{\delta_2}$.

We assume from now on that $V \in (\Sigma_1^{\delta_1})^c$ which means that

$$|V_3|^2 + ((\lambda_1 + \lambda_2)^2 + (\lambda_2 + \lambda_3)^2 + (\lambda_3 + \lambda_1)^2)|V_2|^2 > \delta_1(\lambda_1^2 \lambda_2^2 + \lambda_2^2 \lambda_3^2 + \lambda_3^2 \lambda_1^2)|V_1|^2.$$

One immediately has

$$\begin{split} &|B_1V_1 + B_2V_2 + B_3V_3|^2 \prec |B_1|^2 |V_1|^2 + |B_2|^2 |V_2|^2 + |B_3|^2 |V_3|^2 \\ &\prec (\lambda_1^2\lambda_2^2 + \lambda_2^2\lambda_3^2 + \lambda_3^2\lambda_1^2)|V_1|^2 + ((\lambda_1 + \lambda_2)^2 + (\lambda_2 + \lambda_3)^2 + (\lambda_3 + \lambda_1)^2)|V_2|^2 + |V_3|^2 \\ &\prec ((\lambda_1 + \lambda_2)^2 + (\lambda_2 + \lambda_3)^2 + (\lambda_3 + \lambda_1)^2)|V_2|^2 + |V_3|^2. \end{split}$$

More precisely,

$$|B_1V_1 + B_2V_2 + B_3V_3|^2 \prec ((\lambda_1 + \lambda_2)^2 + (\lambda_2 + \lambda_3)^2 + (\lambda_3 + \lambda_1)^2)|V_2|^2 \quad (28)$$

holds for all $V \in (\Sigma_1^{\delta_1})^c \cap \Sigma_2^{\delta_2}$. We estimate the right-hand side of (24) as

$$\begin{split} &|(\lambda_{2}\lambda_{3})V_{1} - (\lambda_{2} + \lambda_{3})V_{2} + V_{3}|^{2} \\ &+ |(\lambda_{3}\lambda_{1})V_{1} - (\lambda_{3} + \lambda_{1})V_{2} + V_{3}|^{2} + |(\lambda_{1}\lambda_{2})V_{1} - (\lambda_{1} + \lambda_{2})V_{2} + V_{3}|^{2} \\ &\succ \gamma_{1}(|(\lambda_{2} + \lambda_{3})V_{2} - V_{3}|^{2} + |(\lambda_{3} + \lambda_{1})V_{2} - V_{3}|^{2} + |(\lambda_{1} + \lambda_{2})V_{2} - V_{3}|^{2}) \\ &- \gamma_{2}(\lambda_{1}^{2}\lambda_{2}^{2} + \lambda_{2}^{2}\lambda_{3}^{2} + \lambda_{3}^{2}\lambda_{1}^{2})|V_{1}|^{2}, \end{split}$$

for some² constants γ_1 , $\gamma_2 > 0$. By using condition (6) we get the estimate

$$\begin{aligned} &(\lambda_2 - \lambda_1)^2 + (\lambda_3 - \lambda_2)^2 + (\lambda_3 - \lambda_1)^2 \ge \frac{2}{M} (\lambda_1^2 + \lambda_2^2 + \lambda_3^2) \\ &\ge \frac{1}{2M} ((\lambda_1 + \lambda_2)^2 + (\lambda_2 + \lambda_3)^2 + (\lambda_3 + \lambda_1)^2) \end{aligned}$$

and then

$$\begin{split} &|(\lambda_{2}\lambda_{3})V_{1} - (\lambda_{2} + \lambda_{3})V_{2} + V_{3}|^{2} \\ &+ |(\lambda_{3}\lambda_{1})V_{1} - (\lambda_{3} + \lambda_{1})V_{2} + V_{3}|^{2} + |(\lambda_{1}\lambda_{2})V_{1} - (\lambda_{1} + \lambda_{2})V_{2} + V_{3}|^{2} \\ &\succ \gamma_{1}((\lambda_{2} - \lambda_{1})^{2} + (\lambda_{3} - \lambda_{2})^{2} + (\lambda_{3} - \lambda_{1})^{2})|V_{2}|^{2} - \gamma_{2}(\lambda_{1}^{2}\lambda_{2}^{2} + \lambda_{2}^{2}\lambda_{3}^{2} + \lambda_{3}^{2}\lambda_{1}^{2})|V_{1}|^{2} \\ &\succ \gamma_{1}'((\lambda_{1} + \lambda_{2})^{2} + (\lambda_{2} + \lambda_{3})^{2} + (\lambda_{3} + \lambda_{1})^{2})|V_{2}|^{2} - \gamma_{2}(\lambda_{1}^{2}\lambda_{2}^{2} + \lambda_{2}^{2}\lambda_{3}^{2} + \lambda_{3}^{2}\lambda_{1}^{2})|V_{1}|^{2} \\ &\succ (\gamma_{1}' - \gamma_{2}\frac{1}{\delta_{1}}(\delta_{2} + 1))((\lambda_{1} + \lambda_{2})^{2} + (\lambda_{2} + \lambda_{3})^{2} + (\lambda_{3} + \lambda_{1})^{2})|V_{2}|^{2}, \end{split}$$

² For example, we can take $\gamma_1 = \frac{1}{2}$, $\gamma_2 = 1$, using the inequality $|z_1 - z_2|^2 \ge \frac{1}{2}|z_1|^2 - |z_2|^2$.

for some³ constant $\gamma'_1 > 0$. Combining this with (28) we conclude that for any δ_2 and for δ_1 big enough the right-hand side of (24) can be estimated from below by $|V_2|^2$ and, therefore, (24) holds true on $(\Sigma_1^{\delta_1})^c \cap \Sigma_2^{\delta_2}$. Estimate on $(\Sigma_1^{\delta_1})^c \cap (\Sigma_2^{\delta_2})^c$.

Since on $(\Sigma_2^{\delta_2})^c$ we have

$$|V_3|^2 > \delta_2((\lambda_1 + \lambda_2)^2 + (\lambda_2 + \lambda_3)^2 + (\lambda_3 + \lambda_1)^2)|V_2|^2,$$

it follows that

$$|B_1V_1 + B_2V_2 + B_3V_3|^2 \prec |V_3|^2.$$

Then, for $V \in (\Sigma_1^{\delta_1})^c \cap (\Sigma_2^{\delta_2})^c$, for suitable⁴ constants $\gamma_1, \gamma_2, \gamma_3$ (independent of V),

$$\begin{split} &|(\lambda_{2}\lambda_{3})V_{1} - (\lambda_{2} + \lambda_{3})V_{2} + V_{3}|^{2} \\ &+ |(\lambda_{3}\lambda_{1})V_{1} - (\lambda_{3} + \lambda_{1})V_{2} + V_{3}|^{2} + |(\lambda_{1}\lambda_{2})V_{1} - (\lambda_{1} + \lambda_{2})V_{2} + V_{3}|^{2} \\ &\succ \gamma_{3}|V_{3}|^{2} - \gamma_{2}((\lambda_{1} + \lambda_{2})^{2} + (\lambda_{2} + \lambda_{3})^{2} + (\lambda_{3} + \lambda_{1})^{2})|V_{2}|^{2} \\ &- \gamma_{1}\frac{1}{\delta_{1}}(|V_{3}|^{2} + ((\lambda_{1} + \lambda_{2})^{2} + (\lambda_{2} + \lambda_{3})^{2} + (\lambda_{3} + \lambda_{1})^{2})|V_{2}|^{2}) \\ &\succ (\gamma_{3} - \gamma_{1}\frac{1}{\delta_{1}})|V_{3}|^{2} - (\gamma_{2} + \gamma_{1}\frac{1}{\delta_{1}})\frac{1}{\delta_{2}}|V_{3}|^{2}. \end{split}$$

We conclude that for δ_1 and δ_2 big enough,

$$|(\lambda_{2}\lambda_{3})V_{1} - (\lambda_{2} + \lambda_{3})V_{2} + V_{3}|^{2} + |(\lambda_{3}\lambda_{1})V_{1} - (\lambda_{3} + \lambda_{1})V_{2} + V_{3}|^{2} + |(\lambda_{1}\lambda_{2})V_{1} - (\lambda_{1} + \lambda_{2})V_{2} + V_{3}|^{2} \succ |V_{3}|^{2},$$

and, therefore, (24) holds in the area $(\Sigma_1^{\delta_1})^c \cap (\Sigma_2^{\delta_2})^c$.

The next table describes and summarises the proof above:

Area	Estimates in	δ_i
$ \begin{array}{c} \Sigma_1^{\delta_1} \\ (\Sigma_1^{\delta_1})^c \cap \Sigma_2^{\delta_2} \\ (\Sigma_1^{\delta_1})^c \cap (\Sigma_2^{\delta_2})^c \end{array} $	$ V_1 ^2$ $ V_2 ^2$ $ V_3 ^2$	Any δ_1 δ_1 big, any δ_2 δ_1 and δ_2 big

³ Here we use the inequality $|z_1|^2 + |z_2|^2 + |z_3|^2 \ge \frac{1}{4}(|z_1 - z_2|^2 + |z_2 - z_3|^2 + |z_3 - z_1|^2)$. ⁴ Using the inequality $|z_1 - z_2 - z_3|^2 \ge \frac{1}{4}|z_1|^2 - \frac{1}{2}|z_2|^2 - |z_3|^2$.

5.2 The general case m

Inspired by the previous subsection we now deal with the inequality

$$|\mathcal{W}BV| \prec |\mathcal{W}V|$$

for $\mathcal{W} = \mathcal{W}^{(m)}$ and arbitrary $m \ge 2$. This is the topic of the following theorem where the coefficients $\sigma_h^{(m)}(\lambda)$ are defined as in Sect. 3 and $\lambda = (\lambda_1, \lambda_2, \dots, \lambda_m) \in \mathbb{R}^m$ is the vector of the eigenvalues of the matrix $A(t, \xi)$ (or the 0-order normalised roots) satisfying the condition (6).

Theorem 5 Let the entries B_i of the matrix

$$B = \begin{pmatrix} 0 & 0 & 0 & \dots & 0 \\ 0 & 0 & 0 & \dots & 0 \\ \dots & \dots & \dots & \dots & 0 \\ B_1 & B_2 & \dots & \dots & B_m \end{pmatrix}$$

in (13) fulfil the condition

$$|B_j|^2 \prec \sum_{i=1}^m |\sigma_{m-j}^{(m-1)}(\pi_i \lambda)|^2$$
(29)

for $j = 1, \ldots, m$. Then we have

$$|\mathcal{W}BV| \prec |\mathcal{W}V|$$

uniformly over all $V \in \mathbb{C}^m$. More precisely, define

$$\Sigma_{k}^{\delta_{k}} := \{ V \in \mathbb{C}^{m} : |V_{m}|^{2} + \sum_{j=k+1}^{m-1} \sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} \\ \leq \delta_{k} \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2} \},$$

for k = 1, ..., m - 2, and for k = m - 1, define

$$\Sigma_{m-1}^{\delta_{m-1}} := \{ V \in \mathbb{C}^m : |V_m|^2 \le \delta_{m-1} \sum_{i=1}^m |\sigma_1^{(m-1)}(\pi_i \lambda)|^2 |V_{m-1}|^2 \}.$$

Then, there exist suitable $\delta_j > 0, j = 1, ..., m - 1$, such that

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$$|\mathcal{W}BV|^{2} \prec \sum_{i=1}^{m} |\sigma_{m-1}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{1}|^{2},$$
$$|\mathcal{W}V|^{2} \succ \sum_{i=1}^{m} |\sigma_{m-1}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{1}|^{2}$$

on $\Sigma_1^{\delta_1}$,

$$|\mathcal{W}BV|^{2} \prec \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2},$$
$$|\mathcal{W}V|^{2} \succ \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2}$$

on

$$(\Sigma_1^{\delta_1})^{\mathsf{c}} \cap (\Sigma_2^{\delta_2})^{\mathsf{c}} \cap \dots \cap (\Sigma_{k-1}^{\delta_{k-1}})^{\mathsf{c}} \cap \Sigma_k^{\delta_k}$$

for $2 \le k \le m - 1$, and

$$|\mathcal{W}BV|^{2} \prec \sum_{i=1}^{m} |\sigma_{0}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{m}|^{2},$$
$$|\mathcal{W}V|^{2} \succ \sum_{i=1}^{m} |\sigma_{0}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{m}|^{2}$$

on $(\Sigma_1^{\delta_1})^{\mathrm{c}} \cap (\Sigma_2^{\delta_2})^{\mathrm{c}} \cap \cdots \cap (\Sigma_{m-1}^{\delta_{m-1}})^{\mathrm{c}}$.

Note that (29) is a reformulation of the condition (10) on the lower order terms. The proof of Theorem 5 makes use of the following two lemmas.

Lemma 3 For all *i* and *j* with $1 \le i, j \le m$ and k = 1, ..., m - 1, one has

$$\sigma_{m-k}^{(m-1)}(\pi_i\lambda) - \sigma_{m-k}^{(m-1)}(\pi_j\lambda)$$

$$= (-1)^{m-k}(\lambda_j - \lambda_i) \sum_{\substack{i_h \neq i, \ i_h \neq j \\ 1 \le i_1 < i_2 < \dots < i_{m-k-1} \le m}} \lambda_{i_1}\lambda_{i_2} \cdots \lambda_{i_{m-k-1}}.$$
(30)

Proof By definition of $\sigma_{m-k}^{(m-1)}(\pi_i \lambda)$ and $\sigma_{m-k}^{(m-1)}(\pi_j \lambda)$ we have that

$$\sigma_{m-k}^{(m-1)}(\pi_i\lambda) = (-1)^{m-k} \sum_{\substack{1 \le l_1 < l_2 < \cdots < l_{m-k} \le m \\ l_h \neq i}} \lambda_{l_1}\lambda_{l_2}\cdots\lambda_{l_{m-k}}$$

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$$= (-1)^{m-k} \sum_{\substack{1 \le l_1 < l_2 < \cdots < l_{m-k} \le m \\ l_h \ne i, j}} \lambda_{l_1} \lambda_{l_2} \cdots \lambda_{l_{m-k}}$$
$$+ (-1)^{m-k} \lambda_j \sum_{\substack{1 \le l_1 < l_2 < \cdots < l_{m-k-1} \le m \\ l_h \ne i, j}} \lambda_{l_1} \lambda_{l_2} \cdots \lambda_{l_{m-k-1}}$$

and

$$\sigma_{m-k}^{(m-1)}(\pi_{j}\lambda) = (-1)^{m-k} \sum_{\substack{1 \le l_{1} < l_{2} < \dots < l_{m-k} \le m \\ l_{h} \ne j}} \lambda_{l_{1}}\lambda_{l_{2}} \cdots \lambda_{l_{m-k}}}$$
$$= (-1)^{m-k} \sum_{\substack{1 \le l_{1} < l_{2} < \dots < l_{m-k} \le m \\ l_{h} \ne i, j}} \lambda_{l_{1}}\lambda_{l_{2}} \cdots \lambda_{l_{m-k}}}$$
$$+ (-1)^{m-k}\lambda_{i} \sum_{\substack{1 \le l_{1} < l_{2} < \dots < l_{m-k-1} \le m \\ l_{h} \ne i, j}} \lambda_{l_{1}}\lambda_{l_{2}} \cdots \lambda_{l_{m-k-1}}}.$$

This leads immediately to the formula (30).

Lemma 4 For all $k = 1, \ldots, m$, we have

$$\sum_{i=1}^{m} \left| \sum_{j=k+1}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} + \sigma_{m-k}^{(m-1)}(\pi_{i}\lambda) V_{k} \right|^{2} \\ \succ \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2}.$$
(31)

Proof We give a proof by induction on the order m. Setting m = 2 the estimates above makes sense for k = 1. Hence we have to prove that

$$\sum_{i=1}^{2} |\sigma_{0}^{(1)}(\pi_{i}\lambda)V_{2} + \sigma_{1}^{(1)}(\pi_{i}\lambda)V_{1}|^{2} = \sum_{i=1}^{2} |V_{2} + \sigma_{1}^{(1)}(\pi_{i}\lambda)V_{1}|^{2}$$
$$\succ \sum_{i=1}^{2} |\sigma_{1}^{(1)}(\pi_{i}\lambda)|^{2} |V_{1}|^{2}.$$

This is clear since by the condition (6) we have that

$$\sum_{i=1}^{2} |V_2 + \sigma_1^{(1)}(\pi_i \lambda) V_1|^2 \succ |\sigma_1^{(1)}(\pi_1 \lambda) - \sigma_1^{(1)}(\pi_2 \lambda)|^2 |V_1|^2 = (\lambda_2 - \lambda_1)^2 |V_1|^2$$
$$\succ (\lambda_1^2 + \lambda_2^2) |V_1|^2 = \sum_{i=1}^{2} |\sigma_1^{(1)}(\pi_i \lambda)|^2 |V_1|^2.$$

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Assume now that (31) holds for m - 1. Estimating the left-hand side of (31) with the differences between two arbitrary summands⁵ we can write

$$\begin{split} &\sum_{i=1}^{m} \left| \sum_{j=k+1}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} + \sigma_{m-k}^{(m-1)}(\pi_{i}\lambda) V_{k} \right|^{2} \\ & \succ \sum_{1 \leq l_{1} \neq l_{2} \leq m} \left| \sum_{j=k+1}^{m} (\sigma_{m-j}^{(m-1)}(\pi_{l_{1}}\lambda) - \sigma_{m-j}^{(m-1)}(\pi_{l_{2}}\lambda)) V_{j} \right| \\ & + (\sigma_{m-k}^{(m-1)}(\pi_{l_{1}}\lambda) - \sigma_{m-k}^{(m-1)}(\pi_{l_{2}}\lambda)) V_{k} \right|^{2} \\ & = \sum_{1 \leq l_{1} \neq l_{2} \leq m} \left| \sum_{j=k+1}^{m-1} (\sigma_{m-j}^{(m-1)}(\pi_{l_{1}}\lambda) - \sigma_{m-j}^{(m-1)}(\pi_{l_{2}}\lambda)) V_{j} \right| \\ & + (\sigma_{m-k}^{(m-1)}(\pi_{l_{1}}\lambda) - \sigma_{m-k}^{(m-1)}(\pi_{l_{2}}\lambda)) V_{k} \right|^{2}, \end{split}$$

where the terms with j = m cancel. By applying Lemma 3 and the condition (6) we obtain the following bound from below:

$$\begin{split} \sum_{i=1}^{m} \left| \sum_{j=k+1}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} + \sigma_{m-k}^{(m-1)}(\pi_{i}\lambda) V_{k} \right|^{2} &\succ \sum_{1 \leq l_{1} \neq l_{2} \leq m} \left| -(\lambda_{l_{2}} - \lambda_{l_{1}}) V_{m-1} \right. \\ &+ \sum_{j=k+1}^{m-2} (-1)^{m-j} (\lambda_{l_{2}} - \lambda_{l_{1}}) \sum_{\substack{i_{h} \neq l_{1}, i_{h} \neq l_{2} \\ 1 \leq i_{1} < i_{2} < \cdots < i_{m-j-1} \leq m}} (\lambda_{i_{1}}\lambda_{i_{2}} \cdots \lambda_{i_{m-j-1}}) V_{j} \\ &+ (-1)^{m-k} (\lambda_{l_{2}} - \lambda_{l_{1}}) \sum_{\substack{i_{h} \neq l_{1}, i_{h} \neq l_{2} \\ 1 \leq i_{1} < i_{2} < \cdots < i_{m-k-1} \leq m}} (\lambda_{i_{1}}\lambda_{i_{2}} \cdots \lambda_{i_{m-k-1}}) V_{k} \right|^{2} \\ &\succ \sum_{1 \leq l_{1} \neq l_{2} \leq m} (\lambda_{l_{1}}^{2} + \lambda_{l_{2}}^{2}) \left| -V_{m-1} + \sum_{j=k+1}^{m-2} (-1)^{m-j} \sum_{\substack{i_{h} \neq l_{1}, i_{h} \neq l_{2} \\ 1 \leq i_{1} < i_{2} < \cdots < i_{m-j-1} \leq m}} (\lambda_{i_{1}}\lambda_{i_{2}} \cdots \lambda_{i_{m-j-1}}) V_{j} + (-1)^{m-k} \sum_{\substack{i_{h} \neq l_{1}, i_{h} \neq l_{2} \\ 1 \leq i_{1} < i_{2} < \cdots < i_{m-k-1} \leq m}} (\lambda_{i_{1}}\lambda_{i_{2}} \cdots \lambda_{i_{m-k-1}}) V_{k} \right|^{2}. \end{split}$$

Noting that

$$(-1)^{m-j} \sum_{\substack{i_h \neq l_1, i_h \neq l_2\\ 1 \leq i_1 < i_2 < \dots < i_{m-j-1} \leq m}} (\lambda_{i_1} \lambda_{i_2} \cdots \lambda_{i_{m-j-1}}) = -\sigma_{m-1-j}^{(m-2)} (\pi_{l_2}(\pi_{l_1} \lambda))$$

 $\overline{5 \text{ Using the inequality } \sum_{i=1}^{m} |z_i|^2} \succ \sum_{1 \le h_1 \ne h_2 \le m} |z_{h_1} - z_{h_2}|^2.$

for $j = k, \ldots, m - 2$, we write the estimate above as

$$\sum_{i=1}^{m} \left| \sum_{j=k+1}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} + \sigma_{m-k}^{(m-1)}(\pi_{i}\lambda) V_{k} \right|^{2} \succ \sum_{1 \le l_{1} \ne l_{2} \le m} (\lambda_{l_{1}}^{2} + \lambda_{l_{2}}^{2}) \\ \left| -V_{m-1} - \sum_{j=k+1}^{m-2} \sigma_{m-1-j}^{(m-2)}(\pi_{l_{2}}(\pi_{l_{1}}\lambda)) V_{j} - \sigma_{m-1-k}^{(m-2)}(\pi_{l_{2}}(\pi_{l_{1}}\lambda)) V_{k} \right|^{2},$$

where the right hand-side can be written as

$$\sum_{1 \le l_1 \ne l_2 \le m} (\lambda_{l_1}^2 + \lambda_{l_2}^2) \\ \left| V_{m-1} + \sum_{j=k+1}^{m-2} \sigma_{m-1-j}^{(m-2)} (\pi_{l_2}(\pi_{l_1}\lambda)) V_j + \sigma_{m-1-k}^{(m-2)} (\pi_{l_2}(\pi_{l_1}\lambda)) V_k \right|^2 \\ = \sum_{l_1} \lambda_{l_1}^2 \sum_{1 \le l_1 \ne l_2 \le m} \left| V_{m-1} + \sum_{j=k+1}^{m-2} \sigma_{m-1-j}^{(m-2)} (\pi_{l_2}(\pi_{l_1}\lambda)) V_j \right|^2 \\ + \sigma_{m-1-k}^{(m-2)} (\pi_{l_2}(\pi_{l_1}\lambda)) V_k \right|^2 \\ + \sum_{l_2} \lambda_{l_2}^2 \sum_{1 \le l_1 \ne l_2 \le m} \left| V_{m-1} + \sum_{j=k+1}^{m-2} \sigma_{m-1-j}^{(m-2)} (\pi_{l_2}(\pi_{l_1}\lambda)) V_j \right|^2 \\ + \sigma_{m-1-k}^{(m-2)} (\pi_{l_2}(\pi_{l_1}\lambda)) V_k \right|^2.$$
(32)

By now applying the inductive hypothesis to the last two summands in (32) we obtain

$$\begin{split} &\sum_{i=1}^{m} \left| \sum_{j=k+1}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} + \sigma_{m-k}^{(m-1)}(\pi_{i}\lambda) V_{k} \right|^{2} \\ & \succ \sum_{l_{1}} \lambda_{l_{1}}^{2} \sum_{1 \leq l_{1} \neq l_{2} \leq m} |\sigma_{m-1-k}^{(m-2)}(\pi_{l_{2}}(\pi_{l_{1}}\lambda))|^{2} |V_{k}|^{2} \\ & + \sum_{l_{2}} \lambda_{l_{2}}^{2} \sum_{1 \leq l_{1} \neq l_{2} \leq m} |\sigma_{m-1-k}^{(m-2)}(\pi_{l_{1}}(\pi_{l_{2}}\lambda))|^{2} |V_{k}|^{2} \\ & \succ \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2}, \end{split}$$

which completes the proof.

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Proof of Theorem 5 By definition of the matrices W and B we have that $|WBV|^2 \prec$ $|WV|^2$ is equivalent to

$$\left|\sum_{j=1}^{m} B_{j} V_{j}\right|^{2} \prec \sum_{i=1}^{m} \left|\sum_{j=1}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j}\right|^{2}.$$
(33)

Making use of the conditions (29) we have that the following estimate is valid on the area $\Sigma_1^{\overline{\delta}_1}$:

$$\begin{split} \left| \sum_{j=1}^{m} B_{j} V_{j} \right|^{2} &\prec \sum_{j=1}^{m} |B_{j}|^{2} |V_{j}|^{2} \prec \sum_{j=1}^{m} \sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} \\ &\prec |V_{m}|^{2} + \sum_{j=2}^{m-1} \sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} + \sum_{i=1}^{m} |\sigma_{m-1}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{1}|^{2} \\ &\prec (1+\delta_{1}) \sum_{i=1}^{m} |\sigma_{m-1}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{1}|^{2} \prec \sum_{i=1}^{m} |\sigma_{m-1}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{1}|^{2}. \end{split}$$

Setting k = 1 in Lemma 4 we obtain the bound from below

$$\sum_{i=1}^{m} \left| \sum_{j=1}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} \right|^{2} = \sum_{i=1}^{m} \left| \sum_{j=2}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} + \sigma_{m-1}^{(m-1)}(\pi_{i}\lambda) V_{l} \right|^{2}$$
$$\succ \sum_{i=1}^{m} |\sigma_{m-1}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{l}|^{2}.$$

This proves the inequality (33) on $\Sigma_1^{\delta_1}$ for any $\delta_1 > 0$. Let us now assume that $V \in (\Sigma_1^{\delta_1})^c \cap (\Sigma_2^{\delta_2})^c \cap \dots \cap (\Sigma_{k-1}^{\delta_{k-1}})^c \cap \Sigma_k^{\delta_k}$ for $2 \le k \le m-1$. By definition of the regions $\Sigma_h^{\delta_h}$ and taking $\delta_h \ge 1$ for $1 \le h \le k-1$ we have that

$$\sum_{i=1}^{m} |\sigma_{m-(k-1)}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k-1}|^{2}$$

$$< \frac{1}{\delta_{k-1}} \left(|V_{m}|^{2} + \sum_{j=k+1}^{m-1} \sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} + \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2} \right)$$

$$\leq \frac{1}{\delta_{k-1}} (1 + \delta_{k}) \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2},$$

$$\begin{split} &\sum_{i=1}^{m} |\sigma_{m-(k-2)}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k-2}|^{2} \\ &< \frac{1}{\delta_{k-2}} \bigg(|V_{m}|^{2} + \sum_{j=k+1}^{m-1} \sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} \\ &+ \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2} + \sum_{i=1}^{m} |\sigma_{m-(k-1)}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k-1}|^{2} \bigg) \\ &\leq \frac{1}{\delta_{k-2}} \left(1 + \delta_{k} + \frac{1}{\delta_{k-1}} (1 + \delta_{k}) \right) \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2}, \\ &\leq (1 + \delta_{k}) \left(\frac{1}{\delta_{k-1}} + \frac{1}{\delta_{k-2}} \right) \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2}. \end{split}$$

By iteration one can easily prove the following bound

$$\sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_i \lambda)|^2 |V_j|^2 \le (1+\delta_k) \sum_{h=1}^{k-1} \frac{1}{\delta_h} \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_i \lambda)|^2 |V_k|^2, \quad (34)$$

valid on the region $(\Sigma_1^{\delta_1})^c \cap (\Sigma_2^{\delta_2})^c \cap \cdots \cap (\Sigma_{k-1}^{\delta_{k-1}})^c \cap \Sigma_k^{\delta_k}$ for all j with $1 \le j \le k-1$. It follows that

$$\begin{split} & \left| \sum_{j=1}^{m} B_{j} V_{j} \right|^{2} \prec \sum_{j=1}^{m} |B_{j}|^{2} |V_{j}|^{2} \prec \sum_{j=1}^{m} \sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} \\ & \prec \sum_{j=k+1}^{m} \sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} + \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2} \\ & + \sum_{j=1}^{k-1} \sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} \prec \sum_{i=1}^{m} |\sigma_{m-k}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{k}|^{2}. \end{split}$$

We now pass to estimate the right-hand side of (33) making use of Lemma 4 and of the bound (34). We obtain

$$\begin{split} &\sum_{i=1}^{m} \left| \sum_{j=1}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} \right|^{2} \\ & \succ \sum_{i=1}^{m} \gamma_{1} \left| \sum_{j=k+1}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} + \sigma_{m-k}^{(m-1)}(\pi_{i}\lambda) V_{k} \right|^{2} \\ & - \gamma_{2} \sum_{i=1}^{m} \sum_{j=1}^{k-1} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} \end{split}$$

$$\succ \gamma_1 \sum_{i=1}^m |\sigma_{m-k}^{(m-1)}(\pi_i \lambda)|^2 |V_k|^2 - \gamma_2 (1+\delta_k) \sum_{h=1}^{k-1} \frac{1}{\delta_h} \sum_{i=1}^m |\sigma_{m-k}^{(m-1)}(\pi_i \lambda)|^2 |V_k|^2$$

$$= \left(\gamma_1 - \gamma_2 (1+\delta_k) \sum_{h=1}^{k-1} \frac{1}{\delta_h}\right) \sum_{i=1}^m |\sigma_{m-k}^{(m-1)}(\pi_i \lambda)|^2 |V_k|^2.$$

Therefore, the estimate (33) holds in the region $(\Sigma_1^{\delta_1})^c \cap (\Sigma_2^{\delta_2})^c \cap \cdots \cap (\Sigma_{k-1}^{\delta_{k-1}})^c \cap \Sigma_k^{\delta_k}$ for any $\delta_k > 0$ choosing $\delta_1, \delta_2, \dots, \delta_{k-1}$ big enough. We conclude the proof by assuming $V \in (\Sigma_1^{\delta_1})^c \cap (\Sigma_2^{\delta_2})^c \cap \dots \cap (\Sigma_{m-1}^{\delta_{m-1}})^c$. Since

$$|V_m|^2 + \sum_{j=h+1}^{m-1} \sum_{i=1}^m |\sigma_{m-j}^{(m-1)}(\pi_i \lambda)|^2 |V_j|^2 > \delta_h \sum_{i=1}^m |\sigma_{m-h}^{(m-1)}(\pi_i \lambda)|^2 |V_h|^2,$$

for $1 \le h \le m - 1$, arguing as above and taking $\delta_h \ge 1$ we obtain the estimate

$$\sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_i \lambda)|^2 |V_j|^2 \le \sum_{h=1}^{m-1} \frac{1}{\delta_h} |V_m|^2$$
(35)

valid on the region $(\Sigma_1^{\delta_1})^c \cap (\Sigma_2^{\delta_2})^c \cap \cdots \cap (\Sigma_{m-1}^{\delta_{m-1}})^c$ for all j with $1 \leq j \leq m-1$. Hence,

$$\left|\sum_{j=1}^{m} B_{j} V_{j}\right|^{2} \prec \sum_{j=1}^{m} |B_{j}|^{2} |V_{j}|^{2} \prec \sum_{j=1}^{m} \sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} \prec |V_{m}|^{2}$$

and

$$\begin{split} &\sum_{i=1}^{m} \left| \sum_{j=1}^{m} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} \right|^{2} \succ \gamma_{1} |V_{m}|^{2} - \gamma_{2} \sum_{i=1}^{m} \left| \sum_{j=1}^{m-1} \sigma_{m-j}^{(m-1)}(\pi_{i}\lambda) V_{j} \right|^{2} \\ & \succ \gamma_{1} |V_{m}|^{2} - \gamma_{2} \sum_{j=1}^{m-1} \sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_{i}\lambda)|^{2} |V_{j}|^{2} \\ & \succ \left(\gamma_{1} - \gamma_{2}(m-1) \sum_{h=1}^{m-1} \frac{1}{\delta_{h}} \right) |V_{m}|^{2}. \end{split}$$

This means that the inequality (33) holds on $(\Sigma_1^{\delta_1})^c \cap (\Sigma_2^{\delta_2})^c \cap \cdots \cap (\Sigma_{m-1}^{\delta_{m-1}})^c$ for sufficiently large values of $\delta_1, \delta_2, \ldots, \delta_{m-1}$.

6 Well-posedness results

We are now ready to prove the well-posedness results given in Theorem 2. For the advantage of the reader and the sake of clarity we reformulate Theorem 2 as the following Theorem 6 where we make use of the language and notations introduced in Theorem 5.

Theorem 6 Assume $A_j \in C([0, T])$ for all j. If the coefficients satisfy $A_{(j)} \in C^{\infty}([0, T])$, the characteristic roots are real and satisfy (6) and the entries of the matrix B of the lower order terms in (13) fulfil the conditions (29) for ξ away from 0, then the Cauchy problem (1) is well-posed in any Gevrey space. More precisely,

(i) if $A_{(j)} \in C^k([0, T])$ for some $k \ge 2$ and $g_j \in G^s(\mathbb{R}^n)$ for j = 1, ..., m, then there exists a unique solution $u \in C^m([0, T]; G^s(\mathbb{R}^n))$ for

$$1 \le s < 1 + \frac{k}{2(m-1)};$$

(ii) if $A_{(j)} \in C^k([0, T])$ for some $k \ge 2$ and $g_j \in \mathcal{E}'_{(s)}(\mathbb{R}^n)$ for j = 1, ..., m, then there exists a unique solution $u \in C^m([0, T]; \mathcal{D}'_{(s)}(\mathbb{R}^n))$ for

$$1 \le s \le 1 + \frac{k}{2(m-1)}.$$

Proof As usual, the well-posedness in the case of s = 1 follows from the result of Bony and Shapira, so we may assume s > 1. By the finite propagation speed for hyperbolic equations it is not restrictive to take compactly supported initial data and, therefore, to have that the solution u is compactly supported in x.

Combining the energy estimate (15) with the estimates of the first, second and third term in Sect. 4 we obtain the estimate

$$\partial_t E_{\varepsilon}(t,\xi) \le (K_{\varepsilon}(t,\xi) + C_2 \varepsilon \langle \xi \rangle + C_3) E_{\varepsilon}(t,\xi), \tag{36}$$

where $K_{\varepsilon}(t,\xi)$ defined in (19) in Sect. 4.1 has the property (20), i.e.

$$\int_{0}^{T} K_{\varepsilon}(t,\xi) \, dt \le C_1 \varepsilon^{-2(m-1)/k},$$

and C_1 , C_2 , C_3 are positive constants. Thus, (36) holds for $t \in [0, T]$ and $|\xi| \ge R$, with the estimate for the third term provided by Theorem 5. A straightforward application of Gronwall's lemma leads to

$$E_{\varepsilon}(t,\xi) \leq E_{\varepsilon}(0,\xi) e^{C_{1}\varepsilon^{-2(m-1)/k} + C_{2}T_{\varepsilon}(\xi) + C_{3}T}$$
$$\leq E_{\varepsilon}(0,\xi) C_{T} e^{C_{T}(\varepsilon^{-2(m-1)/k} + \varepsilon(\xi))}.$$

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Setting $\varepsilon^{-2(m-1)/k} = \varepsilon \langle \xi \rangle$ we get

$$E_{\varepsilon}(t,\xi) \le E_{\varepsilon}(0,\xi)C_T \mathrm{e}^{C_T \langle \xi \rangle \, \overline{\sigma}}$$

where $\sigma = 1 + k/[2(m-1)]$. Finally, making use of the inequality (17) we arrive at

$$C_m^{-1}\varepsilon^{2(m-1)}|V(t,\xi)|^2 \le E_{\varepsilon}(t,\xi) \le E_{\varepsilon}(0,\xi)C_T e^{C_T\langle\xi\rangle\frac{1}{\sigma}}$$
$$\le C_m|V(0,\xi)|^2 C_T e^{C_T\langle\xi\rangle\frac{1}{\sigma}},$$

which implies

$$|V(t,\xi)| \le C\langle\xi\rangle^{\frac{k}{2\sigma}} e^{C\langle\xi\rangle^{\frac{1}{\sigma}}} |V(0,\xi)|,$$
(37)

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for some new constant C > 0, for $t \in [0, T]$ and $|\xi| \ge R$.

(i) Recall that $V(t,\xi) = \mathcal{F}_{x\to\xi}U(t,x)$, where U is the u_j 's column vector. If the initial data g_j belong to $G_0^s(\mathbb{R}^n)$ from the Fourier transform characterisation of

Gevrey functions ([11, Proposition 2.2]) we have that $|V(0,\xi)| \le c e^{-\delta(\xi)^{\frac{1}{s}}}$ for some constants c > 0 and $\delta > 0$. Hence,

$$|V(t,\xi)| \le C\langle \xi \rangle^{\frac{k}{2\sigma}} \mathrm{e}^{C\langle \xi \rangle^{\frac{1}{\sigma}}} c \, \mathrm{e}^{-\delta\langle \xi \rangle^{\frac{1}{s}}}$$

for all $t \in [0, T]$ and $\xi \in \mathbb{R}^n$. Let $s < \sigma$. Then $V(t, \xi)$ defines a tempered distribution in $\xi'(\mathbb{R}^n)$ such that

$$\begin{aligned} |V(t,\xi)| &\leq Cc\langle\xi\rangle^{\frac{k}{2\sigma}} \mathrm{e}^{C\langle\xi\rangle^{\frac{1}{\sigma}}} \mathrm{e}^{-\frac{\delta}{2}\langle\xi\rangle^{\frac{1}{s}}} \mathrm{e}^{-\frac{\delta}{2}\langle\xi\rangle^{\frac{1}{s}}} \\ &\leq Cc\langle\xi\rangle^{\frac{k}{2\sigma}} \mathrm{e}^{\langle\xi\rangle^{\frac{1}{\sigma}}(C-\frac{\delta}{2}\langle\xi\rangle^{\frac{1}{s}-\frac{1}{\sigma}})} \mathrm{e}^{-\frac{\delta}{2}\langle\xi\rangle^{\frac{1}{s}}}. \end{aligned}$$

It follows that

$$|V(t,\xi)| \le c' \mathrm{e}^{-\frac{\delta}{2}\langle\xi\rangle^{\frac{1}{s}}},$$
(38)

for some $c', \delta > 0$ and for $|\xi|$ large enough. This is sufficient to prove that U(t, x) belongs to the Gevrey class $G^{s}(\mathbb{R}^{n})$ for all $t \in [0, T]$ and that the Cauchy problem (1) has a unique solution $u \in C^{m}([0, T]; G^{s}(\mathbb{R}^{n}))$ for $s < \sigma$ under the assumptions of case (i).

(ii) If the initial data g_j are Gevrey Beurling ultradistributions in E'_(s)(ℝⁿ), from the Fourier transform characterisation of ultradistributions ([11, Proposition 2.13]) we have that there exist δ > 0 and c > 0 such that |V(0, ξ)| ≤ c e^{δ(ξ)¹/s} for all ξ ∈ ℝⁿ. Hence, taking s < σ, we obtain the estimate

$$|V(t,\xi)| \le Cc\langle\xi\rangle^{\frac{k}{2\sigma}} \mathrm{e}^{C\langle\xi\rangle^{\frac{1}{\sigma}}} \mathrm{e}^{+\delta\langle\xi\rangle^{\frac{1}{s}}} \le c' \mathrm{e}^{+\delta'\langle\xi\rangle^{\frac{1}{s}}}$$

for some $c', \delta' > 0$. This proves that the Cauchy problem (1) has a unique solution $u \in C^m([0, T]; \mathcal{D}'_{(s)}(\mathbb{R}^n))$ for $s \leq \sigma$ under the assumptions of case (ii).

We pass to consider the case of analytic coefficients. We prove C^{∞} and distributional well-posedness of the Cauchy problem (1) providing an extension of Theorem 1 in [13] to any space dimension. Our proof makes use of the following lemma on analytic functions, a parameter-dependent version of the statement (61)–(62) in [13].

Lemma 5 Let $f(t, \xi)$ be an analytic function in $t \in [0, T]$, continuous and homogeneous of order 0 in $\xi \in \mathbb{R}^n$. Then,

(i) for all ξ there exists a partition $(\tau_{h(\xi)})$ of the interval [0, T] such that

$$0 = \tau_0 < \tau_1 < \cdots < \tau_{h(\xi)} < \cdots < \tau_{N(\xi)} = T$$

with $\sup_{\xi \neq 0} N(\xi) < \infty$, such that $f(t, \xi) \neq 0$ in each open interval $(\tau_{h(\xi)}, \tau_{h+1(\xi)});$

(ii) there exists C > 0 such that

$$|\partial_t f(t,\xi)| \le C \left(\frac{1}{t - \tau_{h(\xi)}} + \frac{1}{\tau_{h+1(\xi)} - t}\right) |f(t,\xi)|$$

for all $t \in (\tau_{h(\xi)}, \tau_{h+1(\xi)}), \xi \in \mathbb{R}^n$ with $\xi \neq 0$ and $0 \leq h(\xi) \leq N(\xi)$.

Proof Since the function f is homogeneous of order 0 in ξ we can assume $|\xi| = 1$. Excluding the trivial case $f \equiv 0$ we have that $f(t, \xi)$ has a finite number of zeroes in [0, T] and hence we can find a partition $(\tau_{h(\xi)})$ as in (i) such that $f(t, \xi) \neq 0$ in each interval $(\tau_{h(\xi)}, \tau_{h+1(\xi)})$, taking $\tau_{h(\xi)}, 1 \leq h(\xi) \leq N(\xi) - 1$, to be the zeros of $f(\cdot, \xi)$.

Note that the function $N(\xi)$ is locally bounded and, therefore, by homogeneity $\sup_{\xi \neq 0} N(\xi) = \sup_{|\xi|=1} N(\xi) < \infty$. Indeed, if $\sup_{|\xi|=1} N(\xi) = +\infty$ we can find a sequence of points $(\xi_l)_l$ with $|\xi_l| = 1$ and some ξ' with $|\xi'| = 1$ such that $\xi_l \to \xi'$ and $N(\xi_l) \to +\infty$ as $l \to \infty$. It follows that $f(t, \xi')$ must have infinite zeros in t in contradiction with the hypothesis of analyticity on [0, T].

We now work on the interval $(0, \tau_1)$. By the analiticity in t we can write

$$f(t,\xi) = t^{\nu_0(\xi)} (\tau_1 - t)^{\nu_1(\xi)} g(t,\xi)$$

where $g(t, \xi)$ is an analytic function in *t* never vanishing on $[0, \tau_1]$ homogeneous of degree 0 in ξ . Note that the functions v_0 and v_1 are positive and have local maxima at all points (perturbations in ξ in a sufficiently small neighborhood can not increase the multiplicity). Arguing as in [13, p. 566] we write $t|\partial_t f(t, \xi)|$ as

$$\left| f(t) \left(v_0(\xi) - \frac{v_1(\xi)t}{(\tau_1 - t)} + \frac{t \partial_t g(t, \xi)}{g(t, \xi)} \right) \right|.$$

Let us fix ξ_0 with $|\xi_0| = 1$. Taking *t* in $[0, \tau_1/2]$ and ξ in a sufficiently small neighborhood of ξ_0 we have that $\nu_0(\xi) \le c_1$, $\nu_1(\xi)t/(\tau_1 - t) \le c_2$, $|g(t, \xi)| \ge c_0 > 0$ and $t\partial_t g(t, \xi)/g(t, \xi) \le c_3$. Hence,

$$t|\partial_t f(t,\xi)| \le C|f(t,\xi)|$$

on $[0, \tau_1/2]$ for ξ in a neighborhood of ξ_0 . Similarly, one proves that

$$(\tau_1 - t)|\partial_t f(t,\xi)| \le C|f(t,\xi)|$$

on $[\tau_1/2, \tau_1]$ for ξ in a neighborhood of ξ_0 . The homogeneity in ξ combined with a standard compactness argument allows us to extend the inequality

$$|\partial_t f(t,\xi)| \le C \left(\frac{1}{t} + \frac{1}{\tau_1 - t}\right) |f(t,\xi)|$$

to $\mathbb{R}^n \setminus \{0\}$ for $t \in (0, \tau_1)$. Analogously one obtains that

$$|\partial_t f(t,\xi)| \le C \left(\frac{1}{t - \tau_{h(\xi)}} + \frac{1}{\tau_{h+1(\xi)} - t} \right) |f(t,\xi)|$$

when $t \in (\tau_{h(\xi)}, \tau_{h+1(\xi)})$ and $\xi \neq 0$.

In the case of analytic coefficients, Theorem 3 follows from the following Theorem 7.

Theorem 7 If $A_j \in C([0, T])$ for all j, and the coefficients $A_{(j)}$ are analytic on [0, T], the characteristic roots are real and satisfy (6), and the entries of the matrix B of the lower order terms in (13) fulfil the conditions (29) for ξ away from 0, then the Cauchy problem (1) is C^{∞} and distributionally well-posed.

Proof By the finite propagation speed for hyperbolic equations it is not restrictive to assume that the initial data g_j are compactly supported. If the coefficients a_j are analytic in t on [0, T] then by construction the entries of the quasi-symmetriser $Q_{\varepsilon}^{(m)}$ are analytic as well. In particular, by Proposition 1

$$q_{\varepsilon,ij}(t,\xi) = q_{0,ij}(t,\xi) + \varepsilon^2 q_{1,ij}(t,\xi) + \dots + \varepsilon^{2(m-1)} q_{m-1,ij}(t,\xi).$$

We use the partition of the interval [0, T] in Lemma 5 (applied to any $q_{\varepsilon,ij}(t, \xi)$ or more precisely to any $\tilde{q}_{\varepsilon,ij}(t, \xi) = q_{\varepsilon,ij}(\lambda(t, \xi)/|\xi|)$, homogeneous function of order 0 in ξ having the same zeros in t of $q_{\varepsilon,ij}(t, \xi)$). Note that this partition can be chosen independent of ε . Considering the first interval $[0, \tau_1]$ ($\tau_1 = \tau_1(\xi)$) we define

$$E_{\varepsilon}(t,\xi) = \begin{cases} |V(t,\xi)|^2 & \text{for } t \in [0,\varepsilon] \cup [\tau_1 - \varepsilon, \tau_1], \\ (Q_{\varepsilon}(t,\xi)V(t,\xi), V(t,\xi)) & \text{for } t \in [\varepsilon, \tau_1 - \varepsilon]. \end{cases}$$

as in [13, p. 567]. Hence

$$\begin{aligned} \partial_t E_{\varepsilon}(t,\xi) &\leq |\partial_t E_{\varepsilon}(t,\xi)| \leq |((A_1 - A_1^*)V, V)| + |((B - B^*)V, V)| \\ &\leq \left(2 \sup_{t \in [0,T]} \|A_1(t,\xi)\| + 2 \sup_{t \in [0,T]} \|B(t,\xi)\|\right) E_{\varepsilon}(t,\xi) \end{aligned}$$

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on $[0, \varepsilon] \cup [\tau_1 - \varepsilon, \tau_1]$. It follows by the Gronwall inequality that there exists a constant $\alpha > 0$ such that

$$E_{\varepsilon}(t,\xi) \leq \begin{cases} e^{2\alpha\varepsilon\langle\xi\rangle} E_{\varepsilon}(0,\xi) & \text{for } t \in [0,\varepsilon], \\ e^{2\alpha\varepsilon\langle\xi\rangle} E_{\varepsilon}(\tau_1 - \varepsilon,\xi) & \text{for } t \in [\tau_1 - \varepsilon,\tau_1]. \end{cases}$$
(39)

On the interval $[\varepsilon, \tau_1 - \varepsilon]$ we proceed as in the proof for the Gevrey well-posedness under the conditions (29) on the lower order terms for $|\xi| \ge R$. We have

$$\partial_t E_{\varepsilon}(t,\xi) \le \left(\frac{|(\partial_t Q_{\varepsilon}V,V)|}{(Q_{\varepsilon}(t,\xi)V(t,\xi),V(t,\xi))} + C_2\varepsilon\langle\xi\rangle + C_3\right)E_{\varepsilon}(t,\xi).$$
(40)

Since the family $Q_{\varepsilon}(\lambda)$ is nearly diagonal when the roots λ_l satisfy the condition (6) we have that $Q_{\varepsilon} \ge c_0 \text{diag } Q_{\varepsilon}$, i.e.,

$$(Q_{\varepsilon}(t,\xi)V,V) \ge c_0 \sum_{h=1}^m q_{\varepsilon,hh}(t,\xi) |V_h|^2.$$

This fact combined with the inequality

$$\begin{aligned} |q_{\varepsilon,ij}||V_i||V_j| &= (Q_{\varepsilon}e_i, e_j)|V_i||V_j| \le \sqrt{(Q_{\varepsilon}e_i, e_i)(Q_{\varepsilon}e_j, e_j)}|V_i||V_j| \\ &\le \sqrt{q_{\varepsilon,ii}q_{\varepsilon,jj}}|V_i||V_j| \le \sum_{h=1}^m q_{\varepsilon,hh}|V_h|^2 \end{aligned}$$

and Lemma 5 yields

$$\int_{\varepsilon}^{\tau_1-\varepsilon} \frac{|(\partial_t Q_{\varepsilon} V, V)|}{(Q_{\varepsilon}(t,\xi)V(t,\xi), V(t,\xi))} dt \le c_0^{-1} \int_{\varepsilon}^{\tau_1-\varepsilon} \sum_{i,j=1}^m \frac{|\partial_t q_{\varepsilon,ij}(t,\xi)|}{|q_{\varepsilon,ij}(t,\xi)|} dt$$
$$\le C_1 \int_{\varepsilon}^{\tau_1-\varepsilon} \left(\frac{1}{t} + \frac{1}{\tau_1-t}\right) dt = 2C_1 \log \frac{\tau_1-\varepsilon}{\varepsilon} \le 2C_1 \log \frac{T}{\varepsilon},$$

for some constant C_1 independent of t and $\xi \neq 0$. Going back to estimate (40) by Gronwall's lemma we obtain

$$E_{\varepsilon}(t,\xi) \le C_T E_{\varepsilon}(\varepsilon,\xi) e^{C_T \log(1/\varepsilon) + C_T \varepsilon(\xi)},\tag{41}$$

for $[\varepsilon, \tau_1 - \varepsilon]$ and $|\xi| \ge R$. Finally, putting together (39) with (41) we conclude that there exists a constant c > 0 such that

$$E_{\varepsilon}(t,\xi) \leq c E_{\varepsilon}(0,\xi) e^{c(\log(1/\varepsilon) + \varepsilon \langle \xi \rangle)}$$

for all $t \in [0, \tau_1]$ and $|\xi| \ge R$. Hence by applying (17) we have

$$|V(t,\xi)| \le c\varepsilon^{-(m-1)} e^{C_T (\log(1/\varepsilon) + \varepsilon\langle\xi\rangle)} |V(0,\xi)|$$

on $[0, \tau_1]$. An iteration of the same technique on the other subintervals of [0, T] leads to

$$|V(t,\xi)| \le c\varepsilon^{-N(\xi)(m-1)} e^{N(\xi)C_T(\log(1/\varepsilon) + \varepsilon\langle\xi\rangle)} |V(0,\xi)|$$

on [0, T] for $|\xi| \ge R$. Now, setting $\varepsilon = \langle \xi \rangle^{-1}$ we get

$$|V(t,\xi)| \le c \langle \xi \rangle^{N(\xi)(m-1)} \mathrm{e}^{N(\xi)C_T} \langle \xi \rangle^{N(\xi)C_T}$$

Remembering that from Lemma 5 the function $N(\xi)$ is bounded in ξ we conclude that there exist some $\kappa \in \mathbb{N}$ and C > 0 such that

$$|V(t,\xi)| \le C\langle\xi\rangle^{\kappa} |V(0,\xi)| \tag{42}$$

on [0, T] for all $|\xi| \ge R$. It is clear that the estimate (42) implies C^{∞} and distributional well-posedness of the Cauchy problem (1).

Finally, given the energy estimates established above, the proof of Theorem 4 is simple:

Proof of Theorem 4 We observe that the estimates (38) and (42) imply that $V(t,\xi)$ is bounded in ξ if the lower order terms $A(\cdot,\xi)$ are bounded on [0, T]. Coming back to the solution u of (1) and the definition of V we get that the solution u(t, x) is in the class $C^{m-1}([0, T])$ with respect to t. Finally, from the equality $D_t^m u = -\sum_{j=0}^{m-1} A_{m-j}(t, D_x) D_t^j u$ we see that the right hand side is bounded in t, implying that u(t, x) is in $W^{\infty,m}([0, T])$ with respect to t.

We conclude the paper with the following remarks on how the results change if we assume less than the Levi conditions (7). We thank T. Kinoshita for drawing our attention to this question.

We begin by noting that the matrix B of the lower order terms in (13) can be written as

$$B(t,\xi) = \sum_{l=0}^{m-1} B_{-l}(t,\xi),$$

with

$$B_{-l} = \begin{pmatrix} 0 & 0 & 0 & \dots & 0 \\ 0 & 0 & 0 & \dots & 0 \\ \dots & \dots & \dots & \dots & 0 \\ B_{-l,1} & B_{-l,2} & \dots & \dots & B_{-l,m} \end{pmatrix}$$

and

$$B_{-l,j}(t,\xi) = \begin{cases} -\sum_{|\gamma|=m-j-l} a_{m-j+1,\gamma}(t)\xi^{\gamma}\langle\xi\rangle^{j-m}, & \text{for } j \le m-l, \\ 0, & \text{otherwise,} \end{cases}$$

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for j = 1, ..., m. We easily see that the matrix B_{-l} has entries of order -l and the last l entries in the bottom row are equal to 0. Making use of this decomposition of B we can write $((Q_0^{(m)}B - B^*Q_0^{(m)})V, V)$ as

$$\sum_{l=0}^{m-1} ((Q_0^{(m)} B_{-l} - B_{-l}^* Q_0^{(m)}) V, V).$$
(43)

Remark 2 Let $0 \le h \le m - 2$. Let us assume the Levi conditions (7) in the form (29) only on the B_{-l} -matrices up to level h, i.e., instead of (7) assume only that

$$\left|\sum_{l=0}^{\min(h,m-j)} B_{-l,j}\right|^2 \prec \sum_{i=1}^m |\sigma_{m-j}^{(m-1)}(\pi_i\lambda)|^2,$$
(44)

for j = 1, ..., m. In other words, we impose Levi conditions on the coefficients of the equation corresponding to the matrices B_{-l} up to l = h, leaving free the remaining lower order coefficients. Recall that

$$|((Q_0^{(m)}B_l - B_l^*Q_0^{(m)})V, V)| \le 2(m-1)!|\mathcal{W}B_{-l}V||\mathcal{W}V|,$$

where the matrix W is defined in Sect. 3. Under the assumption (44) and the bound (17) from below for the quasi-symmetriser we obtain for (43) the estimate

$$\begin{split} & \left| \sum_{l=0}^{m-1} ((Q_0^{(m)} B_{-l} - B_{-l}^* Q_0^{(m)}) V, V) \right| \\ & \leq \left| \sum_{l=0}^{h} ((Q_0^{(m)} B_{-l} - B_{-l}^* Q_0^{(m)}) V, V) \right| + \left| \sum_{l=h+1}^{m-1} ((Q_0^{(m)} B_{-l} - B_{-l}^* Q_0^{(m)}) V, V) \right| \\ & \leq C_3 E_{\varepsilon} + 2(m-1)! \sum_{l=h+1}^{m-1} |\mathcal{W}B_{-l}V| |\mathcal{W}V| \\ & \leq C_3 E_{\varepsilon} + C_4 \langle \xi \rangle^{-h-1} \varepsilon^{-(m-1)} E_{\varepsilon}. \end{split}$$

This leads to the energy estimate

$$\partial_t E_{\varepsilon}(t,\xi) \le (C_1 \varepsilon^{-2(m-1)/k} + C_2 \varepsilon \langle \xi \rangle + C_3 + C_4 \langle \xi \rangle^{-h-1} \varepsilon^{-(m-1)}) E_{\varepsilon}(t,\xi).$$

The Gevrey well-posedness result of Theorem 2 will still hold true under the relaxed Levi condition (44), e.g., if for $\varepsilon^{-2(m-1)/k} = \varepsilon \langle \xi \rangle$ one has

$$\langle \xi \rangle^{-h-1} \varepsilon^{-(m-1)} \le \langle \xi \rangle^{\frac{1}{\sigma}},$$

with $\sigma = 1 + k/[2(m-1)]$, that is if

$$h+1 \ge \frac{(m-1)(k-2)}{k+2(m-1)}.$$
(45)

In other words, for any fixed $k \in \mathbb{N}$ by involving sufficiently enough matrices B_{-h} in the Levi condition (44) (how many depend on the equation order *m* and the regularity *k* of the coefficients) one can still obtain G^s well-posedness for

$$1 \le s < 1 + \frac{k}{2(m-1)}.$$

More precisely, by rewriting (45) as

$$h \ge m - 2 - \frac{2m(m-1)}{k+2(m-1)},$$

we can take

$$h \ge h_0 := m - 2 - \left[\frac{2m(m-1)}{k+2(m-1)}\right],$$
(46)

for all $m \ge 2$ and $k \ge 2$.

We now focus on the case of second order equations.

Remark 3 From the definition of h_0 in (46), we see that

$$h_0 = -\left[\frac{4}{k+2}\right] = 0$$

if m = 2 and $k \ge 2$. This shows that in the case of second order equations under the hypothesis (6) it is enough to put Levi conditions on the matrix B_0 to prove the Gevrey well-posedness of the Cauchy problem (1) for $1 \le s < 1 + \frac{k}{2(m-1)}$.

Concluding, let's us consider the case of non-analytic but very regular coefficients.

Remark 4 Assume now that the equation coefficients are smooth and that m > 2. This implies that for any a > 0 we can take k large enough such that $\varepsilon^{-2(m-1)/k} \leq \langle \xi \rangle^a$. Hence, $\varepsilon^{-2(m-1)/k} \leq \langle \xi \rangle^{-h-1} \varepsilon^{-(m-1)}$ with h = 0. Setting then $\varepsilon \langle \xi \rangle = \langle \xi \rangle^{-1} \varepsilon^{-(m-1)}$ we get that under the Levi condition (44) with h = 0 the Cauchy problem (1) is well-posed in G^s with

$$1 \le s < 1 + \frac{2}{m-2}$$

In terms of the Gevrey order this result is worse than the one stated in Theorem 2 but it is obtained with Levi conditions only on the coefficients appearing in the matrix B_0 . We note that it is still better than the Bronstein's result due to the extra assumption (6) and the Levi condition (44) with h = 0. Let us give some explanation to the relaxed Levi conditions (44). In particular, using the definition of the terms $B_{-l,j}$ and omitting the squares as in (10), condition (44) is equivalent to

$$\sum_{l=0}^{\min(h_0,m-j)} \left| \sum_{\substack{|\gamma|=m-j-l}} a_{m-j+1,\gamma}(t)\xi^{\gamma} \right| \\
\leq C \sum_{i=1}^{m} \left| \sum_{\substack{1 \le l_1 < \dots < l_{m-j} \le m \\ l_h \neq i \ \forall h}} \lambda_{l_1}(t,\xi) \cdots \lambda_{l_{m-j}}(t,\xi) \right|,$$
(47)

where the powers of $\langle \xi \rangle$ cancel. By comparison with the Levi conditions (7) and the identity (8), it is clear that we impose conditions on less terms of the equation.

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Appendix

Here we show that the Levi conditions (7) can be expressed entirely in terms of the equation coefficients. This is due to the fact that the sum $\sum_{i=1}^{m} |\sigma_{m-j}^{(m-1)}(\pi_i \lambda)|^2$ is, modulo some multiplicative factor *C*, the $q_{j,j}$ entry of the symmetriser $Q_0^{(m)}(\lambda)$. Indeed, using the notations of Sect. 3 we have that

$$Q_0^{(m)}(\lambda) = \sum_{\rho \in \mathcal{P}_m} P_0^{(m)}(\lambda_\rho)^* P_0^{(m)}(\lambda_\rho),$$

where $P_0^{(m)}(\lambda_\rho)^* P_0^{(m)}(\lambda_\rho)$ is the matrix with the *j*, *k*-entry being

$$\sigma_{m-j}^{(m-1)}(\pi_m\lambda_\rho)\sigma_{m-k}^{(m-1)}(\pi_m\lambda_\rho).$$

It follows that

$$q_{j,j}(\lambda) \sim \sum_{i=1}^m (\sigma_{m-j}^{(m-1)}(\pi_i \lambda))^2,$$

where \sim stands for equality modulo a multiplicative positive constant. In particular, $q_{m,m}(\lambda) = m!$. Recalling that Q is the Bezout matrix associated to the polynomials

 $(P, \partial_{\tau} P)$ with

$$P(\tau, t, \xi) = \tau^{m} + \sum_{j=0}^{m-1} \sum_{|\gamma|=m-j} a_{m-j,\gamma}(t) \xi^{\gamma} \tau^{j},$$

from formula (1.13) in [12] we have that, for j = 1, ..., m - 1,

$$q_{j,j} = \langle \xi \rangle^{-2} \left(j b_{(j+1)}^2 - \sum_{p=1}^{m-j-1} 2p \, b_{(j-p+1)} b_{(j+p+1)} \right) + 2(m-j) b_{(2j+1-m)} \langle \xi \rangle^{-1},$$
(48)

where $b_{(k)} = -A_{(m-k+1)}(t,\xi)\langle\xi\rangle^{k-m}$, for k = 1, ..., m, and $b_k = 0$ for k < 1. Finally, inserting (48) into (29) we conclude that the Levi conditions (7) can be written as

$$\begin{split} \left| (A_{m-j+1} - A_{(m-j+1)})(t,\xi) \right|^2 \langle \xi \rangle^{2j-2m} \\ &\leq C \langle \xi \rangle^{-2} \left(j b_{(j+1)}^2 - \sum_{p=1}^{m-j-1} 2p \, b_{(j-p+1)} b_{(j+p+1)} \right) \\ &+ 2C(m-j) b_{(2j+1-m)} \langle \xi \rangle^{-1} \\ &= C \langle \xi \rangle^{-2} \left(j \left(A_{(m-j)}(t,\xi) \langle \xi \rangle^{j+1-m} \right)^2 \right. \\ &- \sum_{p=1}^{m-j-1} 2p \, A_{(m-j+p)}(t,\xi) \langle \xi \rangle^{2j+2-2m} A_{(m-j-p)}(t,\xi) \right) \\ &- 2C(m-j) A_{(2m-2j)}(t,\xi) \langle \xi \rangle^{2j-2m} \end{split}$$

or, equivalently,

$$|(A_{m-j+1} - A_{(m-j+1)})(t,\xi)|^{2} \le C(jA_{(m-j)}(t,\xi)^{2} - \sum_{p=1}^{m-j} 2p A_{(m-j+p)}(t,\xi)A_{(m-j-p)}(t,\xi)).$$

Here, we assume $A_{(j)} = 0$ if $j \ge m + 1$ or j < 0.

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