



Echo Chains as a Linear Mechanism: Norm Inflation, Modified Exponents and Asymptotics

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Abstract

In this article we show that the Euler equations, when linearized around a low frequency perturbation to Couette flow, exhibit norm inflation in Gevrey-type spaces as time tends to infinity. Thus, echo chains are shown to be a (secondary) linear instability mechanism. Furthermore, we develop a more precise analysis of cancellations in the resonance mechanism, which yields a modified exponent in the high frequency regime. This allows us, in addition, to remove a logarithmic constraint on the perturbations present in prior works by Bedrossian, Deng and Masmoudi, and to construct solutions which are initially in a Gevrey class for which the velocity asymptotically converges in Sobolev regularity but diverges in Gevrey regularity.

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1. Introduction

In this article our aim is to develop a further understanding of the long-time asymptotic behavior of the 2D Euler equations near Couette flow

$$\partial_t \omega + y \partial_x \omega + v \cdot \nabla \omega = 0 \tag{1}$$

in an infinite channel $\mathbb{T} \times \mathbb{R}$. As the velocity field is incompressible, all L^p norms of the vorticity ω are conserved and the equation is nonlinearly stable in this sense. Furthermore, the linearized problem exhibits weak convergence of the vorticity in L^2 and as a consequence in the linear problem $v - \langle v \rangle_x$ converges strongly in L^2 as $t \rightarrow \infty$. This behavior is known as linear inviscid damping in analogy to the similar phenomenon of Landau damping in plasma physics. While the linearized problem possesses an explicit solution and exhibits linear inviscid damping for any initial data in H^s , $s \geq 0$, the question of stability and asymptotic behavior of more general shear flows or the nonlinear problem have been a very active area of research in recent years. In particular, we mention the following publications:

- In [3] Bedrossian and Masmoudi established nonlinear inviscid damping for Gevrey 2 regular perturbations around Couette flow in an infinite periodic channel. Their method of proof has further been extended to the setting of Landau damping [4,9].
- These results were recently extended to the case of compactly supported Gevrey regular perturbations to Taylor–Couette flow in [6].
- In [5] the first author and Masmoudi constructed solutions of the Euler equations, such that the initial data is Gevrey 2 close to Couette flow, the solution is well-understood for a finite time and within that finite time the solution exhibits growth along echo chains consistent with a loss of Gevrey 2 regularity. Here, we in particular stress that due to the challenging control of nonlinear corrections, these solutions *do not* capture the full echo chains (the finite time includes about half of a chain) and hence *do not* rule out subsequent asymptotic stability. Furthermore, that work imposes a logarithmic smallness restriction (see Section 4 for a further discussion). Removing this restriction is a key challenge in establishing our modified scattering in Section 5.3.
- Concerning the problem of linear inviscid damping, the case of a finite-periodic channel was shown to behave qualitatively differently in a work by the second author [14] in that stability results are limited to (sharp) $H^{3/2-}$ or weighted H^2 , [12], Sobolev regularity of the vorticity. While sufficient to establish linear inviscid damping with the optimal decay rates of the velocity perturbation, this leaves a large gap compared to the Gevrey regularity requirement of existing nonlinear results. Indeed, in [7,8] Ionescu and Jia further impose a compact support assumption to establish linear stability in Gevrey spaces.

- This gap, together with the physical observation of damping, suggests that damping of the velocity field might be a more robust mechanism than the stability of the vorticity. A first modest step in this direction can be found in [16] where the second author considered special types of forcing of the 2D Euler and Navier-Stokes equations and showed that damping may persist despite instability of the vorticity.

In this article, we stress the point of view that for the question of nonlinear inviscid damping in addition to the question of linear stability of shear flows in Sobolev and Gevrey regularity, one should consider the following two observations:

1. Existing works study the infinite time asymptotic stability of the vorticity in higher Sobolev regularity or Gevrey regularity, from which (linear) inviscid damping then follows as a corollary. However, this is a strictly stronger condition than the physically observed phenomenon of inviscid damping, which is the convergence of the velocity field. Indeed, we show that the linearized problem exhibits solutions which exhibit norm inflation due to complete echo chains. Here, the velocity asymptotically converges despite the divergence of the vorticity as time tends to infinity.
2. In the study of nonlinear inviscid damping, the main challenge is given by cascades of resonances, also known as echo chains. This mechanism is, however, not present in the linearized problem around shear flows due to the decoupling structure in Fourier space of these equations. In this article, we identify echo chains as a *secondary linear* mechanism, where we linearize around an arbitrarily small low-frequency perturbation around Couette flow. More precisely, we show that in arbitrarily small neighborhoods of Couette flow (with respect to local norms) there are *traveling wave-like* solutions of the form

$$\begin{aligned}\omega(t, x, y) &= -1 + c \cos(x - ty), \\ v(t, x, y) &= (y, 0) - \nabla^\perp \frac{1}{1+t^2} c \cos(x - ty),\end{aligned}$$

and study (a simplification of) the linearized problem around this wave.

For simplicity of calculations and presentation, in this article we consider a single mode perturbation by $c \cos(x)$ with c small. However, we think that an extension to more general low-frequency perturbations, while technically tedious, should follow by similar arguments. We show that this (simplified) linearized problem exhibits the same echo chains and norm-inflation results as for the nonlinear problem, globally in time. Thus, *echo chains are a linear mechanism*. Furthermore, we identify a critical regularity threshold at which stability of the vorticity in Gevrey regularity fails, but where damping of the velocity nevertheless persists.

In order to introduce our model consider the Euler equations near Couette flow (1) and change to Lagrangian coordinates (with respect to Couette flow) $(x + ty, y)$. Then our equation is given by

$$\partial_t \omega + \nabla^\perp \Delta_t^{-1} \omega \cdot \nabla \omega = 0,$$

where $\Delta_t = \partial_x^2 + (\partial_y - t\partial_x)^2$ and ω denotes the perturbation. Let now P_N denote a Littlewood–Payley projection to a large dyadic scale N . Then the projection of the equation yields

$$\partial_t P_N \omega + P_N (\nabla^\perp \Delta_t^{-1} \omega \cdot \nabla \omega) = 0. \tag{2}$$

We may then further decompose the nonlinearity by frequency-localizing each factor:

$$\sum_{N_1+N_2 \approx N} P_N (\nabla^\perp \Delta_t^{-1} P_{N_1} \omega \cdot \nabla P_{N_2} \omega). \tag{3}$$

Here, we obtain several regimes. If $N_1 \approx N_2$, we may freely integrate by parts and trade the inverse Laplacian Δ_t^{-1} for time decay. Similarly, if $N_2 \approx N, N_1 \ll N_2$, we may use the incompressibility (to obtain cancellations at leading order) and trade two derivatives for decay of Δ_t^{-1} to obtain a time-decreasing energy estimate.

The main source of possible growth and instability in the nonlinear problem is thus given by the interaction of the *low-frequency part of the vorticity*, $N_2 \ll N$, and the *high-frequency part of the velocity*, $N_1 \approx N$. In our model we now fix the low frequency part $\omega_{\text{low}} = c \cos(x)$ with c small and consider the (simplified) linearized problem for a perturbation, which, with slight abuse of notation, we again denote by ω :

$$\partial_t \omega + c \sin(x) \partial_y \Delta_t^{-1} \omega = 0. \tag{4}$$

We stress that ω_{low} is an exact solution of the nonlinear Euler equations (see Lemma 1), which we call a *traveling wave-like solution*. We remark that ω_{low} is not stationary in Eulerian coordinates, but is stationary in coordinates $(x+ty, y)$ moving with Couette flow (we hence view it as a wave moving with the flow). Our model thus corresponds to a secondary linearization around $v = (y, 0) + c \nabla^\perp \frac{1}{1+t^2} \cos(x-ty)$, where we neglected the transport by $\frac{c \sin(x)}{1+t^2} \partial_y$ for simplicity.

Our main result is then that this model exhibits the same *echo chains* as the full nonlinear Euler equations near Couette flow (in Gevrey regularity, see [2,5]), but further exhibits *modified scattering* and linear *inviscid damping* in the sense the velocity perturbation strongly converges as $t \rightarrow \infty$, but that the vorticity diverges in H^s for any $s > -1$. The secondary linearization captures the nonlinear resonance mechanism and, for a critical class of data, exhibits damping and blow-up at the same time.

Theorem 1. (Summary) *Let $0 < c < 0.2$, then there exists $C > 0$ such that for any $s \in \mathbb{R}$ there exists solutions of (4) with $\omega_0 \in \mathcal{G}_{C, \frac{1}{2}}$ (Gevrey class, see Section 1.1.1) such that $\omega(t)$ converges to a limit ω_∞ in H^σ , $\sigma \leq s$, but diverges in H^σ , $\sigma > s$. In particular, in the case $s \geq 0$ we note that damping, that is asymptotic convergence of the velocity field, holds, while asymptotic stability in Gevrey regularity fails.*

Furthermore, the constant C is (almost) optimal in the sense that there exists a (larger) constant $C_0 > 0$ such that for any $C_1 > 2C_0$ if $\omega_0 \in \mathcal{G}_{C_1, \frac{1}{2}}$, then $\omega(t) \in \mathcal{G}_{C_1-C_0, \frac{1}{2}}$ for all times and $u(t)$ converges in $\mathcal{G}_{C_1-C_0, \frac{1}{2}}$.

As we discuss in the ensuing subsection, a key challenge in this result is that the (nonlinear) growth mechanism of echo chains stops after a finite time for any given chain. Hence, in order to obtain non-trivial asymptotic behavior like blow-up it is necessary to work with countably many echo chains of increasing length. In particular, one has to allow for arbitrarily long chains, which had been ruled out by logarithmic smallness constraints in prior works. A key effort of this article in Sections 4 and 5 is then to remove this constraint and show that solutions behave qualitatively differently in that high frequency regime than in the previously considered low frequency case.

1.1. Outline and Challenges

In their seminal work Bedrossian and Masmoudi [3] established nonlinear asymptotic stability of Couette flow under sufficiently small Gevrey 2 regular perturbations. Furthermore, they showed that the perturbation of the velocity field decays (at algebraic rates) as time tends to infinity – a phenomenon which is known as *inviscid damping* in analogy to Landau damping in plasma physics. Following this result, there has been much renewed interest in inviscid damping, both linear and nonlinear, also around other flows. In particular, it turns out that the linearized problem (also around more general shear flows) is stable already in relatively low Sobolev regularity and that, unless one assumes compact support of perturbations, geometries with boundary do not allow for (linear) stability in high Sobolev regularity.

The reason for this large gap in regularity requirements – Gevrey as opposed to Sobolev – is tied to a physical phenomenon called *fluid echoes* [13]. As we recall in Section 3 here two perturbations at frequencies k and l in x (and suitable frequencies in y) may interact by means of the nonlinearity and introduce an (at first quadratically small) correction at frequency $k + l$. This perturbation may then become resonant at a later time and may thus yield a large contribution to the velocity field. We say that the perturbations at frequency k and l result in an *echo* at frequency $k + l$ at that time. In [3] it is then estimated by means of a toy model that a chain of such echoes, where frequency k causes an echo at frequency $k - 1$, which in turn causes an echo at frequency $k - 2$, \dots could lead to Gevrey 2 norm inflation. Indeed in [5] the first author and Masmoudi showed that such a chain may indeed form in the nonlinear equations (at least from k to $\frac{k}{2}$).

Since the linearized equations around a shear flow (which is independent of x) decouple with respect to frequency in x , these equations do not include the echo mechanism and hence echo chains are commonly understood to be a nonlinear effect. In this article we argue that it might be useful to instead view them as a *secondary linear effect*, where one should linearize around **traveling wave-like** solutions

$$\omega(t, x, y) = -1 + c \cos(x - ty)$$

instead of around a shear flow. We observe that in coordinates $(x - ty, y)$ moving with Couette flow such a wave is stationary and concentrated on two modes $+1, -1$ (in x and 0 in y). The heuristic of our echo chain construction is then the following:

- We study the (simplified) linear problem around this wave, where we initially introduce a perturbation at frequency (k_0, η) .
- These equations include the echo mechanism in the sense that the *underlying wave* (at frequencies $(\pm 1, 1)$) interacts with the *perturbation* at frequency (k_0, η) , which results in *echoes* at frequencies $(k_0 + 1, \eta)$ and $(k_0 - 1, \eta)$ at around a resonant time $t \approx \frac{\eta}{k_0}$.
- In turn the contribution at frequency $(k_0 - 1, \eta)$ by the echo can be viewed as new perturbation at a new initial time.
- The interaction of this perturbation at frequency $(k_0 - 1, \eta)$ with the underlying wave results in echoes at frequencies $(k_0 - 2, \eta)$ and (k_0, η) at around the resonant time $\frac{\eta}{k_0 - 1}$.
- We may iterate along this chain until we reach the last resonance at time $\frac{\eta}{1}$, which corresponds to the largest norm inflation.
- Furthermore, in Theorem 9 we combine infinitely many such chains to show that this inflation is sharp in the sense that there exist Gevrey 2 regular initial data which exhibit infinitely many echo chains and asymptotically blow up in H^s , but still converge in $H^{s-\varepsilon}$. *Damping of the velocity field persists despite blow-up of the vorticity*. When requiring higher Gevrey 2 regularity, that is with a larger constant, this growth is hidden by the rapid frequency decay and the solution remains Gevrey regular for all times and in the limit $t \rightarrow \infty$.

In order to make this heuristic rigorous we first, equivalently, formulate our equation as an infinite ODE system (7) in Section 2, which by the structure of the underlying wave includes only nearest neighbor interaction and is of the form

$$\partial_t \omega(t, l, \eta) + a(t, l + 1, \eta)\omega(t, l + 1, \eta) - a(t, l - 1, \eta)\omega(t, l - 1, \eta) = 0$$

for all l (the equation decouples with respect to η). The main aim of this article, which culminates in Theorem 8 of Section 5.2, is to show that on any time interval I centered around the resonant time $\frac{\eta}{k}$ (see Section 5) *the evolution is largely determined by just the three modes $k + 1, k, k - 1$* . As an important preliminary step, we hence study the associated homogeneous problem for these three modes (that is, neglecting all other modes in the evolution) in Section 5.1. Subsequently, in Section 5.2 we show by a bootstrap approach that the full problem can be approximated by (suitable corrections of) the homogeneous problem, where we separately consider three parts I_1, I_2, I_3 of the time interval I , each in a (sub)subsection.

A main challenge here is given by the fact that already in the homogeneous problem the evolution of the system is highly non-trivial, since the coefficients are time-dependent and potentially very large when considering frequencies η which are very large (compared to k^2). We stress that this regime is essential, since any single echo chain stops after the last resonant time $\frac{\eta}{1}$. In particular, we stress that previous works [3,5] included a **smallness constraint** of the form

$$\log(\eta) \leq C \tag{5}$$

(see equation (10) in Section 3). Any size constraint on η , however, means that one can only deduce norm inflation up to *finite* time (not asymptotic behavior) and furthermore for finite frequencies all norms are equivalent. Therefore, under such

a constraint it is *not* possible to distinguish between norm inflation with respect to different norms and there also is *no* blow-up. A key effort of this article thus lies in removing this constraint and studying the high frequency regime. We first show in Section 4 for a simplified model, that this is not only a technical challenge, but asymptotics indeed *qualitatively differ* in this regime and hence simpler methods of proof (such as in Section 3.1) necessarily fail.

We remark that Section 4 is not essential for the proof of our main result, but unlike the exact model considered in Section 5 allows us to construct explicit solutions (in terms hypergeometric functions and power laws) and hence to show in a more transparent way *why* at high frequencies we obtain a different exponent in the frequency dependence than might be expected from simpler heuristics. This model thus serves as an important preliminary step and heuristic, preceding the analysis of the full model (or the three mode model) in Section 5.2.

The remainder of our article is structured as follows:

- In Section 2 we show that for any $c \in \mathbb{R}$ the pair

$$\begin{aligned}\omega_{\text{low}}(t, x, y) &= -1 + c \cos(x - ty), \\ v_{\text{low}}(t, x, y) &= (y, 0) - \nabla^\perp \frac{1}{1+t^2} c \cos(x - ty),\end{aligned}$$

is a classical solution of the Euler equations (with infinite energy). While not stationary in Eulerian coordinates, we note that ω_{low} is stationary in coordinates $(x + ty, y)$. We call such a function a *traveling wave-like* solution. In particular, we observe that choosing c sufficiently small any small H_{loc}^s (or even locally analytic) neighborhood of Couette flow contains traveling wave-like solutions. Our Theorem 1 then argues that the (simplified) linearization around such waves captures much more of the nonlinear dynamics than the linearization around Couette flow and in particular captures Gevrey norm inflation, “nonlinear” resonances (called fluid echoes) and blow-up. More precisely, our problem (4) then corresponds to a simplification of the linearized Euler equations around this state, where we drop a term involving $\frac{\sin(x)}{1+t^2}$, which we expect to be negligible. Section 2 serves to introduce this model, some changes of coordinates and to establish local and global well-posedness of the model (with highly sub-optimal exponential growth bounds).

- In Section 3 we introduce the echo chain mechanism for a model problem similar to the one considered in [3]. Furthermore, we show that under a smallness assumption (5) on η and c similar to the one imposed in [3,5], also the full model exhibits echo chains and norm inflation. We remark that this part of the article is self-contained and allows for a much shorter and simpler proof than the general case but crucially relies on the smallness assumption. Since this condition limits considerations to finite η (where all norms are equivalent and norm inflation stops after a finite time), a key challenge and improvement of the this article is hence to remove this restriction in the following sections.
- In Section 4 we build an improved model which also takes into account the interaction of neighboring modes in the regime where η is much larger than allowed by the smallness assumption. In particular, we show that the smallness

restriction was not only a technical constraint, but that that the growth of solutions qualitatively differs in the large frequency regime and is connected to estimates on the 1D Schrödinger problem with scaling-critical potentials. The model character of the problem here allows for a more transparent discussion of the growth and decay properties of solutions on different time intervals and illustrates how their interaction leads to a modified dependence on the large frequency parameter.

- In Section 5 we then show that the full problem satisfies the same modified growth. In particular, we prove that for large frequencies η with respect to y (compared to c^{-1}) the dependence on η differs from the one in [3,5] (it is given by a modified power law). Here Section 5.1 introduces the so-called *three mode model*, which is a restriction of the full model to the resonant mode and its neighbors at around the resonant time $t \approx \frac{\eta}{k}$. In Section 5.2 we then show that these three modes determine the growth of the full model during this time interval, which is the main part of our proof. In Section 5.3 we combine the growth on each time interval into a chain to build solutions which exhibit norm inflation. Furthermore, we construct solutions which are global in time and exhibit inviscid damping (that is convergence of the velocity) but whose vorticity asymptotically diverges.

We remark that in the high regularity, small c , η regime the nonlinear problem and our model problem both asymptotically approach transport dynamics. It remains a challenging problem to determine how the full nonlinear Euler equations near Couette flow behave outside this regime and whether they also exhibit modified asymptotics similar to this linear model.

1.1.1. Notation We use the following notational conventions:

- The spatial Fourier transform of a function $\omega(t, x, y) \in L^2(\mathbb{T} \times \mathbb{R})$ is denoted by $\tilde{\omega}(t, k, \eta) \in L^2(\mathbb{Z} \times \mathbb{R})$.
- The Gevrey class $\mathcal{G}_{C, \frac{1}{s}}$ is defined in terms of the Fourier transform, that is $u \in \mathcal{G}_{C, \frac{1}{s}}$ if $\exp(C|\eta|^s)\tilde{u} \in L^2$. Here, we omit the k dependence in the exponent since for the functions we consider only the region $|k| \leq |\eta|$ is of interest.
- We use $a \lesssim b$ to denote that there exists an absolute constant $C > 0$ such that $|a| \leq C|b|$. In particular, we omit absolute value signs in our notation.
- In our calculations $c \in (0, 0.2)$ and $\eta \in \mathbb{R}$ can be treated as arbitrary but fixed parameters. While the low frequency regime allows for rather simple arguments (see Section 3), in the high frequency regime where $|\eta|$ is much larger than c^{-1} much finer control of cancellations is necessary.
- There we use $a \approx b$ to denote that there exists a constant $C_1 > 1$, such that $C_1^{-1}a \leq b \leq C_1a$ such that $|C_1 - 1| \leq 0.01$ for the regime of c and η we are considering (with a possibly even smaller deviation for c smaller and/or η larger). For example, in the regime of large η it holds that $\eta^{2-c} + \eta \approx \eta^{2-c}$.

2. Local Wellposedness and Asymptotic Stability

Our choice of the low-frequency vorticity $\omega_{low}(t, x, y) = -1 + c \cos(x + ty)$ (including Couette flow and Eulerian coordinates) is motivated by the Fourier coupling introduced by $\cos(x)$ and the fact that ω_{low} is a solution of the nonlinear Euler equations, which for small c can be considered as a very small analytic perturbation of Couette flow.

Lemma 1. *Let $c \in \mathbb{R}$ and define*

$$\begin{aligned} \omega_{low}(t, x, y) &= -1 + c \cos(x - ty), \\ v_{low}(t, x, y) &= (y, 0) - \nabla^\perp \frac{1}{1+t^2} c \cos(x - ty). \end{aligned}$$

Then (ω_{low}, v_{low}) is a classical solution of the Euler equations on $\mathbb{T} \times \mathbb{R}$ (with infinite energy) and $\omega(t, x + ty, y)$ does not depend on time.

In analogy to dispersive equations we call ω_{low} a *traveling wave-like* solution.

Proof. In the study of stationary solutions to the Euler equations one observes that the equation

$$0 = v \cdot \nabla \omega = \nabla^\perp (-\Delta)^{-1} \omega \cdot \nabla \omega$$

implies that the gradients of the stream function $\phi = (-\Delta)^{-1} \omega$ and of the vorticity are parallel. Therefore, they locally share level sets and locally one can express ω as a function of ϕ . Conversely, if there exists a smooth function f such that $\omega = f(\phi)$, then by the chain rule $v \cdot \nabla \omega = 0$ and one obtains a stationary solution (see for example [1]).

In the following we use the fact that $\cos(x - ty)$ is an eigenfunction of the the Laplacian and adapt this method to construct time-dependent solutions $\omega(t, x, y) = \omega(0, x - ty, y)$ which move with Couette flow. We observe that for every $t \in \mathbb{R}$ the stream function of the perturbation $\frac{1}{1+t^2} c \cos(x - ty) = (-\Delta)^{-1} c \cos(x - ty)$ and the perturbation of the vorticity $c \cos(x - ty)$ have the same level sets (and are even colinear). Therefore the gradients are parallel and

$$\nabla^\perp \frac{1}{1+t^2} c \cos(x - ty) \cdot \nabla \omega_{low} = 0.$$

Inserting this fact into the Euler equations it follows that

$$\partial_t \omega_{low} + v_{low} \cdot \nabla \omega_{low} = \partial_t \omega_{low} + y \partial_x \omega_{low} = 0,$$

which concludes the proof. \square

In our equation (4) we now consider the linearization around ω_{low} , where we omit the transport by $\frac{c \sin(x)}{1+t^2} \partial_y$ for simplicity. Since our equation (4) involves multiplication by a sine, it has an explicit characterization in Fourier variables:

$$\begin{aligned} \partial_t \tilde{\omega}(t, k, \eta) + \frac{c\eta}{(k-1)^2 + (\eta - (k-1)t)^2} \tilde{\omega}(t, k-1, \eta) \\ - \frac{c\eta}{(k+1)^2 + (\eta - (k+1)t)^2} \tilde{\omega}(t, k+1, \eta) = 0. \end{aligned}$$

In particular, we note that this equation decouples with respect to η (while linearizations around shear flows instead decouple with respect to k). Hence, we may consider $\eta \in \mathbb{R}$ to be a given parameter and introduce a new time variable τ :

$$t = \tau \eta. \tag{6}$$

With respect to this variable (4) can be equivalently expressed as

$$\begin{aligned} \partial_\tau \omega(\tau, k, \eta) + c \frac{1}{(k-1)^2} \frac{1}{\eta^{-2} + (\frac{1}{k-1} - \tau)^2} \omega(\tau, k-1, \eta) \\ - c \frac{1}{(k+1)^2} \frac{1}{\eta^{-2} + (\frac{1}{k+1} - \tau)^2} \omega(\tau, k+1, \eta) = 0. \end{aligned} \tag{7}$$

In particular, we note that in this formulation all *critical times* are given by $\frac{1}{j}$ for $j \in \mathbb{N}$ and thus independent of η . Furthermore, there are no critical times after time $\tau = 1$, which implies that norm inflation or instability results are restricted to the evolution on small times and that the evolution is asymptotically stable after time $\tau = 2$.

Proposition 1. *Let X be a weighted L^2 (or l^2) space on the Fourier side (e.g. a fractional Sobolev or Gevrey space) such that the weight $\rho(k, \eta)$ satisfies $\sup_\eta \frac{\rho(k \pm 1, \eta)}{\rho(k, \eta)} \leq C_1 < \infty$. Then there exists a constant $C = C(cC_1)$ such that any solutions of (7) satisfy*

$$\|\omega(\tau)\|_X \leq C \|\omega(2)\|_X$$

for all $\tau \geq 2$. Furthermore, $\omega(\tau)$ strongly converges in X to a limit $\omega_\infty \in X$ as $t \rightarrow \infty$.

Proof of Proposition 1. We may estimate the multipliers by

$$\begin{aligned} \frac{c}{l^2} \frac{1}{\eta^{-2} + (\frac{1}{l} - \tau)^2} \\ \leq c \frac{1}{0 + (2 - \frac{1}{l} + (\tau - 2))^2} \\ \leq \frac{c}{(1 + (\tau - 2))^2}, \end{aligned}$$

irrespective of the size of η or of $l \in \mathbb{Z}, l \neq 0$. We further remark that if $l\eta < 0$ no restriction is necessary, since then there cannot be any cancellation. We may estimate

$$\begin{aligned} \partial_\tau \|\omega\|_X &\leq \frac{c}{1 + (\tau - 2)^2} (\|\omega(k+1)\|_X + \|\omega(k-1)\|_X) \\ &\leq \frac{2cC_1}{1 + (\tau - 2)^2} \|\omega\|_X. \end{aligned}$$

The result then immediately follows from integrating this inequality with

$$C := \exp\left(2cC_1 \left\| \frac{1}{1 + (\tau - 2)^2} \right\|_{L^1}\right).$$

□

We remark that the admissible class of spaces X includes fractional Sobolev spaces H^s and Gevrey spaces $\mathcal{G}_{C, \frac{1}{s}}$, but not the homogeneous fractional Sobolev spaces due to the quotient $\frac{|\eta \pm 1|^s}{|\eta|^s}$ degenerating for $\eta \downarrow 0$.

In order to analyze possible norm inflation, in the following we hence focus on characterizing the evolution on the time interval $[0, 2]$. Here, as a first easy, non-optimal estimate, we may roughly control the multiplier by

$$\sup_{l \neq 0} \frac{|c\eta^2|}{l^2} = c\eta^2,$$

and thus obtain

Theorem 2. *Let X be as in Proposition 1. Then if the initial data is localized on the mode η , it holds that*

$$\|\omega(\tau)\|_X \leq \exp(2C_1 c\eta^2 \min(2, \tau)) \|\omega_0\|_X,$$

for all $\tau \geq 0$. Furthermore, $\omega(\tau)$ strongly converges in X as $\tau \rightarrow \infty$.

Proof of Theorem 2. By Proposition 1 it suffices to establish a bound for $\tau \leq 2$. Here, we estimate the Fourier multiplier in (7) by $c\eta^2$ and thus obtain that, for ω_0 localized at frequency η , it holds that

$$\begin{aligned} \|\partial_\tau \omega\|_X &\leq c\eta^2 (\|S_{+1}\omega\|_X + \|S_{-1}\omega\|_X) \\ &\leq 2C_1 c\eta^2 \|\omega\|_X, \end{aligned}$$

where $S_{\pm 1}$ denotes a shift in the k dependence. The result then follows by Gronwall’s lemma. \square

We stress that these estimates are far from optimal. However, they allow us to control the evolution until a small positive time $\tau_0 > 0$, after which a more detailed study of resonance chains establishes sharper bounds.

Indeed, in Section 5, we show that this linearized model attains the same Gevrey 2 class norm inflation results as the nonlinear problem around Couette flow considered in [3, 5]. Furthermore, the solutions exist globally in time and exhibit modified scattering and damping.

In order to introduce ideas, in the following Section 3 we consider the setting where c is very small compared to η^{-1} . This setting allows for highly simplified proofs and serves to introduce the resonance mechanism in a clear way. In Section 4 we then introduce an improved “toy model”, which includes further cancellation mechanisms compared to the one studied in [3] and shows what modifications to asymptotic convergence should be expected for large η (compared to c^{-1}). The results of Section 5 then show that the full linear problem exhibits the same kind of growth as this model and we further discuss implications for scattering and damping.

3. Echoes, Paths and Norm Inflation

Theorem 2 of the preceding section proves wellposedness of the evolution equation for frequency localized initial data, but provides a very rough upper bound on the possible growth. In this section we study the evolution of equation (7)

$$\begin{aligned} \partial_\tau \omega(\tau, k, \eta) + c \frac{1}{(k-1)^2} \frac{1}{\eta^{-2} + (\tau - \frac{1}{k-1})^2} \omega(\tau, k-1, \eta) \\ - c \frac{1}{(k+1)^2} \frac{1}{\eta^{-2} + (\tau - \frac{1}{k+1})^2} \omega(\tau, k+1, \eta) = 0. \end{aligned} \tag{8}$$

for $0 \leq \tau \leq 2$ in more detail in order to obtain a finer description of the associated norm-inflation mechanism. More precisely, in this section we consider the setting where η satisfies a smallness assumption (10) as in [3,5], which allows for a short and simple proof of norm inflation. We remark that this section is self-contained and that the low frequency results of this section are not required for our main results. It is instead intended to provide a simple introduction to the main growth mechanism and to motivate why removing the smallness assumption is a highly non-trivial, challenging problem.

In order to develop a heuristic for the growth mechanism, formally consider the equation with $\omega(\tau, k-1, \eta)$, $\omega(\tau, k+1, \eta)$ frozen in time. Then integrating the equation we obtain that, for two times $\tau_0 < \tau_1$, it holds that

$$\begin{aligned} \omega(\tau_1, k, \eta) - \omega(\tau_0, k, \eta) \\ \approx \omega(\tau_0, k-1, \eta) \int_{\tau_0}^{\tau_1} c \frac{1}{(k-1)^2} \frac{1}{\eta^{-2} + (\tau - \frac{1}{k-1})^2} d\tau \\ - \omega(\tau_0, k+1, \eta) \int_{\tau_0}^{\tau_1} c \frac{1}{(k+1)^2} \frac{1}{\eta^{-2} + (\tau - \frac{1}{k+1})^2} d\tau. \end{aligned}$$

In reference to fixed point iterations, we also call this heuristic the first Duhamel iteration (which inserts the initial datum into the fixed point map. See also Section 3 for higher iterations.). If τ_0 and τ_1 are suitably chosen the growth factor can thus be expected to be comparable to

$$\frac{1}{(k-1)^2} \int_{-\infty}^{\infty} \frac{c}{\eta^{-2} + (\tau - \frac{1}{k-1})^2} d\tau = c\pi\eta \frac{1}{(k-1)^2} = c\pi \frac{\eta}{(k-1)^2}.$$

Similarly to the nonlinear equations studied in [3,5], the main norm inflation mechanism of (7) is then for mode (k, η) to induce growth of the mode $(k-1, \eta)$ at time $\tau \approx \frac{1}{k}$, which in turn induces growth of the mode $(k-2, \eta)$ at time $\tau \approx \frac{1}{k-1}$. Iterating this heuristic along a chain until time $\tau \approx 1$, this suggests a total growth factor of

$$\frac{\eta}{k^2} \frac{\eta}{(k-1)^2} \cdots \frac{\eta}{1^2} = \frac{\eta^k}{(k!)^2}. \tag{9}$$

This factor achieves its maximum with respect to k for $k \approx \sqrt{\eta}$, which yields $\exp(C\sqrt{\eta})$ and thus a norm inflation result consistent with Gevrey 2 regularity.

While single echoes or small numbers of echoes are well-studied in the physical and numerical literature [11, 13] we stress here that long chains of echoes correspond to Gevrey regularity. Indeed, one of the main challenges in the following will be to construct infinitely many chains of arbitrarily long lengths in such a way that the evolution can be controlled for all times.

Compared to the nonlinear results of [5], we highlight the following differences of the present setting:

- Due to the difficulties involved in controlling a nonlinear evolution around growing perturbations, [5] establishes such a growth chain until time $t \approx \frac{1}{2k} \approx \frac{1}{2\sqrt{\eta}} \ll 1$. Our results established for the linear equation (7) instead capture a full echo chain as well as the subsequent asymptotic stability as $\tau \rightarrow \infty$.
- The more explicit structure of the associated Duhamel iteration scheme allows us to establish more precise upper and lower bounds and in particular quantify the dependence on the parameter c , i.e. the size of the low-frequency part. Furthermore, as we discuss in Section 4, the evolution for η much larger than k^2 introduces additional logarithmic corrections, which we show in Section 5 to result in a correction to the asymptotic evolution. In particular, prior works included a constraint of the form

$$c \log(1 + |\eta|) \lesssim 1,$$

where a change in exponent by c is then comparable to multiplication by a constant.

- Identifying the sharp growth behavior and removing this constraint allows us to construct families of initial data, for which the velocity field converges while the vorticity diverges (see Section 5)! We view this as an indication that nonlinear inviscid damping, understood as asymptotic convergence of the velocity field, might still hold in lower regularity despite norm inflation results for the vorticity.

As a first step, in the following theorems we improve the bounds by $\exp(C\eta^2)$ to Gevrey 2-type estimates of the above form. Similarly to [2, 5] in this section we impose a logarithmic smallness condition

$$c \leq C \ln(|\eta|)^{-1}, \tag{10}$$

which greatly simplifies this proof. In particular, it allows us to clearly demonstrate the growth mechanism, obtain explicit (approximate) solutions and point out where these methods of proof would fail for high frequencies.

We stress that for any fixed c this smallness condition prevents us from considering (sequences of) arbitrarily large η , which is necessary to establish instability and modified scattering results. A key effort of this article in Sections 4 and 5.2 is to remove this constraint. Here, it turns out that large frequencies not only pose technical challenges, but that for larger η the asymptotics indeed differ. In Section 5.3 we show that taking into account these modified asymptotics our linear system exhibits norm inflation and echo chains for all η , as well as modified scattering.

Theorem 3. (Resonance chain) *Let $\eta > 1$, $l \in \mathbb{N}$ with $\frac{\eta}{l^2} \geq 1$ and let $0 < c \leq \min((\ln(1 + \eta^2))^{-1}, \frac{1}{100})$. Let further $\tau_0 = \frac{1}{2}(\frac{1}{l+1} + \frac{1}{l})$ and $\tau_1 = 1.5$. Then the solution $\omega(\tau)$ of (7) with*

$$\omega(\tau_0, k) = \delta_{k,l}$$

satisfies

$$\begin{aligned} \omega(\tau_1, k) &\geq c^l \left(1 - \frac{c}{1-c}\right)^l \frac{\eta^l}{(l!)^2} \text{ if } k \in \{1, 3\}, \\ \omega(\tau_1, k) &\leq c^l \left(1 + \frac{c}{1-c}\right)^l \frac{\eta^l}{(l!)^2} \text{ if } k \in \{1, 3\}, \\ |\omega(\tau_1, k)| &\leq c^{|k-l|} c^l \left(1 + \frac{c}{1-c}\right)^l \frac{\eta^l}{(l!)^2} \text{ if } k \notin \{1, 3\}. \end{aligned}$$

Furthermore, there exists $\omega_\infty(k)$ such that $\omega(t, k) \rightarrow \omega_\infty(k)$ as $\tau \rightarrow \infty$ and $\|\omega_\infty\|_X \leq C_X \|\omega(\tau_1)\|_X$. In particular, the associated velocity field converges as $\tau \rightarrow \infty$.

We remark that such a result also yields bounds for frequency-localized initial data globally in time by using Theorem 2 to control the evolution until time τ_0 . In order to prove Theorem 3 we use precise description of the evolution operator around single resonances $\tau \approx \frac{1}{k}$ for $k = l, l - 1, \dots, 1$, which we then combine to establish the over all growth.

3.1. Single Resonance Estimates

We note that in equation (7) modes only directly interact with their nearest neighbors. Thus, if we prescribe that $\omega(\tau_0) = e^{ik_0x + i\eta y}$, the value of $\tilde{\omega}(\tau_1, k, \eta)$ involves sequences $\gamma = (k_0, k_0 - 1, \dots, k)$, where a mode γ_i influences a mode γ_{i+1} by nearest neighbor interaction. Indeed, any finite Duhamel iteration can be associated with a sum over all such sequences with an upper bound on their length and integrals of the form

$$\iint_{\tau_0 \leq s_1 \leq s_2 \leq \dots \leq \tau_1} \prod_{i=1}^{|\gamma|} \frac{1}{\gamma_i^2} \frac{c}{\eta^{-2} + (s_i - \frac{1}{\gamma_i})^2} ds_i. \tag{11}$$

In order to simplify discussion, we introduce some notation.

Definition 1. Let $k, k_0 \in \mathbb{Z}$ with $\text{sgn}(k) = \text{sgn}(k_0)$. Then a *path* γ from k_0 to k is a finite sequence $\gamma = (\gamma_i)_{i=1}^n$ with $\gamma_1 = k_0, \gamma_n = k$ and $|\gamma_{i+1} - \gamma_i| = 1$. We call the integral (11) the *integral* associated to γ , $I[\gamma]$ and call $n - 1 =: |\gamma|$ the *length* of γ .

We remark that equation (7) introduces a post-multiplication by $\sin(x)$ and a single Duhamel iteration hence results in paths $(k_0, k_0 + 1)$ and $(k_0, k_0 - 1)$, which we thus denote as paths of length 1.

Given a time interval $\frac{1}{l+1} < \tau_0 < \frac{1}{l} < \tau_1 < \frac{1}{l-1}$, we observe that integrals tend to be small unless $\gamma_i = l$ for some i . Indeed, if $\gamma_i \neq l$, then $|\frac{1}{\gamma_i} - \frac{1}{l}| \geq \frac{c}{l^2}$ and we may estimate each integral with respect to s_i by

$$\begin{aligned} & \frac{1}{l^2} \int_{|\tau - \frac{1}{l}| \geq \frac{1}{2l(l-1)}} \frac{c}{\eta^{-2} + (\tau - \frac{1}{l})^2} d\tau \\ & \leq \frac{1}{l^2} \int_{s \geq \frac{1}{2l(l-1)}} \frac{c}{\tau^2} d\tau \leq 4c, \end{aligned} \tag{12}$$

uniformly in η and in l .

Definition 2. Let $l \in \mathbb{N}$ and define $T_0 := \frac{1}{2}(\frac{1}{l+1} + \frac{1}{l}) < \frac{1}{l} < \frac{1}{2}(\frac{1}{l-1} + \frac{1}{l}) =: T_1$. Let further $k, k_0 \in \mathbb{N}$ and let γ be a path from k_0 to k . We call γ_i *resonant* if $\gamma_i = l$ and *non-resonant* else. A path γ is called *non-resonant* if γ_i is non-resonant for all i .

Proposition 2. (Single resonance bound) *Let $l \in \mathbb{N}$, $\eta > 0$ and let*

$$\tau_0 := \frac{1}{2} \left(\frac{1}{l+1} + \frac{1}{l} \right) < \frac{1}{l} < \frac{1}{2} \left(\frac{1}{l} + \frac{1}{l-1} \right) =: \tau_1,$$

and suppose that the constant $c > 0$ is sufficiently small that $2c \leq \ln(1 + \eta)^{-1}$.

1. Let $k_0 \in \mathbb{N}$, $k_0 \neq l$ and let $\omega(\tau, \cdot)$ be the solution of (7) with

$$\omega(\tau_0, k) = \delta_{k, k_0}.$$

Then it holds that

$$|\omega(\tau_1, k) - \omega(\tau_0, k)| \leq C(c^{|k-l|+|k_0-l|} \frac{\eta}{l^2} + c^{|k-k_0|})$$

for all $k \in \mathbb{N}$. Here, with slight abuse of notation $|k - l| = \max(|k - l|, 2)$ denotes the length of the shortest path connecting k and l . This agrees with the absolute value if $k \neq l$, but is 2 if $l = k$, since the shortest path $(k, k \pm 1, k)$ has length 2.

2. Let $\omega(\tau, \cdot)$ be the solution of (7) with

$$\omega(\tau_0, k) = \delta_{k, l}.$$

Then it holds that

$$\begin{aligned} \omega(\tau_1, l \pm 1) & \approx \frac{\eta}{l^2}, \\ |\omega(\tau_1, l) - 1| & \leq Cc^2 \frac{\eta}{l^2}, \\ |\omega(\tau_1, k)| & \leq C(c^{|k-l|} \left(1 + \frac{\eta}{l^2} \right)) \text{ else.} \end{aligned}$$

Proof of Proposition 2. We argue by a Duhamel-iteration approach and summing over all paths starting in k_0 and ending in k . Here, we first note that if a path is non-resonant, we may estimate its contribution by $c^{|\gamma|}$. If we then estimate the number of such paths from above by $2^{|\gamma|}$, the contribution by all non-resonant paths is bounded by

$$\sum_{|\gamma| \geq \text{dist}} (2c)^{|\gamma|} = \frac{1}{1 - 2c} (2c)^{\text{dist}},$$

where dist is the length of the shortest (non-resonant) path connecting k and k_0 . Our main challenge in the following is thus going to be to control resonances and in particular the interaction between multiple resonances (that is, $\gamma_i = l$ for several indices i).

Ad 1 Let $k_0 \neq l$ be given and let $k \in \mathbb{N}$. Let further γ be a path starting in k_0 and ending in k and denote $j = |\gamma|$. We have already discussed the case of all non-resonant paths above, so suppose that $\gamma_i = l$ for some i and let $(i_\kappa)_{\kappa=1}^n$ denote all such indices. We note that we first have to reach l starting from k_0 and thus $i_1 \geq |k_0 - l|$. Similarly it follows that $|i_n - j| \leq |k - l|$. Furthermore, $i_\kappa \leq i_{\kappa+1} - 2$, since consecutive entries in a path are distinct.

In order to bound $I[\gamma]$, first consider two subsequent resonances indices i_κ and $i_{\kappa+1}$. Keeping these indices fixed, we first integrate over the intermediate indices s^1, \dots, s^l and consider

$$\int_{s_{i_\kappa} \leq s^1 \leq \dots \leq s^l \leq s_{i_{\kappa+1}}} \prod_{j=1}^l \frac{1}{\gamma_j^2} \frac{c}{\eta^{-2} + \left(s^j - \frac{1}{\gamma_j}\right)^2} ds_j \leq (Cc)^l |s_{i_{\kappa+1}} - s_{i_\kappa}|,$$

where we used that by assumption all these $\gamma_j \neq l$ are non-resonant and can be bounded by powers of c and bounded one of the integrals by $c(s_{i_{\kappa+1}} - s_{i_\kappa})$.

Repeating this argument for all κ , we only need to consider the integrals with respect to the resonances:

$$(cC)^j \iint_{T_0 \leq s_{i_1} \leq s_{i_2} \leq \dots \leq s_{i_n} \leq T_1} \frac{1}{\eta^{-2} + \left(s_{i_n} - \frac{1}{\eta}\right)^2} \prod_{\kappa=1}^{n-1} \frac{s_{i_{\kappa+1}} - s_{i_\kappa}}{\eta^{-2} + \left(s_{i_\kappa} - \frac{1}{\eta}\right)^2}.$$

Rescaling and shifting all s . as $\eta^{-1}(s - \frac{1}{\eta})$, we obtain a factor η^{-j} from the Jacobian, a factor η^{-j+1} from the linear factors, and a factor η^{2j} from the denominators, and thus, in total:

$$(cC)^j \eta \iint_{\eta(T_0 - \frac{1}{\eta}) \leq s_{i_1} \leq s_{i_2} \leq \dots \leq s_{i_n} \leq \eta(T_1 - \frac{1}{\eta})} \frac{1}{1 + s_{i_n}^2} \prod_{\kappa=1}^{n-1} \frac{s_{i_{\kappa+1}} - s_{i_\kappa}}{1 + s_{i_\kappa}^2}.$$

It remains to estimate the integral over this now η dependent region.

Expanding the product in the numerator and looking at each summand separately we need to consider monomials in the numerator. If a factor s_{i_κ} does not

appear in the numerator we simply bound by

$$\int_{\mathbb{R}} \frac{1}{1 + s_{i_k}^2} = \pi.$$

If it appears once, we may compute

$$\int_s \frac{|s|}{1 + s^2} = \frac{1}{2} \ln(1 + s^2) + c.$$

If it appears twice, we bound

$$\int_{s_{i_{k-1}} \leq s_{i_k} \leq s_{i_{k+1}}} \frac{s_{i_k}^2}{1 + s_{i_k}^2} \leq s_{i_{k+1} - s_{i_{k-1}}},$$

but in this way i_k does not appear in the integral anymore and we can repeat the above argument with this resonance removed. Hence, we only need to consider monomials where each s_{i_k} appears either to power 1 or does not appear:

$$\iint_{\eta(T_0 - \frac{1}{l}) \leq s_{i_1} \leq s_{i_2} \leq \dots \leq s_{i_n} \leq \eta(T_1 - \frac{1}{l})} \prod_{i_k \text{ appears}} \frac{|s_{i_k}|}{1 + s_{i_k}^2} \prod_{i_k \text{ does not}} \frac{1}{1 + s_{i_k}^2}.$$

As discussed above the integrals over the “does not” case can be bounded by π^{j_2} , where $j_2 = j - j_1$ denotes the number of such cases. The integral over the first product is bounded by a power of a logarithm:

$$\begin{aligned} & \frac{1}{j_1!} \left(\frac{1}{2} \ln \left(1 + \left(\eta \left(T_1 - \frac{1}{l} \right) \right)^2 \right) + \frac{1}{2} \ln \left(1 + \left(\eta \left(T_0 - \frac{1}{l} \right) \right)^2 \right) \right)^{j_1} \\ & \approx \frac{1}{j_1!} \left(C \ln \left(1 + \frac{\eta}{l^2} \right) \right)^{j_1}. \end{aligned} \tag{13}$$

Recalling the additional prefactor $c^{|\gamma|}$, summing over all such j_1 then leads to a bound of this contribution by

$$\left(1 + \frac{\eta}{l^2} \right)^{Cc} = \exp \left(Cc \ln \left(1 + \frac{\eta}{l^2} \right) \right) \leq \exp(\tilde{C}),$$

where we used the logarithmic smallness assumption

$$c \ln \left(1 + \frac{\eta}{l^2} \right) \ll 1.$$

Hence, this exponential can be treated as a perturbation provided c and thus \tilde{C} is sufficiently small.

Ad 2 We proceed similarly as in case 1, but note that

$$\int_{\tau_0}^{\tau_1} \frac{1}{l^2} \frac{c}{\eta^{-2} + \left(s - \frac{1}{l} \right)^2} \approx c \frac{\eta}{l^2}.$$

All longer paths starting in l and ending in $l - 1$ have length at least 3 and hence the sum over all these paths can again be controlled by

$$\frac{c^2}{1 - c} \frac{c\eta}{l^2}.$$

We hence obtain the claimed comparability with factors $1 \pm \frac{c^2}{1-c}$. For $k = l$, we note that the shortest non-trivial path has length 2 and for $k \notin \{l - 1, l, l + 1\}$ we argue as before except that all paths start in $\gamma_1 = l$. \square

We remark that by using the linearity of the problem this further yields a convolution-type bound for $\omega(\tau_0)$ not being localized on a single mode. In the following we then combine these mode-wise upper and lower bounds as well as the local well-posedness of Theorem 2 to construct an explicit example of a function that exhibits the Gevrey norm inflation up to time $\tau = 1$ and afterwards is asymptotically stable and hence also exhibits inviscid damping. In Section 5 we show with considerable technical effort that a similar but distinct result also hold if η does not satisfy a logarithmic bound.

3.2. Proof of Theorems 3

Using the single-resonance results of Proposition 2 we are now ready to prove Theorem 3.

We proceed in multiple steps:

1. We use local well-posedness to control the evolution from time $\tau = 0$ up to a small positive time $\tau^* = \tau^*(c, \eta) \ll 1$.
2. We next choose $l \in \mathbb{N}$ such that $\frac{1}{l} > \tau^*$ and l is maximal with this property and prescribe $\omega(\tau)$ at the time $\tau_0 = \frac{1}{2}(\frac{1}{l+1} + \frac{1}{l}) < \frac{1}{l}$ to be localized at frequencies (l, η) . Proposition 2 then allows us to control the evolution up to time $\tau_1 = \frac{1}{2}(\frac{1}{l} + \frac{1}{l-1}) > \frac{1}{l}$, where we in particular obtain upper and lower bounds on the modes $l \pm 1$ and upper bounds on all other modes.
3. We then iterate this control another $l - 1$ times and establish upper bounds on all modes and lower bounds along our chain of resonances $(l, l - 1, l - 2, \dots, 1)$ until after time $\tau = 1$.
4. By our construction of the coordinate τ , after this time no resonances appear anymore and we may use Proposition 1 to control the long-time asymptotic behavior.

Proof of Theorem 3. We note that it suffices to establish upper and lower bounds on ω until time $\tau = 1.5$, since after that time asymptotic stability and convergence of the velocity field follow by Proposition 1.

Let thus $l \in \mathbb{N}$ be given, to be fixed later. By Theorem 2 we may find initial data such that at time $\tau_l := \frac{1}{2}(\frac{1}{l-1} + \frac{1}{l})$ it holds that

$$|\omega(\tau_l, l)| \geq 0.5 \max |\omega(\tau_l)| =: 0.5\theta. \tag{14}$$

We then apply Proposition 2 using the linearity and a triangle inequality to obtain a convolution bound. That is, for $k \notin \{l - 1, l, l + 1\}$, we obtain that

$$\begin{aligned} |\omega(\tau_{l+1}, k) - \omega(\tau_l, k)| &\leq C \left(c^{|k-l|} \left(1 + \frac{\eta}{l^2} \right) \right) |\tilde{\omega}(\tau_l, l)| \\ &\quad + \sum_{k_0 \neq l} |\tilde{\omega}(\tau_l, k_0)| C \left(c^{|k-l|+|k_0-l|} \frac{\eta}{l^2} + c^{|k-k_0|} \right) \\ &\leq \theta C \left(c^{|k-l|} \left(1 + \frac{\eta}{l^2} \right) \right) + \theta C c^{|k-l|} \frac{1}{1-c} \frac{\eta}{l^2} + \frac{\theta}{1-c}. \end{aligned}$$

If $k = l$, the first term is replaced by

$$C c^2 \frac{\eta}{l^2}.$$

Finally, if $k \in \{l - 1, l + 1\}$, we obtain that

$$|\omega(\tau_{l+1}, k) \mp C_{\pm} c \frac{\eta}{l^2} \omega(\tau_l, l) - \omega(\tau_l, k)| \leq \theta C c^{|k-l|} \frac{1}{1-c} \frac{\eta}{l^2} + \frac{\theta}{1-c},$$

where $C_{\pm 1} \approx 1$. Recalling that $\omega(\tau_l, l)$ achieves θ within a factor 2 and using that $c \frac{\eta}{l^2} \gg 1$, it follows that

$$\pm c \frac{\eta}{l^2} \omega(\tau_l, l) \approx \omega(\tau_{l+1}, l \pm 1) \geq 0.5 \max |\omega(\tau_{l+1})|.$$

also achieves the new maximum within a factor 2 and is larger than the previous maximum by a factor $c \frac{\eta}{l^2}$.

We may thus repeat our application of Proposition 2 iteratively until for $k = 1$, we obtain that

$$\omega(\tau_1, 1) \approx \frac{c^l \eta^l}{(l!)^2} \omega(\tau_l, l)$$

achieves the full growth along a chain and again is comparable to the supremum at that time. \square

While this result is quite useful and shows the growth mechanism, the logarithmic bound

$$c \ln \left(1 + \frac{\eta}{l^2} \right) \ll 1,$$

prevents us from considering η arbitrarily large. In particular, reverting the change of variables $t \mapsto \tau$, this means that after a final time $t = \max \eta$ there are no more resonances and the evolution is asymptotically stable. Furthermore, since we can only consider a finite set in η all norms are equivalent and thus the associated growth while suggestive of a Gevrey regularity class is consistent with any norm.

We remark that also [3,5] contain a similar constraint. While one might at first hope that this is a purely technical constraint, in our proof we saw various logarithmic terms appearing, where in particular the contribution by the path $\gamma = (l, l - 1, l, l - 1)$ is also bounded below. It is hence a large, non-negligible correction,

which can not be treated perturbatively. Thus, our main goal in the following is to understand how the dynamics change and to remove this restriction. To this end, in Section 4 we first introduce an improved model problem and new methods of proof, which allow us to consider η arbitrarily large. Then in Section 5 we show that this behavior also holds in the full problem.

4. An Improved Model Problem

In order to better study the effect of resonances and, in particular, the behavior for large η and/or large c , in the following we consider an abridged model which considers the evolution for $\frac{1}{l+1} < \tau < \frac{1}{l-1}$ and only considers the resonant mode l and its neighbor $l - 1$:

$$\begin{aligned} \partial_\tau \omega(\tau, l - 1) - \frac{c}{\eta^{-2} + (\tau - \frac{1}{l})^2} \frac{1}{l^2} \omega(\tau, l) &= 0, \\ \partial_\tau \omega(\tau, l) + \frac{c}{\eta^{-2} + (\tau - \frac{1}{l-1})^2} \frac{1}{(l - 1)^2} \omega(\tau, l - 1) &= 0. \end{aligned}$$

As remarked in Section 1.1, this section is not strictly necessary for the proof of our main result in Section 5. However, its model structure allows us to explain why the behavior at large frequencies necessarily strongly differs and how growth and decay on several time intervals connect to each other to yield a modified power law growth. Indeed in Section 5 we show that also the full model behaves similar to this two-mode model (or more accurately like the three-mode model involving the modes $l + 1, l, l - 1$), where our strategy of proof closely mirrors the one of Proposition 3 of this section (with each step of the proof now corresponding to a subsection). This section hence serves as an important preliminary step, highlighting the main challenges and illustrating the effects of large frequencies violating the smallness constraint.

The present model is very similar to the one considered in [3], except that our model does not estimate the coefficients by absolute values but rather keeps the signs, which corresponds to exploiting the real-valuedness of $\sin(x)$ (and hence anti-symmetry of the imaginary parts of the Fourier coefficients). We further remark that a similar model for a single echo has been studied in [11] both analytically and numerically, but not for chains and only in the small parameter regime.

As we will see in the following this yields important cancellation properties and further exposes a connection of this problem to the 1D Schrödinger problem with scaling critical potential.

As a simplification of calculations, we use that $\frac{1}{l-1}$ is non-resonant and approximate $\tau - \frac{1}{l-1} \approx \frac{1}{l^2}$ and thus consider the following *two-mode system*:

$$\begin{aligned} \partial_\tau u - \frac{c}{\eta^{-2} + \tau^2} \frac{1}{l^2} v(\tau) &= 0, \\ \partial_\tau v + \frac{c}{\eta^{-2} + l^{-4}} \frac{1}{l^2} u(\tau) &= 0. \end{aligned} \tag{15}$$

Here we also shifted the time to $\tau \in [-l^{-2}, l^{-2}]$ for notational convenience.

In order to introduce ideas, we first consider the case where η is not larger than l^2 . There, we approximate $\eta^{-2} + \tau^2 \approx \eta^{-2} + l^{-4} \approx \eta^{-2}$ and obtain the following lemma:

Lemma 2. (Small η case) *Let $\eta > 0, l \in \mathbb{Z}$ and $c \in \mathbb{R}$ be given and consider the following approximation to system (15):*

$$\partial_\tau \begin{pmatrix} u \\ v \end{pmatrix} + \begin{pmatrix} 0 & -c\frac{\eta^2}{l^2} \\ c\frac{\eta^2}{l^2} & 0 \end{pmatrix} \begin{pmatrix} u \\ v \end{pmatrix} = 0, \tag{16}$$

for $\tau \in (-l^{-2}, l^{-2})$. Then, the unique solution is given by

$$\begin{pmatrix} u(\tau) \\ v(\tau) \end{pmatrix} = \begin{pmatrix} \cos\left(c\frac{\eta}{l^2}\tau\right) & \sin\left(c\frac{\eta}{l^2}\tau\right) \\ -\sin\left(c\frac{\eta}{l^2}\tau\right) & \cos\left(c\frac{\eta}{l^2}\tau\right) \end{pmatrix} \begin{pmatrix} u(0) \\ v(0) \end{pmatrix}. \tag{17}$$

In particular, it follows that $|u(\tau)|^2 + |v(\tau)|^2$ is a conserved quantity.

Proof. While the solution follows immediately by computation of the matrix exponential, for later reference we note an alternative proof arguing at the level of second derivatives. Formally differentiating the equation for u by τ and using the second equation, we obtain

$$\partial_\tau^2 u + c^2 \frac{\eta^4}{l^4} u = 0,$$

which has a general solution

$$\alpha \cos\left(c\frac{\eta}{l^2}\tau\right) + \beta \sin\left(c\frac{\eta}{l^2}\tau\right)$$

for constants $\alpha, \beta \in \mathbb{R}$. The first-order equation

$$\partial_\tau u - c\frac{\eta}{l^2}v = 0$$

then further determines $\partial_\tau u$ and allows us to determine α and β in terms of $u(0), v(0)$. \square

Having discussed the case of small η , in the following we are interested in the setting where $\eta \geq l^2$ is potentially very large. Here, we may again consider decoupled second order equations instead of a system of first order equations:

$$\begin{aligned} \partial_\tau^2 u + \frac{c^2}{\eta^{-2} + \tau^2} u &= 0, \\ \partial_\tau(\eta^{-2} + \tau^2)\partial_\tau v + c^2 v &= 0. \end{aligned} \tag{18}$$

We note that the equation for u is given by a stationary Schrödinger equation with potential $\frac{c^2}{\eta^{-2} + \tau^2}$, which is a mollified version of the scaling critical potential $\frac{c^2}{\tau^2}$.

The equation satisfied by v instead is a degenerate elliptic problem. However, we note that (15) allows us to determine v in terms of $\partial_\tau u$. Hence, in the following it suffices to study the evolution of u .

In Section 3 we showed that a Duhamel iteration converges, with the dominant terms being given by the nearest neighbor paths $(l, l - 1, l, l - 1, l, \dots)$ of this section, but had to require that c is sufficiently small to control logarithmic corrections. The aim in what follows is to study the system (15) and show that for large η better estimates hold and that only an absolute bound on c is required.

We recall that (18) is posed on the interval $(-l^{-2}, l^{-2})$. We may rescale our time variable as $\tau = l^{-2}t$, so that $t \in (-1, 1)$ and obtain

$$\partial_t^2 u + \frac{c^2}{\eta^{-2}l^4 + t^2} u = 0. \tag{19}$$

We stress that this equation depends only on c and $\frac{\eta}{l^2}$, but not on η and l separately. For simplicity of notation, in what follows we abbreviate

$$\xi := \frac{\eta}{l^2}. \tag{20}$$

Once we have computed u , we may recover v using that by (15)

$$\partial_t u = l^{-2} \partial_\tau u = \frac{c}{\eta^{-2} + \tau^2} \frac{1}{l^4} v. \tag{21}$$

We remark that for η much larger than l^2 at time $\tau = \frac{1}{l^4}$ it holds that $\frac{c}{\eta^{-2} + l^{-4}} l^{-4} \approx \frac{c}{l^{-4}} l^{-4} = c$. We further note that here v is evaluated at time $\tau = \frac{1}{l^2} t$, but corresponds to the mode $(l - 1, \eta)$ of the vorticity. The equation (19) has an explicit solution in terms of special functions, which we use to establish the following theorem:

Theorem 4. *Let $\xi = \frac{\eta}{l^2} \in \mathbb{R}$, then the problem (19)*

$$\partial_t^2 u + \frac{c^2}{\xi^{-2} + t^2} u = 0$$

on $(-1, 1)$ has an explicit scattering matrix $M = M(c, \xi)$ (see Proposition 3) such that

$$\begin{pmatrix} u(1) \\ \partial_t u(1) \end{pmatrix} = M \begin{pmatrix} u(-1) \\ \partial_t u(-1) \end{pmatrix}.$$

In particular, if $0 < c < \frac{1}{2}$ and $\xi \gg c^{-1}$, and we further assume that $\partial_t u(-1) \geq 0.5|u(-1)|$, then it follows that

$$\begin{pmatrix} u(1) \\ \partial_t u(1) \end{pmatrix} \approx \xi^\gamma u(-1) \begin{pmatrix} \pi c^2 \\ \pi c^2 \end{pmatrix}, \tag{22}$$

where $\gamma = \sqrt{1 - 4c^2}$.

We thus observe that the logarithmic correction seen in Section 3 here manifests in a modified exponent. As the explicit solution in terms of hypergeometric functions is technically involved but rather opaque, the proof of Theorem 4 is given in ‘‘Appendix A’’.

In what follows we instead discuss the underlying mechanism by approximating $\frac{c^2}{\xi^{-2}+t^2}$ by $\frac{c^2}{t^2}$ and $c^2\xi^2$ on sub-intervals of $(-1, 1)$. The interactions between these regimes then yields a correction to the exponent. We remark that in this section our splitting corresponds to a more rough approximation in order to introduce ideas. In Section 5.1 we choose our splitting more carefully and consider also the influence of other modes.

Let thus $\xi \in \mathbb{R}$ be given and suppose that $\xi \gg c^{-1}$ is large (otherwise we may consider the system (17)). Then on the intervals $(-1, \xi^{-1})$ and $(\xi^{-1}, 1)$, it holds that

$$\frac{1}{\xi^{-2} + t^2} \approx \frac{1}{t^2},$$

and we hence consider the second order ODEs

$$\partial_t^2 u + \frac{c^2}{t^2} u = 0 \tag{23}$$

on $\xi^{-1} < |t| < 1$ and

$$\partial_t^2 u + c^2 \xi^2 u = 0 \tag{24}$$

on $|t| < \xi^{-1}$.

In order to understand the mapping properties of Theorem 4 we thus consider the following scheme:

1. We prescribe initial data $(u(-1), \partial_t u(-1))$ and solve (23) to obtain $(u(-\xi^{-1}), \partial_t u(-\xi^{-1}))$ in Lemma 4
2. We then solve (24) to obtain $(u(\xi^{-1}), \partial_t u(\xi^{-1}))$ in Lemma 3
3. Finally, we solve (23) on $(\xi^{-1}, 1)$ by again using Lemma 4.
4. Combining these three maps, we show in Proposition 3 that the map $(u, \partial_t u)|_{t=-1} \mapsto (u, \partial_t u)|_{t=1}$ has singular values of size ξ^γ with $\gamma = \gamma(c) < 1$.

In the following Section 4.1 we then iterate this evolution in l to obtain sharp Gevrey 2-type norm inflation results for this model problem.

We remark that (23) is the scaling critical Schrödinger problem, which has a general solution

$$c_1 |t|^{\frac{1}{2}(1+\sqrt{1-4c^2})} + c_2 |t|^{\frac{1}{2}(1-\sqrt{1-4c^2})},$$

provided $0 < c < \frac{1}{2}$. We in particular note that the exponents $\gamma_1 = \frac{1}{2} + \sqrt{\frac{1}{4} - c^2}$ and $\gamma_2 = \frac{1}{2} - \sqrt{\frac{1}{4} - c^2}$ are strictly between 0 and 1 and $\gamma_1 + \gamma_2 = 1$.

Lemma 3. *Let $c \neq 0, \eta > 0, k \in \mathbb{Z} \setminus \{0\}$, then the solution $u \in C^2$ of*

$$\partial_t^2 u + c^2 \xi^2 u = 0 \tag{25}$$

in $(-\xi^{-1}, \xi^{-1})$ satisfies

$$\begin{pmatrix} u(\xi^{-1}) \\ \partial_t u(\xi^{-1}) \end{pmatrix} = \begin{pmatrix} \cos(2c) & \frac{1}{c} \xi^{-1} \sin(2c) \\ -c \xi \sin(2c) & \cos(2c) \end{pmatrix} \begin{pmatrix} u(-\xi^{-1}) \\ \partial_t u(-\xi^{-1}) \end{pmatrix}.$$

In particular, we observe that for $\xi \gg c^{-2}$ the bottom left matrix entry is by far the largest, while $\cos(2c) \approx 1$.

Proof of Lemma 3. We observe that a general solution of (25) is given by $c_1 \sin(c\xi t) + c_2 \cos(c\xi t)$. One may then verify that

$$\begin{pmatrix} u(t) \\ \partial_t u(t) \end{pmatrix} = \begin{pmatrix} \cos(c\xi(t + \xi^{-1})) & \frac{1}{c} \xi^{-1} \sin(c\xi(t + \xi^{-1})) \\ -c \xi \sin(c\xi(t + \xi^{-1})) & \cos(c\xi(t + \xi^{-1})) \end{pmatrix} \begin{pmatrix} u(-\xi^{-1}) \\ \partial_t u(-\xi^{-1}) \end{pmatrix}$$

satisfies the initial conditions. \square

On the exterior intervals we similarly obtain an explicit solution, but with very different dependence on ξ .

Lemma 4. *Let $0 < c < \frac{1}{2}$ and $\xi > 1$ then the solution of $u \in C^2$ of*

$$\partial_t^2 u + \frac{c^2}{t^2} u = 0 \text{ in } (-1, 1) \setminus (-\xi^{-1}, \xi^{-1}) \tag{26}$$

satisfies

$$\begin{pmatrix} u(\xi^{-1}) \\ \partial_t u(\xi^{-1}) \end{pmatrix} = \begin{pmatrix} \xi^{-\gamma_1} & \xi^{-\gamma_2} \\ \gamma_1 \xi^{1-\gamma_1} & \gamma_2 \xi^{1-\gamma_2} \end{pmatrix} \begin{pmatrix} 1 & 1 \\ \gamma_1 & \gamma_2 \end{pmatrix}^{-1} \begin{pmatrix} u(1) \\ \partial_t u(1) \end{pmatrix},$$

and

$$\begin{pmatrix} u(-\xi^{-1}) \\ \partial_t u(-\xi^{-1}) \end{pmatrix} = \begin{pmatrix} \xi^{-\gamma_1} & \xi^{-\gamma_2} \\ -\gamma_1 \xi^{1-\gamma_1} & -\gamma_2 \xi^{1-\gamma_2} \end{pmatrix} \begin{pmatrix} 1 & 1 \\ -\gamma_1 & -\gamma_2 \end{pmatrix}^{-1} \begin{pmatrix} u(-1) \\ \partial_t u(-1) \end{pmatrix},$$

where $\gamma_{1,2} = \frac{1}{2} \pm \sqrt{\frac{1}{4} - c^2}$.

Proof of Lemma 4. We note that a general solution of (26) is given by

$$c_{1,\pm} |t|^{\gamma_1} + c_{2,\pm} |t|^{\gamma_2},$$

where $c_{,\pm}$ are constant on each connected component of the domain. We may then again verify that

$$\begin{pmatrix} u(t) \\ \partial_t u(t) \end{pmatrix} = \begin{pmatrix} |t|^{\gamma_1} & |t|^{\gamma_2} \\ \text{sgn}(t) \gamma_1 |t|^{\gamma_1-1} & \text{sgn}(t) \gamma_2 |t|^{\gamma_2-1} \end{pmatrix} \begin{pmatrix} 1 & 1 \\ \text{sgn}(t) \gamma_1 & \text{sgn}(t) \gamma_2 \end{pmatrix}^{-1} \begin{pmatrix} u(1) \\ \partial_t u(1) \end{pmatrix}$$

satisfies the boundary conditions. \square

Proposition 3. Let $0 < c < \frac{1}{2}$ and let $u \in C^1$ be a solution of (23) and (24). Then there exists an explicitly computable matrix $M = M(c, \xi)$ such that

$$\begin{pmatrix} u(1) \\ \partial_t u(1) \end{pmatrix} = M \begin{pmatrix} u(-1) \\ \partial_t u(-1) \end{pmatrix}. \tag{27}$$

In particular, if $0 < c < \frac{1}{2}$ and $\xi \gg c^{-1}$, and we further assume that $|u(-1)| \geq 2|u'(-1)|$, then it follows that

$$\begin{pmatrix} u(1) \\ \partial_t u(1) \end{pmatrix} \approx \xi^\gamma u(-1) \begin{pmatrix} 8c^2 \\ 8c^2 \end{pmatrix}, \tag{28}$$

where $\gamma = \sqrt{1 - 4c^2}$.

Proof of Proposition 3. Combining Lemma 3 and Lemma 4, we obtain that our data at $t = -1$ and $t = 1$ are related by

$$\begin{pmatrix} 1 & 1 \\ \gamma_1 & \gamma_2 \end{pmatrix} \begin{pmatrix} \xi^{-\gamma_1} & \xi^{-\gamma_2} \\ \gamma_1 \xi^{1-\gamma_1} & \gamma_2 \xi^{1-\gamma_2} \end{pmatrix}^{-1} \begin{pmatrix} \cos(2c) & \frac{1}{c} \xi^{-1} \sin(2c) \\ -c \xi \sin(2c) & \cos(2c) \end{pmatrix} \\ \begin{pmatrix} \xi^{-\gamma_1} & \xi^{-\gamma_2} \\ -\gamma_1 \xi^{1-\gamma_1} & -\gamma_2 \xi^{1-\gamma_2} \end{pmatrix} \begin{pmatrix} 1 & 1 \\ -\gamma_1 & -\gamma_2 \end{pmatrix}^{-1}.$$

We now compute

$$\begin{pmatrix} \xi^{-\gamma_1} & \xi^{-\gamma_2} \\ \gamma_1 \xi^{1-\gamma_1} & \gamma_2 \xi^{1-\gamma_2} \end{pmatrix}^{-1} = -\frac{1}{\gamma} \begin{pmatrix} \gamma_2 \xi^{1-\gamma_2} & -\xi^{-\gamma_2} \\ -\gamma_1 \xi^{1-\gamma_1} & \xi^{-\gamma_1} \end{pmatrix}, \\ \begin{pmatrix} 1 & 1 \\ -\gamma_1 & -\gamma_2 \end{pmatrix}^{-1} = \frac{1}{\gamma} \begin{pmatrix} -\gamma_2 & -1 \\ \gamma_1 & 1 \end{pmatrix},$$

where we used $\gamma_1 + \gamma_2 = 1$ and computed the determinants as

$$\gamma_1 - \gamma_2 =: \gamma = \sqrt{1 - 4c^2} \approx 1.$$

Since $\xi \gg 1$, $\gamma_1 = \frac{1}{2} + \sqrt{\frac{1}{4} - c^2} \approx 1 - c^2/2 \approx 1$, $\gamma_2 = \frac{1}{2} - \sqrt{\frac{1}{4} - c^2} \approx c^2 \ll 1$ it follows that the largest powers of ξ are dominant. We may thus approximate:

$$\begin{pmatrix} u(-\xi^{-1}) \\ \partial_t u(-\xi^{-1}) \end{pmatrix} = \frac{1}{\gamma} \begin{pmatrix} \xi^{-\gamma_1} & \xi^{-\gamma_2} \\ -\gamma_1 \xi^{1-\gamma_1} & -\gamma_2 \xi^{1-\gamma_2} \end{pmatrix} \begin{pmatrix} -\gamma_2 & -1 \\ \gamma_1 & 1 \end{pmatrix} \begin{pmatrix} u(-1) \\ u'(-1) \end{pmatrix} \\ \approx \begin{pmatrix} \xi^{-\gamma_1} & \xi^{-\gamma_2} \\ -\gamma_1 \xi^{1-\gamma_1} & -\gamma_2 \xi^{1-\gamma_2} \end{pmatrix} \begin{pmatrix} -\gamma_2 u(-1) - \partial_t u(-1) \\ u(-1) + \partial_t u(-1) \end{pmatrix} \\ \approx u(-1) \begin{pmatrix} \xi^{-\gamma_2} \\ -\gamma_2 \xi^{1-\gamma_2} \end{pmatrix}.$$

Next, it follows that

$$\begin{pmatrix} u(\xi^{-1}) \\ \partial_t u(\xi^{-1}) \end{pmatrix} = \begin{pmatrix} \cos(2c) & \frac{1}{c} \xi^{-1} \sin(2c) \\ -c \xi \sin(2c) & \cos(2c) \end{pmatrix} \begin{pmatrix} u(-\xi^{-1}) \\ \partial_t u(-\xi^{-1}) \end{pmatrix} \\ \approx u(-1) \begin{pmatrix} \xi^{-\gamma_2} \\ -(2c^2 + \gamma_2) \xi^{1-\gamma_2} \end{pmatrix},$$

where we omitted $\gamma_2 \frac{\sin(2c)}{c} = \mathcal{O}(c^2)$ as a small perturbation to 1 and approximated $\sin(2c) \approx 2c$. Finally, it follows that

$$\begin{aligned} \begin{pmatrix} u(1) \\ \partial_t u(1) \end{pmatrix} &= -\frac{1}{\gamma} \begin{pmatrix} 1 & 1 \\ \gamma_1 & \gamma_2 \end{pmatrix} \begin{pmatrix} \gamma_2 \xi^{1-\gamma_2} & -\xi^{-\gamma_2} \\ -\gamma_1 \xi^{1-\gamma_1} & \xi^{-\gamma_1} \end{pmatrix} \begin{pmatrix} u(\xi^{-1}) \\ u'(\xi^{-1}) \end{pmatrix} \\ &\approx -u(-1) \begin{pmatrix} 1 & 1 \\ \gamma_1 & \gamma_2 \end{pmatrix} \begin{pmatrix} (2\gamma^2 + 2c^2)\xi^{1-2\gamma_2} \\ -1\xi^0 \end{pmatrix} \\ &\approx u(-1)\xi^{1-2\gamma_2} \begin{pmatrix} 8c^2 \\ 8c^2 \end{pmatrix}. \end{aligned}$$

□

We thus see that, while the evolution on the small interval $(-\xi^{-1}, \xi^{-1})$ yields a singular value of size ξ^1 , the conjugation with the power law evolution on $|t| > \xi^{-1}$ yields a much smaller singular value $\xi^\gamma \ll \xi^1$.

Having establish a precise description of the evolution for times close to a single resonant time $\tau \approx \frac{1}{l}$, in the following we consider an iterated model to study the norm inflation in Gevrey spaces and the associated asymptotic behavior.

4.1. Model Echo Chains and Modified Exponents

In Section 3 we have studied chains of echoes for the linear problem (4) and have established norm inflation with a factor $\exp(C\sqrt{\eta})$. However, in that case our proof limited us to considering only η such that $c \ln(1 + \eta^2)$ is not too large.

In what follows we instead consider an iterated version of the model of Section 4, which does not possess such an obstruction. In particular, combining the behavior of infinitely many modes η_j with $\eta_j \rightarrow \infty$ we construct solutions which exhibit norm inflation for arbitrarily large times and do not converge as time tends to infinity. However, despite the failure of the convergence of the vorticity, the velocity field is shown to converge.

We briefly recall the approximations made in the preceding sections of this article. We started with the 2D Euler equations close to Couette flow

$$\partial_t \omega + y \partial_x \omega + v \cdot \nabla \omega = 0,$$

and focused on perturbations of the form $\omega = c \cos(x + ty) + \varepsilon \omega_*$. Omitting the transport by $c \frac{\sin(x+ty)}{1+t^2}$ and changing to variables $(x + ty, y)$, we obtain

$$\partial_t \omega + c \sin(x) \partial_y \Delta_t^{-1} \omega + \varepsilon \nabla^\perp \Delta_t^{-1} \omega \cdot \nabla \omega = 0.$$

In particular, we note that this equation formally conserves $\|c \cos(x) + \varepsilon \omega\|_{L^2}$. Omitting the nonlinearity by letting $\varepsilon \downarrow 0$, we lose this conserved quantity, but obtain an explicit Fourier problem with nearest-neighbor-interactions:

$$\begin{aligned} \partial_t \tilde{\omega}(t, k, \eta) + c \frac{\eta}{(k-1)^2 + (\eta - (k-1)t)^2} \tilde{\omega}(t, k-1, \eta) \\ - c \frac{\eta}{(k+1)^2 + (\eta - (k+1)t)^2} \tilde{\omega}(t, k+1, \eta) = 0. \end{aligned} \tag{29}$$

Here with slight abuse of notation we replaced c by $c/2$ for brevity. This further highlights resonant times, where $\eta - (k \pm 1)t \approx 0 \Leftrightarrow t \approx \frac{\eta}{k \pm 1}$. After studying the system (29) in Sections 2 and 3, in Section 4 we further introduced a model problem (15) that focuses solely on resonant modes v ((l, η) such that $\eta - lt \approx 0$) and their neighbors u :

$$\begin{aligned} \partial_\tau u - \frac{c}{\eta^{-2} + \tau^2} \frac{1}{l^2} v(\tau) &= 0, \\ \partial_\tau v + \frac{c}{\eta^{-2} + l^{-4}} \frac{1}{l^2} u(\tau) &= 0. \end{aligned} \tag{30}$$

In particular, we showed that $v|_{\tau=\tau_1}$ is approximately of size $c(\frac{\eta}{2})^\gamma u|_{\tau=\tau_0}$.

Building on the single resonance results of Theorem 4, we construct the following iteration scheme:

- Let $k \in \mathbb{N}, k > 1$ and $\eta \in \mathbb{R}$ be given and define $\tau_k = \frac{1}{2}(\frac{1}{k-1} + \frac{1}{k})$. We then prescribe $(u, \partial_\tau u)|_{\tau_k}$ and use equation (19) to determine $(u, \partial_t u)|_{\tau_{k-1}}$.
- Relabeling v of the previous step as u of the case $k - 1$ and using (21) we prescribe

$$\begin{pmatrix} u_{k-1} \\ \partial_t u_{k-1} \end{pmatrix} |_{\tau=\tau_{k-1}} := \begin{pmatrix} c^{-1} \partial_t u_k \\ cu_k \end{pmatrix} |_{\tau=\tau_{k-1}}.$$

We then again use equation (19) to determine $(u, \partial_t u)|_{\tau_{k-2}}$.

- We iterate this procedure until we reach τ_1 , where we define $\tau_0 = 1.5$.

Recalling the construction of the model problem of Section 4, $(u, \partial_t u)|_{\tau=\tau_k}$ corresponds to prescribing $(\tilde{\omega}(t, k, \eta), \tilde{\omega}(t, k - 1, \eta))$ at time $t = \frac{1}{2}(\frac{\eta}{k-1} + \frac{\eta}{k})$ and $(u, u')|_{\tau=\tau_0}$ corresponds to the value of the modes $(1, \eta)$ and $(0, \eta)$ at time $t = 1.5\eta$.

Theorem 5. *Let $k \in \mathbb{N}$ and $\eta \in \mathbb{R}$ be given and prescribe $(u, \partial_t u)|_{\tau=\tau_k} = (1, 0)$. Then, the above iteration scheme yields that*

$$(u, \partial_t u)|_{\tau=\tau_0} \approx (c, c^2)c^{k-1} \left(\frac{\eta^k}{(k!)^2} \right)^\gamma,$$

where $\gamma = \sqrt{1 - 4c^2} \neq 1$. In particular, choosing k maximally for c, η fixed, we obtain a growth factor

$$\max_k c^k \left(\frac{\eta^k}{(k!)^2} \right)^\gamma \sim e^{2\gamma\sqrt{c^{1/\gamma}\eta}}$$

consistent with a Gevrey regularity class.

Proof. Using the result of Theorem 4 (specifically equation (22)), we observe that

$$\begin{pmatrix} u_{k-1} \\ \partial_t u_{k-1} \end{pmatrix} |_{\tau=\tau_{k-1}} \approx \partial_t u_k(\tau_k) \xi^\gamma \begin{pmatrix} 8c \\ 8c^3 \end{pmatrix}.$$

In particular, we note that $u_{k-1} \approx c^{-2} \partial_t u_{k-1}$ and we may thus apply Theorem 4 again. We repeat this process another $k - 1$ times, where $\xi = \frac{\eta}{k^2}$ changes in each step, and thus obtain the claimed growth factor

$$c^k \left(\frac{\eta^k}{(k!)^2} \right)^\gamma.$$

Considering $c\eta^\gamma$ large and k large, by Stirling’s approximation it holds that

$$c^k \left(\frac{\eta^k}{(k!)^2} \right)^\gamma \sim c^k \left(\eta^k \frac{e^{2k}}{2\pi k k^{2k}} \right)^\gamma.$$

Choosing k as approximately $\sqrt{c^{\frac{1}{\gamma}} \eta}$ (rounding up or down), we obtain a cancellation of $c^k \left(\frac{\eta^k}{k^{2k}} \right)^\gamma = 1$ and thus

$$e^{2\gamma \sqrt{c^{1/\gamma} \eta}} \left(\frac{1}{2\pi \sqrt{c^{1/\gamma} \eta}} \right)^\gamma$$

as the maximal growth factor. \square

As a corollary, for this model we can construct initial data in a critical Gevrey regularity class.

Theorem 6. *Consider the chained two mode model with $0 < c < \frac{1}{2}$ and $\gamma = \sqrt{1 - 4c^2}$. Then there exists $s = s(\gamma)$ and $C > 0$ such that for every $\varepsilon > 0$ and every $\sigma_0 \in \mathbb{R}$, there exists initial data $u_0 \in \mathcal{G}_{C, \frac{1}{s}}$ (see Section 1.1.1 for a definition) such that*

$$\|u_0\|_{\mathcal{G}_{C, \frac{1}{s}}} < \varepsilon, \tag{31}$$

and such that for every $\tilde{C} > 0$,

$$\lim_{t \rightarrow \infty} \|u(t)\|_{\mathcal{G}_{s, \tilde{C}}} = \infty.$$

Furthermore, $u(t)$ does converge in H^σ , $\sigma < \sigma_0$, but diverges in H^σ , $\sigma > \sigma_0$.

In particular, choosing $-1 < \sigma_0 < 0$, we find initial data, arbitrarily small in the critical Gevrey regularity class, such that the vorticity does not converge in L^2 as $t \rightarrow \infty$, but the velocity field does converge.

Proof. For any given η , let k_η be the associated maximizer of the growth factor obtained in Theorem 5 and let $g(\eta)$ denote that growth factor.

Let now $\psi \in \cap_{\sigma < \sigma_0} H^\sigma(\mathbb{R})$ be given and prescribe as initial data

$$u = \int_\eta \frac{1}{g(\eta)} \hat{\psi}(\eta) e^{i\eta y + ik_\eta x}.$$

Since $g(\eta) \approx \exp(C\sqrt{\eta})$ for $C = C(c)$, this function is in a Gevrey class.

Then by construction, the $k = 1$ mode will asymptotically be given by

$$\int_{\eta} g(\eta) \frac{1}{g(\eta)} \hat{\psi}(\eta) e^{i\eta y + ix} = \psi(y) e^{ix}$$

□

Building on the insights obtained in this model problem in the next Section we consider the full problem.

5. Echo Chains as a Linear Mechanism and Modified Scattering

In the previous section we have shown that for large η a linear growth factor $\frac{\eta}{i^2}$ cannot be expected to be accurate anymore. Indeed, the logarithmic corrections in the Duhamel iteration are much larger than the prior “leading term”. Instead we expect to see a modified exponent, which is less than 1 due to cancellations with neighboring modes. In what follows we consider the full model (7)

$$\begin{aligned} \partial_{\tau} \omega(\tau, k, \eta) + c \frac{1}{(k-1)^2} \frac{1}{\eta^{-2} + (\frac{1}{k-1} - \tau)^2} \omega(\tau, k-1, \eta) \\ - c \frac{1}{(k+1)^2} \frac{1}{\eta^{-2} + (\frac{1}{k+1} - \tau)^2} \omega(\tau, k+1, \eta) = 0 \end{aligned} \tag{32}$$

on a time-interval around a single resonance $\frac{1}{k_0}$,

$$\left(\frac{\frac{1}{k_0+1} + \frac{1}{k_0}}{2}, \frac{\frac{1}{k_0} + \frac{1}{k_0-1}}{2} \right).$$

The main aim of this section, which is achieved in Theorem 8 of Section 5.2, is then to show that the evolution of this full model is largely determined by the evolution of the three modes $k_0 + 1, k_0, k_0 - 1$. To this end, in Section 5.1 we first introduce a model problem focusing on just these three modes, which should be seen as a refinement of the two mode model of Section 4. Here, as in Section 4 a key challenge lies in capturing the asymptotic dependence on $\xi = \frac{\eta}{k^2}$ for ξ very large. Subsequently, we employ a bootstrap approach in the (sub)subsections of Section 5.2 to show that this three mode model indeed captures the full growth mechanisms and that all other modes can be controlled. Taking these results as building blocks, in Section 5.3 we combine the resonances to construct echo chains which exhibit Gevrey 2 growth and furthermore combine countably many echo chains to construct solutions with critical stability/blow-up behavior.

Considering that

$$\frac{1}{k_0 \pm 1} - \frac{1}{k_0} = \frac{1}{k_0(k_0 \pm 1)} \leq \frac{2}{k_0^2},$$

and the central role of $\tau = \frac{1}{k_0}$, compared to Section 4 we again change variables as

$$t = k_0^2 \left(\tau - \frac{1}{k_0} \right) \in \left(-\frac{k_0}{2(k_0 + 1)}, \frac{k_0}{2(k_0 - 1)} \right) =: (t_0, t_1). \quad (33)$$

Then $\partial_\tau = k_0^2 \partial_t$ and hence (7) reads as

$$\partial_t \omega(t, k, \eta) + a(k - 1) \omega(t, k - 1, \eta) - a(k + 1) \omega(t, k + 1, \eta) = 0, \quad (34)$$

where

$$a(k_0) = \frac{c}{\left(\frac{\eta}{k_0^2} \right)^{-2} + t^2}, \quad (35)$$

and we abbreviate

$$\xi := \frac{\eta}{k_0^2}, \quad (36)$$

and for $k \neq k_0$,

$$a(k) = \frac{1}{k^2 k_0^2} \frac{c}{\eta^{-2} + \left(\frac{1}{k} - \frac{1}{k_0} - k_0^{-2} t \right)^2} \quad (37)$$

$$= \frac{c}{\eta^{-2} k_0^2 k^2 + \left(k_0 - k - \frac{k}{k_0} t \right)^2}. \quad (38)$$

In particular, we note that since we only consider the resonant interval around k_0 and are thus far from other resonant times $|a(k)| \leq 4c$ for any $k \neq k_0$. In contrast, $a(k_0)$ at time $t = 0$ is of size $c\xi^2 \gg 1$.

We remark that in Section 4 we made several simplifications compared to the full model:

1. We approximated $a(k_0 \pm 1) \approx \pm c$, which allowed us to compute explicit solutions. In the following we need to show that this is a valid approximation, that is the evolution of the full problem can be estimated above and below by the approximate evolution.
2. In Section 4 we considered a two-mode model involving just k_0 and $k_0 - 1$. Instead we show that the precise behavior is more accurately captured by the three-mode model involving $k_0 - 1, k_0, k_0 + 1$, which yields a change of the exponent γ (involving $2c^2$ in place of c^2).
3. In view of the sizes of $a(k)$ in the following Section 5.1 we at first again neglect all except the three modes $k_0 - 1, k_0, k_0 + 1$. We call this the *three-mode model*. In contrast to the two-mode model of Section 4 we here do not approximate the coefficient functions. In Section 5.2 we then discuss the full problem incorporating all modes and prove that indeed all other modes can be treated as perturbations in a bootstrap approach.

5.1. The Three-Mode Model

In this section we introduce the solution operator of the homogeneous three-mode model:

$$\partial_t u(k) + a(k + 1)u(k + 1) - a(k - 1)u(k - 1) = 0, \tag{39}$$

$$a(k) = \begin{cases} a(k_0 \pm 1) = \frac{c}{\xi^{-2}(\frac{k_0 \pm 1}{k_0})^2 + (1 \pm \frac{k_0 \pm 1}{k_0} t)^2}, \\ a(k_0) = \frac{c}{\xi^{-2} + t^2}, \\ 0. \end{cases} \tag{40}$$

That is, we consider only the modes $k_0 - 1, k_0, k_0 + 1$ and neglect all other modes as “inhomogeneities”. This allows for a clearer discussion of the growth and decay mechanisms and serves to introduce the techniques of proof used in the different regimes. It is also an important preliminary step before studying the behavior of the full problem in Theorem 8 of Section 5.2. More precisely, we will (with considerable technical effort) consider the full problem as a perturbation by an inhomogeneity, which we control by means of a bootstrap approach.

We recall that, by the results of Section 4, the heuristic of the approximate model $a(k_0 \pm 1) \approx c$ suggests a power law behavior of solutions. As for the present case of exact coefficients an explicit solution is not feasible anymore, so we establish a comparison estimate. We argue in multiple steps:

- By symmetry it holds that $u(k_0 + 1) + u(k_0 - 1) = \text{const.}$ We hence to some extent reduce to a two-mode model. However, as $a(k_0 + 1) \neq a(k_0 - 1)$ in this model the problem does not completely decouple.
- We first establish a rough power law upper bound on the growth of solutions on $(-t_0, -t)$ as $t \downarrow 0$. This is achieved by concatenating multiple small time estimates.
- Subsequently, we iteratively improve this bound to $t^{\gamma_2 - 1}$, $\gamma_2 = \frac{1}{2} - \sqrt{\frac{1}{4} - 2c^2}$ upper and lower bounds similar as in Section 4 (here in the formula for γ_2 the c^2 is replaced by $2c^2$ due to the third mode). This step relies on reformulations of the ODE system as second order ODEs and integrating these.
- We then show that the resonant mode is decreasing as t^{γ_2} while the neighboring modes grow like $t^{\gamma_2 - 1}$ both with an upper and lower bound.
- Combining these results, we construct solution operators on the intervals $I_1 = (t_0, -\frac{d}{\xi})$, $I_3 = (\frac{d}{\xi}, t_1)$ with $d = c^{-1}$. On the resonant interval $I_2 = (-\frac{d}{\xi}, \frac{d}{\xi})$ we instead use a Duhamel iteration argument to construct the solution operator.
- Concatenating the solution operators we obtain the solution operator from time t_0 to t_1 and show that it exhibits analogous $(\frac{\eta}{k_0^2})^\nu$ growth behavior to the model problem of Section 4. In Section 5.2 we then show that this behavior persists in the full problem.

The main results of this section are summarized in the following theorem:

Theorem 7. (The full solution, homogeneous case) *Let $\xi = \frac{\eta}{k_0^2} \gg 1$ and $0 < c < 0.2$ be given and let $t_0 = -\frac{1}{2} \frac{k_0}{k_0 + 1}$ and $t_1 = \frac{1}{2} \frac{k_0}{k_0 - 1}$. We consider the ODE system*

(39) for $u_1 = \frac{u(k_0+1)-u(k_0-1)}{2}$, $u_2 = u(k_0)$, $u_3 = \frac{u(k_0+1)+u(k_0-1)}{2}$:

$$\partial_t u + M(t)u = 0, \tag{41}$$

where

$$M(t) = \begin{pmatrix} 0 & a(k_0) & a(k_0 + 1) - a(k_0 - 1) \\ -(a(k_0 + 1) - a(k_0 - 1)) & 0 & -a(k_0 + 1) + a(k_0 - 1) \\ 0 & 0 & 0 \end{pmatrix}$$

on the interval (t_0, t_1) .

Suppose that at time $t_0 = -\frac{1}{2} \frac{k_0}{k_0+1}$ it holds that $u(k_0) \geq 0.5 \max(u)$. Then at time $t_1 = \frac{1}{2} \frac{k_0}{k_0-1}$ it holds that

$$\begin{aligned} u_1(t_1) &\approx c^{2-2\gamma_2} \xi^\gamma u_2(t_0), \\ u_2(t_1) &\approx c^{4-2\gamma_2} \xi^\gamma u_2(t_0), \\ u_3(t_1) &= u_3(t_0), \end{aligned}$$

where $\gamma_2 = \frac{1}{2} - \sqrt{\frac{1}{4} - 2c^2}$ and $\gamma = 1 - 2\gamma_2 = \sqrt{1 - 4c^2}$.

We note, in particular, that the exponent here is different from 1 (which was not visible in prior works due to the logarithmic constraints) and that at time t_1 our solution u satisfies the assumptions of this theorem with k_0 replaced by $k_0 - 1$. Thus, we may iteratively apply this theorem until we reach the frequency 1 and obtain the following corollary (in Section 5.3 we return to this echo chain behavior in the context of the full problem and also discuss (modified) asymptotic behavior as $t \rightarrow \infty$):

Corollary 1. (Echo chain, homogeneous case) *Let $\eta \gg 1$ and $k_0 \ll \eta$ be given. We then consider the iterated system where initially set $\xi = \frac{\eta}{k_0^2}$ and solve (41) with*

$$u(t_0) = \begin{pmatrix} 0 \\ 1 \\ 0 \end{pmatrix}$$

In the next step, we then decrease $k_0 \mapsto k_0 - 1$ and set $\xi = \frac{\eta}{(k_0-1)^2}$ and again solve (41) with new initial data being given by

$$\begin{pmatrix} u_1(t_0) \\ u_2(t_0) \\ u_3(t_0) \end{pmatrix}_{new} = \begin{pmatrix} 0 \\ u_1(t_1) \\ 0 \end{pmatrix}_{old}$$

We iterate this procedure another $k_0 - 2$ times until we have reached $k = 1$. At this time it then holds that

$$u_1(1) \approx (c^{2-2\gamma_2})^{k_0} \left(\frac{\eta^{k_0}}{(k_0!)^2} \right)^\gamma.$$

In particular, choosing k_0 maximally for given η we obtain that $u_1(1) \approx C_1 \exp(C_2 \sqrt{\eta})$ attains a Gevrey 2 norm inflation factor:

We remark that the precise choice of u_1, u_3 in the update step here does not change the result as long as $u_2(t_0)$ is (comparable to) the largest entry and thus Theorem 7 can be applied.

Proof of Corollary 1. Let k_0 be given and let $t_0[k_0] = \frac{1}{2}(\frac{1}{k_0+1} - \frac{1}{k_0})$, $t_1[k_0] = \frac{1}{2}(\frac{1}{k_0-1} - \frac{1}{k_0})$. Then at time $t_0[k_0]$ it holds that $u_2 \geq 0.5 \max(u)$ and we may thus apply Theorem 7 to conclude that at time $t_1[k_0]$,

$$u_1(t_1[k_0]) \approx c^{2-2\gamma_2} \xi^\gamma u_2(t_0[k_0]) \geq 0.5 \max(u(t_1[k_0])).$$

We then decrease k_0 by 1 and obtain that, at the new initial time $t_0[k_0 - 1] = t_1[k_0]$, by relabeling the above,

$$u_2(t_0[k_0 - 1]) = u_1(t_1[k_0]) \geq 0.5 \max(u)$$

again satisfies the assumptions of Theorem 7. Iterating this procedure until we reach $k_0 = 1$ then yields the result. \square

Our proof proceeds by considering the three intervals $I_1 = (t_0, -\frac{d}{\xi})$, $I_2 = (-\frac{d}{\xi}, \frac{d}{\xi})$, $I_3 = (\frac{d}{\xi}, t_1)$, where $d = c^{-1}$. The right boundary datum of each interval then serves as the left boundary datum of the next. We thus follow a similar strategy as in the proof of Theorem 4. In the Section 5.1.1 we derive upper and lower bounds for the evolution on the intervals I_1 and I_3 , which we expect to be similar to the ones derived for the two mode model in Lemma 4. That is, we expect to derive asymptotics in terms of $\xi^{-\gamma_1}$, $\xi^{-\gamma_2}$ and $\xi^{1-\gamma_1}$, $\xi^{1-\gamma_2}$. In Section 5.1.2 we study the evolution on the inner interval I_2 , where we expect an evolution similar to Lemma 4, and thus to encounter powers ξ^{+1} and ξ^{-1} . Combining these asymptotics as in Proposition 3, we then obtain a combined exponent γ , which concludes our proof of Theorem 7.

5.1.1. The Interval I_1 As a first step in our proof of Theorem 7 we consider the evolution of the evolution of the homogeneous problem (39) on the interval

$$I_1 = \left(t_0, -\frac{d}{\xi} \right).$$

Similarly to the results of Lemma 3 of Section 4 we here aim to show that for given initial data $(u_1, u_2, u_3)|_{t=t_0}$ at the right end of the interval u_1 has grown and u_2 has decreased by an appropriate power of ξ . Our results are summarized in the following proposition:

Proposition 4. (Left interval, homogeneous case) *Let $\xi = \frac{\eta}{k} \gg 1$ and $0 < c < 0.2$ be given.*

We consider the problem (39) on the interval $I_1 = (t_0, -\frac{d}{\xi})$ with $d = c^{-1}$. Then the unique solution $u(t)$ satisfies

$$u_3(t) = u_3(t_0).$$

Furthermore, it holds that

$$\begin{aligned} |u_1(t)| &\leq Ct^{\gamma_2-1} \max(u(t_0)), \\ |u_2(t)| &\leq Ct^{\gamma_2} \max(u(t_0)). \end{aligned}$$

where $\gamma_1 = \frac{1}{2} + \sqrt{\frac{1}{4} - 2c^2}$, $\gamma_2 = \frac{1}{2} - \sqrt{\frac{1}{4} - 2c^2}$.

If in addition $|u_2(t_0)| \geq 4c|u(t_0)|$, then at time $t = -\frac{d}{\xi}$ it holds that

$$\begin{aligned} u_1\left(-\frac{d}{\xi}\right) &\approx \left(\frac{d}{\xi}\right)^{\gamma_2-1} \frac{c}{\gamma_2 - \gamma_1} u_2(t_0), \\ u_2\left(-\frac{d}{\xi}\right) &\approx \left(\frac{d}{\xi}\right)^{\gamma_2} \frac{\gamma_1}{\gamma_2 - \gamma_1} u_2(t_0). \end{aligned}$$

For later reference we note that $u_1(-\frac{d}{\xi}) \gg u_2(-\frac{d}{\xi})$ and that $u_1(-\frac{d}{\xi})(\frac{d}{\xi}) \approx c u_2(-\frac{d}{\xi})$.

We may write the homogeneous three-mode model (39) as

$$\partial_t \begin{pmatrix} u(k_0 - 1) \\ u(k_0) \\ u(k_0 + 1) \end{pmatrix} + \begin{pmatrix} 0 & a & 0 \\ -b_1 & 0 & b_2 \\ 0 & -a & 0 \end{pmatrix} \begin{pmatrix} u(k_0 - 1) \\ u(k_0) \\ u(k_0 + 1) \end{pmatrix} = 0, \tag{42}$$

where $a = \frac{c}{\xi^{-2} + t^2}$, $b_1 \neq b_2 \approx c$. For this structure it is apparent that $u_3 := \frac{1}{2}(u(k_0 - 1) + u(k_0 + 1))$ is conserved. However, compared to the two mode model studied in Lemma 3 here u_3 serves as an inhomogeneity for the remaining modes, as shown in the following lemma. This inhomogeneity then leads to a particular solution constructed in Lemma 7. Furthermore, unlike in the model setting of Section 4 the resulting problem does not allow for explicit solutions. Instead we show by successive improvement in Lemma 6 that the solutions can be bounded by an appropriate power law.

Lemma 5. (Reduction 1) Consider the problem (42), then it holds that

$$u(k - 1) + u(k + 1) = \text{const.}$$

We hence introduce the change of unknowns and notation

$$\begin{aligned} u_1 &= \frac{1}{2}(u(k_0 - 1) - u(k_0 + 1)), \\ u_2 &= u(k_0), \\ u_3 &= \frac{1}{2}(u(k_0 - 1) + u(k_0 + 1)), \end{aligned}$$

where u_3 is invariant under the evolution and (u_1, u_2) solve

$$\partial_t \begin{pmatrix} u_1 \\ u_2 \end{pmatrix} + \begin{pmatrix} 0 & a \\ -b & 0 \end{pmatrix} \begin{pmatrix} u_1 \\ u_2 \end{pmatrix} = \begin{pmatrix} (b_1 + b_2)u_3 \\ (b_1 - b_2)u_3 \end{pmatrix}, \tag{43}$$

where $b = b_1 + b_2$.

Proof of Lemma 5. We observe that

$$\partial_t(u(k_0 + 1) + u(k_0 - 1)) = -au(k_0) + au(k_0) = 0$$

and hence u_3 is conserved. The equation satisfied by u_1, u_2 is then just a reformulation of equation (42). \square

In order to solve (43) we first focus on the special case $u_3 = 0$,

$$\partial_t \begin{pmatrix} u_1 \\ u_2 \end{pmatrix} + \begin{pmatrix} 0 & a \\ -b & 0 \end{pmatrix} \begin{pmatrix} u_1 \\ u_2 \end{pmatrix} = 0. \tag{44}$$

That is, we study the homogeneous solutions of (43). In Lemma 7 we then construct a particular solution for the case $u_3 \neq 0$.

As a first heuristic, note that on $I_1 = (t_0, -\frac{d}{\xi})$ it seems reasonable to approximate

$$a(t) = \frac{c}{\xi^{-2} + t^2} \approx \frac{c}{t^2}, \tag{45}$$

$$b(t) \approx 2c. \tag{46}$$

This approximated problem can be explicitly solved and suggests that

$$\begin{pmatrix} u_1 \\ u_2 \end{pmatrix} \approx \begin{pmatrix} \frac{\gamma_1}{c} |t|^{\gamma_1-1} & \frac{\gamma_2}{c} |t|^{\gamma_2-1} \\ |t|^{\gamma_1} & |t|^{\gamma_2} \end{pmatrix} \begin{pmatrix} \alpha \\ \beta \end{pmatrix}, \tag{47}$$

for suitable constants α, β and $\gamma_1 = \frac{1}{2} + \frac{1}{2}\sqrt{1 - 8c^2}$, $\gamma_2 = \frac{1}{2} + \frac{1}{2}\sqrt{1 - 8c^2}$. In the following we will show that this heuristic is indeed valid in the sense that the actual solution operator has the same asymptotic power law behavior as $|t|$ becomes small.

In order to establish these asymptotics we make use of a self-improving estimate. That is we will assume for the moment that there exists some $\sigma < \infty$ such that

$$|u_1(t)| \leq C|t|^{-\sigma} |u(t_0)|, \tag{48}$$

and show that then u_1 and u_2 necessarily also satisfy the upper lower bounds given in Proposition 4. This a priori assumption is established at the end of this Section in Lemma 8.

Lemma 6. (Improvement) *Consider the problem (44) on the interval $I_1 = (t_0, -\frac{d}{\xi})$ and suppose that u_1 satisfies the upper bound (48).*

Then there exists a constant $C > 0$ (independent of ξ, ε , possibly dependent on c), such that

$$|u_1(t)| \leq C|t|^{\gamma_2-1} |u(t_0)|, \tag{49}$$

$$|u_2(t)| \leq C|t|^{\gamma_2} |u(t_0)|, \tag{50}$$

where $\gamma_2 = \frac{1}{2} - \frac{1}{2}\sqrt{1 - 8c^2}$.

Furthermore, if $|u_2(t_0)| \geq 0.5|u(t_0)|$, then at time $t = -\frac{d}{\xi}$ it holds that

$$\begin{aligned} |u_1(t)| &\geq C/2 \left(\frac{d}{\xi}\right)^{\gamma_2-1} |u_2(t_0)|, \\ |u_2(t)| &\leq C/2 \left(\frac{d}{\xi}\right)^{\gamma_2} |u_2(t_0)|. \end{aligned} \tag{51}$$

An analogous result holds on the interval $I_3 = (\frac{d}{\xi}, t_1)$.

Proof of Lemma 6. We note that $|b(t) - 2c| \leq 10ct$ by Taylor’s formula and hence treat it as a perturbation.

A toy model In order to introduce our method of proof let us first consider the approximated problem where we replace $a(t)$ by exactly $\frac{c}{t^2}$. Then our problem can be written as

$$\partial_t u + \begin{pmatrix} 0 & \frac{c}{t^2} \\ -2c & 0 \end{pmatrix} u = \begin{pmatrix} 0 \\ \mathcal{O}(c|t)u_1 \end{pmatrix},$$

where by assumption (48) it holds that $\mathcal{O}(c|t)u_1 = c\mathcal{O}(|t|^{1-\sigma})$. We then consider a variation of constants ansatz:

$$u(t) = \begin{pmatrix} \frac{\gamma_1}{c} |t|^{\gamma_1-1} & \frac{\gamma_2}{c} |t|^{\gamma_2-1} \\ |t|^{\gamma_1} & |t|^{\gamma_2} \end{pmatrix} \begin{pmatrix} \alpha \\ \beta \end{pmatrix}.$$

Plugging this in, we get that

$$\begin{aligned} &\begin{pmatrix} \frac{\gamma_1}{c} |t|^{\gamma_1-1} & \frac{\gamma_2}{c} |t|^{\gamma_2-1} \\ |t|^{\gamma_1} & |t|^{\gamma_2} \end{pmatrix} \partial_t \begin{pmatrix} \alpha \\ \beta \end{pmatrix} = \begin{pmatrix} 0 \\ \mathcal{O}(ct)u_1(t) \end{pmatrix} \\ \Leftrightarrow \partial_t \begin{pmatrix} \alpha \\ \beta \end{pmatrix} &= \frac{c}{\gamma} \begin{pmatrix} |t|^{\gamma_2} & -\frac{\gamma_2}{c} |t|^{\gamma_2-1} \\ -|t|^{\gamma_1} & \frac{\gamma_1}{c} |t|^{\gamma_1-1} \end{pmatrix} \begin{pmatrix} 0 \\ \mathcal{O}(ct)u_1(t) \end{pmatrix}. \end{aligned}$$

Now using that $\mathcal{O}(ct)|u_1| \leq Cc|t|^{1-\sigma}$, we get that

$$\begin{aligned} |\partial_t \alpha| &\leq Cc|t|^{\gamma_2-1}|t|^{1-\sigma} = Cc|t|^{\gamma_2-\sigma}, \\ |\partial_t \beta| &\leq Cc|t|^{\gamma_1-\sigma}. \end{aligned}$$

Integrating, we obtain that

$$\begin{aligned} |\alpha| &\leq C + C|t|^{1+\gamma_2-\sigma}, \\ |\beta| &\leq C + C|t|^{1+\gamma_1-\sigma}. \end{aligned}$$

Plugging this into our ansatz it follows that

$$\begin{aligned} |u_1(t)| &\leq |t|^{\gamma_1-1}(C + |t|^{1+\gamma_2-\sigma}) + |t|^{\gamma_2-1}(C + |t|^{1+\gamma_1-\sigma}) \\ &\leq C|t|^{\gamma_2-1} + C|t|^{1-\sigma}, \\ |u_2(t)| &\leq |t|^{\gamma_1}(C + |t|^{1+\gamma_2-\sigma}) + |t|^{\gamma_2}(C + |t|^{1+\gamma_1-\sigma}) \\ &\leq C|t|^{\gamma_1} + C|t|^{2-\sigma} + C|t|^{2-\sigma}. \end{aligned}$$

Here, we used that $\gamma_1 + \gamma_2 = 1$, but may also more roughly estimate $\gamma_1 + \gamma_2 \geq 2\gamma_2 > 0$.

In particular, we observe that

$$|u_1(t)| \leq C|t|^{\gamma_2-1} + C|t|^{2\gamma_2-\sigma} \leq C|t|^{-\sigma}$$

is a strict improvement over (48) if $\sigma > 1 - \gamma_2$ initially. We may then repeat our argument with (48) with this smaller σ and successively improve until $\sigma = 1 - \gamma_2$.

With this improved upper bound at hand, we now can establish the lower bounds (51). For this we observe that

$$\begin{aligned} \begin{pmatrix} \alpha(t_0) \\ \beta(t_0) \end{pmatrix} &= S(t_0)^{-1} \begin{pmatrix} u_1(t_0) \\ u_2(t_0) \end{pmatrix} \\ &\approx \frac{c}{\sqrt{1-8c^2}} \begin{pmatrix} \gamma_1/ct_0^{\gamma_2-1} & -t_0^{\gamma_2} \\ -\gamma_2/ct_0^{\gamma_1-1} & t_0^{\gamma_1} \end{pmatrix} \begin{pmatrix} u_1(t_0) \\ u_2(t_0) \end{pmatrix}. \end{aligned}$$

Under the assumption of the lower bound on $u_2(t_0)$ it thus follows that $\beta(t_0)$ is bounded below and of comparable size to $u_2(t_0)$. It thus only remains to show that $(\alpha(t), \beta(t))$ do not deviate too much from this. Here, we use the above developed bounds on $\partial_t \alpha$ and $\partial_t \beta$ to get

$$\begin{aligned} |\partial_t \beta| &\leq \frac{c}{\gamma} \frac{\gamma_1}{c} |t^{\gamma_1}| |u_1(t)| |t| Cc, \\ \rightsquigarrow |\beta(t) - \beta(t_0)| &= \mathcal{O}(c), \\ |\partial_t \alpha| &\leq \frac{c}{\gamma} \frac{\gamma_2}{c} |t|^{\gamma_2-1} Cc |t| |u_1(t)| \\ \rightsquigarrow |\alpha(t) - \alpha(t_0)| &= \mathcal{O}(c^3), \end{aligned}$$

where we used that $\gamma_2 = \frac{1}{2} - \sqrt{\frac{1}{4} - 2c^2} = \mathcal{O}(c^2)$. Since c is small, $\alpha(t), \beta(t)$ hence indeed do not deviate too much and the result follows.

Proof in the actual model This proof is analogous with the only difference being that the solution matrix is now not anymore given by power laws but rather by explicit hypergeometric/Legendre functions, which are bounded above and below by such power laws.

We note that the problem

$$\partial_t \psi + \begin{pmatrix} 0 & \frac{c}{\xi^{-2}+t^2} \\ -2c & 0 \end{pmatrix} \psi = 0,$$

can be decoupled into two second order equations and that ψ_1 solves

$$\partial_t (\xi^{-2} + t^2) \partial_t \psi_1 + 2c^2 \psi_1 = 0.$$

After some changes of coordinates one sees that this is a Legendre type equation

$$(1 - t^2)y''(t) - 2ty'(t) + v(v + 1)y(t) = 0$$

for imaginary t , which has the general solution

$$\psi_1(t) = C_1 P\left(-\frac{1}{2} + \frac{1}{2}\sqrt{1-8c^2}, it\xi\right) + C_2 Q\left(-\frac{1}{2} + \frac{1}{2}\sqrt{1-8c^2}, it\xi\right),$$

where P is the Legendre function of the first kind and Q is the Legendre function of the second kind (see the NIST Digital Library of Mathematical Functions [10] <https://dlmf.nist.gov/14>). We here write $P(v, z) = P_v(z) = P_v^0(z)$. Most importantly, we are interested in the interval of t where $t\xi \in (c^{-1}, \xi)$ is very large and thus can use the asymptotics (<https://dlmf.nist.gov/14.8>)

$$\begin{aligned} P(v, z) &\sim c_1 z^v, \\ Q(v, z) &\sim c_2 z^{-v-1}, \end{aligned}$$

which for our purposes becomes $t^{\gamma_1-1}, t^{\gamma_2-1}$. Here, $c_{1,2}$ are explicit (if complicated) coefficients in terms of v .

The corresponding values of ψ_2 can then be obtained from

$$\psi_2 = -\frac{\xi^{-2} + t^2}{c} \partial_t \psi_1,$$

and the recurrence relation (see <https://dlmf.nist.gov/14.10>)

$$(1-x^2) \frac{d}{dx} P(v, x) = v P(v-1)(x) - vx P_v(x).$$

Note that here $v \sim -\gamma_1, -\gamma_2$ was negative. Thus, for our purposes, these functions asymptotically behave as linear combinations of $(\xi t)^{\gamma_1}, (\xi t)^{\gamma_2}$.

Given these homogeneous solutions, we may thus again make the ansatz that

$$u(t) = S(t) \begin{pmatrix} \alpha(t) \\ \beta(t) \end{pmatrix},$$

where all entries of $S(t)$ are comparable to the power laws obtained above and $\det(S(t)) = \det(S(t_0))$ is bounded away from zero. We remark that the conservation of the determinant follows from

$$\partial_t S + MS = 0,$$

and note that the trace of M vanishes. \square

Given the solution operator of problem (44), we return to problem (43) and introduce a particular solution to account for u_3 .

Lemma 7. (Particular Solution) *Consider the problem (43) with initial conditions*

$$u_1(t_0) = u_2(t_0) = 0.$$

Then in the notation of Lemma 6 it holds that

$$\begin{aligned} |u_1\left(-\frac{d}{\xi}\right)| &\leq c|u_3(t_0)|C\left(\frac{d}{\xi}\right)^{\gamma_2-1}, \\ |u_2\left(-\frac{d}{\xi}\right)| &\leq c|u_3(t_0)|C\left(\frac{d}{\xi}\right)^{\gamma_2}. \end{aligned}$$

An analogous result holds on $I_3 = (\frac{d}{\xi}, t_1)$.

Proof of Lemma 7. Similarly to the proof of Lemma 6 we consider a variation of constants ansatz but this time using the solution operator $S(t)$ of the exact problem (that is, not approximating b). Here, we introduce

$$f = \begin{pmatrix} (b_1 + b_2)u_3 \\ (b_1 - b_2)u_3 \end{pmatrix}$$

in order to simplify notation. In particular recall that u_3 is conserved and hence $|f(t)| \leq 2c|u_3(t_0)|$ for all times.

We hence consider the ansatz

$$u(t) = S(t) \begin{pmatrix} \alpha(t) \\ \beta(t) \end{pmatrix},$$

where $\alpha(t_0) = \beta(t_0) = 0$ and $S(t)$ satisfies the power law bounds discussed in the proof of Lemma 6. In particular, $|S^{-1}(t)| \leq C|t|^{\gamma_2-1}$ is integrable in time (by using the conserved determinant and Cramer's rule).

Thus, it follows that

$$\begin{aligned} S(t)\partial_t \begin{pmatrix} \alpha(t) \\ \beta(t) \end{pmatrix} &= f, \\ \Leftrightarrow \partial_t \begin{pmatrix} \alpha(t) \\ \beta(t) \end{pmatrix} &= S(t)^{-1}f, \end{aligned}$$

and hence,

$$(\alpha(t), \beta(t)) = \int_{t_0}^t \partial_s(\alpha(s), \beta(s)) \, ds$$

is bounded by $Cc|u_3(t_0)|$. \square

Finally, we return to the proof of the a priori bound by a power law (48).

Lemma 8. Consider the problem (44) on the interval $(t_0, -\delta)$ with $\xi^{-1} \ll \delta \ll 1$, then there exists a constant $C > 0$ and $\sigma > 0$ such that

$$|u_1(t)| \leq C|t|^{-\sigma} \|(u_1, u_2)(t_0)\|$$

for all $t \in (t_0, -\delta)$. Here, C and σ may depend on c but are uniform in ξ and δ .

Analogously, for the interval (δ, t_1) it holds that

$$|u_1(t)| \leq C|t|^{-\alpha} \|(u_1, u_2)(t_1)\|.$$

Proof of Lemma 8. Our proof is based on concatenating estimates on a large number of very small intervals, where we can use perturbative arguments.

Thus, consider an interval $I = (t_0, t_0 + \nu) \subset (-1, \delta)$ with t_0 given and ν to be specified later. Then on this interval it holds that

$$|b(t) - b(t_0)| \leq 4c\nu.$$

We thus write our system as

$$\partial_t \begin{pmatrix} u_1 \\ u_2 \end{pmatrix} + \begin{pmatrix} 0 & a \\ -b(t_0) & 0 \end{pmatrix} \begin{pmatrix} u_1 \\ u_2 \end{pmatrix} = \begin{pmatrix} 0 & 0 \\ b(t) - b(t_0) & 0 \end{pmatrix} \begin{pmatrix} u_1 \\ u_2 \end{pmatrix}.$$

Considering the right-hand-side as a perturbation, we make the ansatz

$$\begin{pmatrix} u_1 \\ u_2 \end{pmatrix} = S(t) \begin{pmatrix} \alpha(t) \\ \beta(t) \end{pmatrix}, \tag{52}$$

where $S(t)$ is the solution operator to the problem with b replaced by $b(t_0)$. Then by Lemma 6 we know that component-wise

$$S(t) \sim \begin{pmatrix} -\frac{\gamma_1}{b(t_0)}|t|^{\gamma_1-1} & -\frac{\gamma_2}{b(t_0)}|t|^{\gamma_2-1} \\ |t|^{\gamma_1} & |t|^{\gamma_2} \end{pmatrix},$$

where $\gamma_{1,2}$ solve $\gamma_i(\gamma_i - 1) = cb(t_0)$.

We further recall that $\det(S(t))$ is conserved due the vanishing trace in equation (44) and with the above comparison of the following size:

$$\det(S) \sim \frac{\gamma_2 - \gamma_1}{b(t_0)}.$$

Hence, we may easily invert this matrix using Cramer’s rule.

Plugging in this ansatz, we obtain that

$$\partial_t \begin{pmatrix} \alpha(t) \\ \beta(t) \end{pmatrix} = S(t)^{-1} \begin{pmatrix} 0 & 0 \\ b(t) - b(t_0) & 0 \end{pmatrix} S(t) \begin{pmatrix} \alpha(t) \\ \beta(t) \end{pmatrix}.$$

We now consider the associated Duhamel integral and show that for ν sufficiently small it yields a small perturbation to the identity. That is, we estimate

$$\begin{aligned} & \int_{t_0}^{t_0+\nu} \left\| S(t)^{-1} \begin{pmatrix} 0 & 0 \\ b(t) - b(t_0) & 0 \end{pmatrix} S(t) \right\|_{op} dt \\ & \leq \frac{b_0}{\gamma_1 - \gamma_2} 4c\nu \int_{t_0}^{t_0+\nu} \left\| \begin{pmatrix} |t|^{\gamma_2} & \frac{\gamma_2}{b(t_0)}|t|^{\gamma_2-1} \\ -|t|^{\gamma_1} & -\frac{\gamma_1}{b(t_0)}|t|^{\gamma_1-1} \end{pmatrix} \begin{pmatrix} 0 & 0 \\ 1 & 0 \end{pmatrix} \begin{pmatrix} -\frac{\gamma_1}{b(t_0)}|t|^{\gamma_1-1} & -\frac{\gamma_2}{b(t_0)}|t|^{\gamma_2-1} \\ |t|^{\gamma_1} & |t|^{\gamma_2} \end{pmatrix} \right\|_{op} \\ & = \frac{b_0}{\gamma_1 - \gamma_2} 4c\nu \int_{t_0}^{t_0+\nu} \left\| \begin{pmatrix} 0 & \frac{\gamma_2}{b(t_0)}|t|^{\gamma_2-1} \\ 0 & -\frac{\gamma_1}{b(t_0)}|t|^{\gamma_1-1} \end{pmatrix} \begin{pmatrix} 0 & 0 \\ 1 & 0 \end{pmatrix} \begin{pmatrix} -\frac{\gamma_1}{b(t_0)}|t|^{\gamma_1-1} & -\frac{\gamma_2}{b(t_0)}|t|^{\gamma_2-1} \\ 0 & 0 \end{pmatrix} \right\|_{op}. \end{aligned}$$

As $|t| < 1$, we obtain an upper bound of the integrand by

$$\begin{aligned} & \frac{1}{b(t_0)^2} \left(\gamma_2^2 |t|^{2\gamma_2-2} + 2\gamma_1\gamma_2 |t|^{-1} + \gamma_1^2 |t|^{2\gamma_1-2} \right) \\ & \leq 4c^2 |t|^{-2} + 4|t|^{-1} + 4c^{-2} |t|^{2\gamma_2-1}, \end{aligned}$$

where we used that $c \leq b(t_0) \leq 4c$, $\gamma_1 + \gamma_2 = 1$, $\gamma_2 \approx 4c^2$ (by Taylor’s approximation). Here, the last term is integrable in time, the middle term yields a logarithm and the first term is estimated rather roughly.

We may hence control

$$\int_{t_0}^{t_0+\nu} \left\| S(t)^{-1} \begin{pmatrix} 0 & 0 \\ b(t) - b(t_0) & 0 \end{pmatrix} S(t) \right\|_{op} \tag{53}$$

$$\leq \frac{b_0}{\gamma_1 - \gamma_2} 4c\nu \left(\frac{c^2}{t_0 + \nu} + \ln \left(\frac{t_0 + \nu}{t_0} \right) + 1 \right). \tag{54}$$

Here we used that $t_0 < t_0 + \nu < 0$ to bound the first integral by the larger term. The case $0 < t_0 < t_0 + \nu$ is switched accordingly. We hence note that we cannot choose $t_0 + \nu$ arbitrarily small, but rather need to require that $\frac{t_0+\nu}{t_0}$ (respectively its reciprocal) is not too large. Choosing ν such that

$$t_0 + \nu = \frac{1}{2}t_0, \tag{55}$$

we thus obtain that (53) is bounded by $4c \ll 1$, and hence the map

$$\begin{pmatrix} \alpha(t_0) \\ \beta(t_0) \end{pmatrix} \mapsto \begin{pmatrix} \alpha(t_0 + \nu) \\ \beta(t_0 + \nu) \end{pmatrix}$$

is comparable to the identity within a factor 2, provided ν satisfies (55). It thus follows that $|u(t)|$ also only grows by a constant factor bounded by 100 on the interval $(t_0, t_0 + \nu) = (t_0, t_0/2)$. In order to establish the desired bound on all of $I_1 = (t_0, \delta)$, we partition our interval as

$$(t_0, t_0/2), (t_0/2, t_0/4), \dots, (2\delta, \delta).$$

In order to reach a given $t \in I_1$, we then concatenate $\lfloor \ln(t/t_0) / \ln(2) \rfloor$ intervals and thus obtain that

$$\|(u_1(t), u_2(t))\| \leq 100^{|\ln(t/t_0)/\ln(2)|} \|(u_1, u_2)(t_0)\| = |t/t_0|^{-\ln(100)/\ln(2)} \|(u_1, u_2)(t_0)\|.$$

This concludes the proof of our rough upper bound with $\sigma = \ln(100)/\ln(2)$. \square

5.1.2. The Interval I_2 In order to study the evolution of (39) on the middle interval I_2 we use a different approach based on the convergent limit of iterated Duhamel integration. This method of proof also readily extends to the full (inhomogeneous) problem, which is discussed in more detail in Proposition 6. We again expect the behavior to be similar to the solution derived in Lemma 3 for the two mode model and thus exhibit growth and decay by powers ξ^{+1}, ξ^{-1} , as formulated in the following lemma:

Lemma 9. (Middle interval, homogeneous case) *In the homogeneous case on the middle interval $I_2 = (-\frac{d}{\xi}, \frac{d}{\xi})$, u_3 is again conserved and*

$$\begin{pmatrix} u_1 \left(\frac{d}{\xi} \right) \\ u_2 \left(\frac{d}{\xi} \right) \end{pmatrix} \approx \begin{pmatrix} 1 & Cc\xi \\ \frac{2cd}{\xi} & 1 \end{pmatrix} \begin{pmatrix} u_1 \left(-\frac{d}{\xi} \right) \\ u_2 \left(-\frac{d}{\xi} \right) \end{pmatrix},$$

where $C = \arctan(\xi t) \Big|_{-\frac{d}{\xi}}^{\frac{d}{\xi}} = 2 \arctan(d)$.

In particular, at time $t_2 = +\frac{d}{\xi}$ it holds that

$$u_1(t_2) \approx Cc\xi u_2(t_1) \approx 2Cc^{1-\gamma_2}\xi^{1-\gamma_2} \frac{\gamma_1}{\gamma_2 - \gamma_1} u_2(t_0),$$

$$u_2(t_2) \approx u_2(t_1) \approx 2c^{-\gamma_2}\xi^{-\gamma_2} \frac{\gamma_1}{\gamma_2 - \gamma_1} u_2(t_0).$$

Proof. An explicit solution can be given in terms of hypergeometric functions, as we explore in the ‘‘Appendix A’’. A more useful and shorter proof for our purposes considers an expansion in terms of the Duhamel iteration. Here, it turns out that the first Duhamel iteration is dominant and all higher order Duhamel iterations can be estimated in a geometric series. Since the statement of the inhomogeneous case includes this result as a special case, in order to avoid duplication, we refer to the proof of Proposition 6 for details. Concerning the asymptotics, we recall from Proposition 4 that

$$u_1(t_1) \approx c^{2-\gamma_2}\xi^{1-\gamma_2} u_2(t_0),$$

$$u_2(t_1) \approx c^{-\gamma_2}\xi^{-\gamma_2} u_2(t_0).$$

Hence, it follows that

$$\begin{aligned} & \begin{pmatrix} 1 & Cc\xi \\ \frac{2cd}{\xi} & 1 \end{pmatrix} \begin{pmatrix} u_1(t_1) \\ u_2(t_1) \end{pmatrix} \\ & \approx u_2(t_0) \begin{pmatrix} 1 & cC\xi \\ \frac{2}{\xi} & 1 \end{pmatrix} \begin{pmatrix} c^{2-\gamma_2}\xi^{1-\gamma_2} \\ 2c^{-\gamma_2} \end{pmatrix} \\ & \approx u_2(t_0) \begin{pmatrix} 2Cc^{1-\gamma_2}\xi^{1-\gamma_2} \\ 2c^{-\gamma_2}\xi^{-\gamma_2} \end{pmatrix}. \end{aligned}$$

□

5.1.3. The Interval I_3 It remains to discuss the evolution on the interval $I_3 = (\frac{d}{\xi}, t_1)$. Since this interval roughly corresponds to a mirror image of I_1 , we follow a similar method of proof as in Section 5.1.1, except that the time direction is in a sense reversed. For this reason, we additionally have to consider the inverse of the solution matrix.

Lemma 10. (Right interval, homogeneous case) *Let $\xi \gg 1$ and $0 < c < 0.2$ be given. We consider the problem (39) on the interval $I_3 = (\frac{d}{\xi}, t_1)$. Then under the assumptions of Theorem 7 it holds that*

$$u_1(t_1) \approx c^{-1-\gamma_2}\xi^{-\gamma_2} u_1\left(\frac{d}{\xi}\right) - c^{1-\gamma_2}\xi^{1-\gamma_2} u_2\left(\frac{d}{\xi}\right) \approx \xi^\gamma c^{1-2\gamma_2} \xi^\gamma u_2(t_0)$$

$$u_2(t_1) \approx 2c^{1-\gamma_2}\xi^{-\gamma_2} u_1\left(\frac{d}{\xi}\right) - 2c^{2-\gamma_2}\xi^{1-\gamma_2} u_2\left(\frac{d}{\xi}\right) \approx \xi^\gamma c^{2-2\gamma_2} u_2(t_0),$$

$$u_3(t_1) = u_3(t_0).$$

We note that $u_1\left(\frac{d}{\xi}\right)$ is multiplied by a negative power of ξ while $u_2\left(\frac{d}{\xi}\right)$ is multiplied with a positive power of ξ .

Proof of Lemma 10. This proof is largely analogous to the one of Proposition 4 in the sense that the inverse solution operator $u(t_1) \mapsto u\left(\frac{d}{\xi}\right)$ satisfies the same (asymptotic) estimates.

As a first heuristic let us again consider the approximated two-mode model, where

$$\begin{aligned} \begin{pmatrix} u_1\left(\frac{d}{\xi}\right) \\ u_2\left(\frac{d}{\xi}\right) \end{pmatrix} &\approx \begin{pmatrix} \left(\frac{d}{\xi}\right)^{\gamma_1-1} & \left(\frac{d}{\xi}\right)^{\gamma_2-1} \\ \frac{\gamma_2}{c}\left(\frac{d}{\xi}\right)^{\gamma_1} & \frac{\gamma_1}{c}\left(\frac{d}{\xi}\right)^{\gamma_2} \end{pmatrix} S(t_1)^{-1} \begin{pmatrix} u_1(t_1) \\ u_2(t_1) \end{pmatrix} \\ \Leftrightarrow \begin{pmatrix} u_1(t_1) \\ u_2(t_1) \end{pmatrix} &= S(t_1) \begin{pmatrix} \left(\frac{d}{\xi}\right)^{\gamma_1-1} & \left(\frac{d}{\xi}\right)^{\gamma_2-1} \\ \frac{\gamma_2}{c}\left(\frac{d}{\xi}\right)^{\gamma_1} & \frac{\gamma_1}{c}\left(\frac{d}{\xi}\right)^{\gamma_2} \end{pmatrix}^{-1} \begin{pmatrix} u_1\left(\frac{d}{\xi}\right) \\ u_2\left(\frac{d}{\xi}\right) \end{pmatrix}. \end{aligned}$$

Now note that

$$\begin{pmatrix} \left(\frac{d}{\xi}\right)^{\gamma_1-1} & \left(\frac{d}{\xi}\right)^{\gamma_2-1} \\ \frac{\gamma_2}{c}\left(\frac{d}{\xi}\right)^{\gamma_1} & \frac{\gamma_1}{c}\left(\frac{d}{\xi}\right)^{\gamma_2} \end{pmatrix}^{-1} = \frac{c}{\gamma_1 - \gamma_2} \begin{pmatrix} \frac{\gamma_1}{c}\left(\frac{d}{\xi}\right)^{\gamma_2} & -\left(\frac{d}{\xi}\right)^{\gamma_2-1} \\ -\frac{\gamma_2}{c}\left(\frac{d}{\xi}\right)^{\gamma_1} & \left(\frac{d}{\xi}\right)^{\gamma_1-1} \end{pmatrix}.$$

Thus, in this model, we may compute that

$$\begin{aligned} \begin{pmatrix} u_1(t_1) \\ u_2(t_1) \end{pmatrix} &= \begin{pmatrix} t_1^{\gamma_1} & t_1^{\gamma_2} \\ \frac{\gamma_2}{c}t_1^{\gamma_1-1} & \frac{\gamma_1}{c}t_1^{\gamma_2-1} \end{pmatrix} \begin{pmatrix} \frac{\gamma_1}{c}\left(\frac{d}{\xi}\right)^{\gamma_2} & -\left(\frac{d}{\xi}\right)^{\gamma_2-1} \\ -\frac{\gamma_2}{c}\left(\frac{d}{\xi}\right)^{\gamma_1} & \left(\frac{d}{\xi}\right)^{\gamma_1-1} \end{pmatrix} \begin{pmatrix} u_1\left(\frac{d}{\xi}\right) \\ u_2\left(\frac{d}{\xi}\right) \end{pmatrix} \\ &\approx \begin{pmatrix} \frac{\gamma_1}{c}\left(\frac{d}{\xi}\right)^{\gamma_2} & -\left(\frac{d}{\xi}\right)^{\gamma_2-1} \\ 2c\left(\frac{d}{\xi}\right)^{\gamma_2} & -\frac{\gamma_2}{c}\left(\frac{d}{\xi}\right)^{\gamma_2-1} \end{pmatrix} \begin{pmatrix} u_1\left(\frac{d}{\xi}\right) \\ u_2\left(\frac{d}{\xi}\right) \end{pmatrix} \\ &= \begin{pmatrix} \frac{\gamma_1}{c}d^{\gamma_2} & -d^{\gamma_2-1} \\ 2c d^{\gamma_2} & -\frac{\gamma_2}{c}d^{\gamma_2-1} \end{pmatrix} \begin{pmatrix} \xi^{-\gamma_2}u_1\left(\frac{d}{\xi}\right) \\ \xi^{1-\gamma_2}u_2\left(\frac{d}{\xi}\right) \end{pmatrix} \\ &\approx \begin{pmatrix} \frac{\gamma_1}{c}d^{\gamma_2} & -d^{\gamma_2-1} \\ 2c d^{\gamma_2} & -\frac{\gamma_2}{c}d^{\gamma_2-1} \end{pmatrix} \begin{pmatrix} 2Cc^{1-\gamma_2}\frac{\gamma_1}{\gamma_2-\gamma_1} \\ 2c^{-\gamma_2}\frac{\gamma_1}{\gamma_2-\gamma_1} \end{pmatrix} \xi^\gamma u_2(t_0), \end{aligned}$$

where we considered smaller powers of ξ as errors and denoted $\gamma = 1 - 2\gamma_2 = 2\sqrt{\frac{1}{4} - 2c^2} = \sqrt{1 - 8c^2}$. Choosing $d = c^{-1}$ and recalling that $\gamma_2 \approx 2c^2$ by Taylor's approximation, $\gamma_1 \approx \gamma \approx 1$, and we may further approximate to get

$$\begin{aligned} &\begin{pmatrix} \frac{\gamma_1}{c}d^{\gamma_2} & -d^{\gamma_2-1} \\ 2c d^{\gamma_2} & -\frac{\gamma_2}{c}d^{\gamma_2-1} \end{pmatrix} \begin{pmatrix} 2Cc^{1-\gamma_2}\frac{\gamma_1}{\gamma_2-\gamma_1} \\ 2c^{-\gamma_2}\frac{\gamma_1}{\gamma_2-\gamma_1} \end{pmatrix} \\ &\approx \begin{pmatrix} c^{-1-\gamma_2} & -c^{1-\gamma_2} \\ 2c^{1-\gamma_2} & -2c^{2-\gamma_2} \end{pmatrix} \begin{pmatrix} 2Cc^{1-\gamma_2} \\ 2c^{-\gamma_2} \end{pmatrix} \\ &= \begin{pmatrix} Cc^{2\gamma_2} - 2c^{1-2\gamma_2} \\ 2Cc^{2-2\gamma_2} - 4c^{2-2\gamma_2} \end{pmatrix} \approx \begin{pmatrix} -2c^{1-2\gamma_2} \\ \tilde{C}c^{2-2\gamma_2} \end{pmatrix}. \end{aligned}$$

In particular, we stress that $u_1(t_1) \geq c^{-1}u_2(t_1)$ is much bigger provided that c is sufficiently small.

It remains to discuss the extension to the full problem. We note that under the transformation $(u_1(t), u_2(t), u_3(t)) \mapsto (-u_1(-t), u_2(-t), u_3(-))$ we obtain an analogous problem as the one considered on I_1 in Section 5.1.1 except that t_0 is replaced by $-t_1$ and $b(t)$ is replaced by $b(-t)$. By the same arguments as in Section 5.1.1 it hence follows that the solution operator $u(t) := S(t)u(t_1)$ is of the form

$$S = \begin{pmatrix} S' & b \\ 0 & 0 & 1 \end{pmatrix},$$

where S' is asymptotically well-approximated by our two-mode model heuristic.

We may then explicitly compute the inverse of S as

$$S(t)^{-1} = \begin{pmatrix} S'^{-1} & \tilde{b} \\ 0 & 0 & 1 \end{pmatrix},$$

$$\begin{pmatrix} \tilde{b} \\ \tilde{c} \end{pmatrix} = -S'^{-1} \begin{pmatrix} b \\ c \end{pmatrix},$$

where we can use Cramer’s rule for inverting S . As we know that $S(\frac{d}{\xi})$ is well-approximated by the power laws the result for the full model hence follows by the exact same argument as we discussed for the approximate model above. \square

In the following section we use multiple bootstrap arguments to show that this behavior persists in the full problem.

5.2. The Full Model

Building on the results for the three-mode model of Section 5.1 we show in the following that the full dynamics have the same asymptotic behavior. In particular, this shows that if the logarithmic constraint imposed in prior works is removed this imposes not only technical challenges but results in different asymptotics. In Section 5.3 we crucially rely on this improvement to construct solutions with initial data in Gevrey regularity which exhibit infinitely many cascades of arbitrarily long length. Furthermore, they *diverge* in (any) Sobolev regularity but their force field nevertheless *converges*. Damping persists despite the failure of scattering. The estimates for a single resonance in the full problem are summarized in the following theorem:

Theorem 8. (Summary) *Let $c \in \mathbb{R}$ with $0 < |c| < 0.1$ and let $\xi \gg 1$. Let u be the solution of the full model (34):*

$$\partial_t u(k) + a(k + 1)u(k + 1) - a(k - 1)u(k - 1) = 0,$$

on (t_0, t_1) . Suppose that at time t_0 it holds that $u(k_0) \geq 0.5 \max |u| =: 0.5\theta$ and define $\gamma = \sqrt{1 - 8c^2}$. Then there exists a constant $C > 0$ such that it holds that, at time $t = t_1$,

$$|u(k)| \leq C\xi^\gamma\theta$$

for all k and

$$|u(k_0 - 1)| \geq \frac{C}{2}\xi^\gamma\theta.$$

In particular, u at time t_1 satisfies the assumptions imposed at time t_0 . We may therefore apply this theorem again on the next time interval.

As in Section 5.1 we split our proof into considering the three intervals I_1, I_2, I_3 and then combining the respective estimates to establish Theorem 8. On the intervals I_1 and I_3 we employ a bootstrap approach and a variation of constants ansatz to control the effect of non-resonant modes. In the center interval I_2 we instead are able to explicitly construct a solution operator by Duhamel iteration. We stress that if the logarithmic smallness condition is imposed this approach can be used on the full interval (see Section 3) but that this cannot work without this constraint. Indeed, higher order Duhamel iterations are much larger (by powers of logarithms) than the first iteration and their cancellations yield the modified exponent.

5.2.1. The Interval I_1 As in Section 5.1.1, our main aim of this section is to show that on the interval I_1 now for the full solution the mode $u(k_0 \pm 1)$ increases with a power law $\xi^{1-\gamma_2}$, while $u(k_0)$ decreases with a power law $\xi^{-\gamma_2}$, with γ_2 as in Theorem 8. Compared to Proposition 4 a main challenge here is given by the coupling to other modes than $k_0 - 1, k_0, k_0 + 1$, which appear as a right-hand-side in the inhomogeneous problem and themselves evolve. Our strategy in the following is thus to assume that we can control all modes in an appropriate way at least for some small positive time, which is quantified in the bootstrap ansatz (59) to (63). Using the equations we then show that all controls improve on that time interval and hence must remain valid on the whole interval I_1 . We summarize our results in the following proposition:

Proposition 5. (Left interval, inhomogeneous case) *Let $c \in \mathbb{R}$ with $0 < |c| < 0.1$ and let $\xi \gg 1$. Let u be as in Theorem 8. Then at $-\frac{d}{\xi} = -\frac{1}{c\xi}$ it holds that*

$$u(k_0) \approx \theta c \left(\frac{d}{\xi}\right)^{\gamma_2}, \tag{56}$$

$$u(k_0 \pm 1) \approx \theta c \left(\frac{d}{\xi}\right)^{\gamma_2-1}, \tag{57}$$

$$u(k) \leq C\theta, \tag{58}$$

where $C = (e^c - 1)$

Proof. By linearity without loss of generality we may set $\theta = 1$.

Then at least for small positive time it holds that

$$u(k_0) \approx cu_2(t_0)|t|^{\gamma_2}, \tag{59}$$

$$(u(k_0 + 1) - u(k_0 - 1)) \approx cu_2(t_0)|t|^{\gamma_2-1}, \tag{60}$$

$$\left(\frac{u(k_0 + 1) + u(k_0 - 1)}{2}\right)_{t_0}^t \leq e^{4.1c(t-t_0)} - 1 + (t^{\gamma_2} - t_0^{\gamma_2}), \tag{61}$$

$$u(k) \leq e^{2.1c(t+1)} - 1 + C(t^{\gamma_2} - t_0^{\gamma_2}) \text{ else,} \tag{62}$$

$$\|u \cdot 1_{|k-k_0| \geq 2}\|_{l^2}(t) \leq e^{4.1c(t-t_0)} \|u \cdot 1_{|k-k_0| \geq 2}\|_{l^2}(t_0). \tag{63}$$

Here $1_{|k-k_0| \geq 2}(k) \in \{0, 1\}$ denotes the characteristic function and we denoted $u_2(t_0) = u(k_0)|_{t=t_0}$ in analogy to the notation of the homogeneous model of Section 5.1.

In what follows we argue by bootstrap that the maximal time satisfying these estimates is given by $T = \frac{1}{c\xi}$.

Indeed, suppose there were $T < \frac{1}{c\xi}$ maximal with these properties. We then show that all conditions (59) to (62) do not achieve equality at time T and that hence T could be chosen larger by continuity, which contradicts the maximality.

Ad (61) In a fashion similar as to Lemma 5, we may use the symmetry of the problem to compute that

$$\begin{aligned} \partial_t(u(k_0 + 1) + u(k_0 - 1)) &= -a(k_0 + 2)u(k_0 + 2) + a(k_0 - 2)u(k_0 - 2) \\ &\Rightarrow |u(k_0 + 1) + u(k_0 - 1)|_{t_0}^t \leq c \int_{t_0}^T (|u(k_0 + 2)| \\ &\quad + |u(k_0 - 2)|) \\ &\leq 2c \int_{t_0}^T e^{4.1c(t-t_0)} + (t^{\gamma_2} - t_0^{\gamma_2}). \end{aligned}$$

We may then roughly control $|t^{\gamma_2} - t_0^{\gamma_2}| \leq 1 \leq e^{4.1c(t-t_0)}$ and thus control the integral by

$$\frac{4}{4.1} \left(e^{2.1c(T-t_0)} - 1 \right) < \left(e^{4.1c(T-t_0)} - 1 \right).$$

Thus, equality in (61) is not attained at time T .

Ad (62) First suppose in addition that $k \neq k_0 + 2, k_0 - 2$. Then it holds that

$$\begin{aligned} u(k)|_{t_0}^T &= \int_{t_0}^T -a(k + 1)u(k + 1) + a(k - 1)u(k - 1) \leq 4c \int_{t_0}^T e^{4.1c(t-t_0)} dt \\ &= \frac{4}{4.1} \left(e^{2c(T-t_0)} - 1 \right) < e^{2c(T-t_0)} - 1, \end{aligned}$$

where we again estimated $|t^{\gamma_2} - t_0^{\gamma_2}| \leq 1 \leq e^{4.1c(t-t_0)}$ and used that $|a(k)| \leq c$ unless $k = k_0$.

If $k = k_0 \pm 2$, $\int_{t_0}^T \partial_t u(k)$ additionally involves

$$\int_{t_0}^T a(k_0 \pm 1)u(k_0 \pm 1) dt \leq ct^{\gamma_{1,2}}|_{t_0}^T,$$

which is controlled by $t^{\gamma_{1,2}} - t_0^{\gamma_{1,2}}$.

Ad (63) Taking a time-derivative of the energy, that is the left-hand-side, we get that

$$\partial_t E(t) \leq 4cE(t) + c|u(k_0 + 2)||u(k_0 + 1)| \leq 4cE(t) + c(e^{4.1ct} + 2)(t^{\gamma_2} - t_0^{\gamma_2}).$$

The estimate hence follows by Gronwall's lemma or by multiplying with e^{-4ct} and then integrating.

The main part of the proof is thus given by the proof of the estimates for $u(k_0)$ and $u(k_0 + 1) - u(k_0 - 1)$.

Ad (59) and (60) Following a similar argument as in the proof of Proposition 4, we study the inhomogeneous problem for $(\frac{u(k_0+1)-u(k_0-1)}{2}, u(k_0)) =: (u_1, u_2)$:

$$\begin{aligned} \partial_t \begin{pmatrix} u_1 \\ u_2 \end{pmatrix} + \begin{pmatrix} 0 & a \\ b & 0 \end{pmatrix} \begin{pmatrix} u_1 \\ u_2 \end{pmatrix} &= - \begin{pmatrix} a(k_0 + 2)u(k_0 + 2)/2 - a(k_0 - 2)u(k_0 - 2)/2 \\ (a(k_0 + 1) - a(k_0 - 1))(u(k_0 + 1) + u(k_0 - 1))/2 \end{pmatrix} \\ &=: f. \end{aligned}$$

As in Lemma 6 we make a variation of constants ansatz, where for simplicity of notation write our calculations in terms of the power law solutions of the approximate model. By the estimates of Section 5.1, the solution operator $S(t)$ of the full model is comparable, and hence the exact same proof immediately applies:

$$\begin{pmatrix} u_1 \\ u_2 \end{pmatrix} = S(t) \begin{pmatrix} \alpha(t) \\ \beta(t) \end{pmatrix} \approx \alpha(t) \begin{pmatrix} t^{\gamma_1-1} \\ t^{\gamma_1} \frac{\gamma_2}{c} \end{pmatrix} + \beta(t) \begin{pmatrix} t^{\gamma_2-1} \\ t^{\gamma_2} \frac{\gamma_1}{c} \end{pmatrix}.$$

Solving at time t_0 , we obtain

$$\begin{aligned} \begin{pmatrix} \alpha(t_0) \\ \beta(t_0) \end{pmatrix} &= S(t_0)^{-1} \begin{pmatrix} u_1(t_0) \\ u_2(t_0) \end{pmatrix} \\ &\approx \frac{1}{\sqrt{1-8c^2}} \begin{pmatrix} \gamma_1 t_0^{\gamma_2-1} & -ct_0^{\gamma_2} \\ -\gamma_2 t_0^{\gamma_1-1} & ct_0^{\gamma_1} \end{pmatrix} \begin{pmatrix} u_1(t_0) \\ u_2(t_0) \end{pmatrix}. \end{aligned}$$

Here, in the proof of Proposition 4 we concluded by noting that only $\beta(t_1) = \beta(t_1) \approx cu_2(t_0)$ is relevant at t_1 (due to the smaller powers ξ for α). In order to establish the desired bounds in this inhomogeneous case, we hence need to show that α, β do not deviate from this too much. We compute

$$\begin{aligned} \partial_t \begin{pmatrix} \alpha \\ \beta \end{pmatrix} &= -c \begin{pmatrix} t^{\gamma_1-1} & t^{\gamma_2-1} \\ t^{\gamma_1} \frac{\gamma_2}{c} & t^{\gamma_2} \frac{\gamma_1}{c} \end{pmatrix}^{-1} f \\ &= -\frac{c^2}{\sqrt{1-8c^2}} \begin{pmatrix} t^{\gamma_2} \frac{\gamma_1}{c} & -t^{\gamma_2-1} \\ -t^{\gamma_1} \frac{\gamma_2}{c} & t^{\gamma_1-1} \end{pmatrix} f. \end{aligned}$$

By estimates (62) and (61) we note that $|f| \leq c$ and we further note that the integrals of $t^{\gamma_1}, t^{\gamma_1-1}$ are uniformly bounded, while $\int t^{\gamma_2-1} \leq \frac{1}{\gamma_2} \leq c^{-2}$, where we used that $\gamma_2 \approx 2c^2$. Hence, in total it follows that $(\alpha, \beta)|_{t_0}^t \leq C(c, c^3)$ for all times and therefore $(\alpha(t), \beta(t)) \approx (\alpha(t_0), \beta(t_0))$, which implies the result. \square

5.2.2. The Interval I_2 As in Section 5.1.2, we next study the evolution, now of the full solution, on the middle interval $I_2 = (-\frac{d}{\xi}, \frac{d}{\xi})$. Here again our main aim is to quantify the growth and decay of the three modes $k_0 + 1, k_0, k_0 - 1$ and control all other modes. In particular, we again are focused on the power law dependence in terms ξ for ξ large, which arises from the evolution on I_2 and the one on I_1 studied in the previous section.

Since the interval I_2 is shorter the larger ξ is, in what follows we rely on a characterization of the solution in terms of a Duhamel series (see also Section 3 for the associated notion of paths). We then first study the evolution for initial data concentrated on the resonant mode k_0 in Lemma 11, which includes with a power law ξ^1 . If the initial data is instead concentrated on a non-resonant mode as in Lemma 12 the contribution is small. Combining both cases and using the linearity of the solution operator in the initial data, we obtain the following main proposition of this subsection:

Proposition 6. (Middle interval, inhomogeneous case) *Let u be as in Theorem 8. Then at time $\frac{d}{\xi}$ it holds that*

$$\begin{aligned} u(k_0 \pm 1)|_{-d/\xi}^{d/\xi} &\approx Cc\xi u(k_0, t_1) \approx \theta Cc^{1-\gamma_2}\xi^{1-\gamma_2}, \\ u(k_0)|_{-d/\xi}^{d/\xi} &\approx \frac{c}{\xi}u(k_0 \pm 1, t_1) \approx \theta c^{-\gamma_2}\xi^{-\gamma_2}, \\ u(k)|_{-d/\xi}^{d/\xi} &\leq \frac{cd}{\xi}u(k_0 \pm 1, t_1) \leq \theta c^{1-\gamma_2}\xi^{1-\gamma_2}, \end{aligned}$$

where $C = 2 \arctan(d) \approx \pi$.

We here use the Duhamel iteration, which we phrase as properties of the solution map for single mode initial data. Arbitrary initial data can then be realized as a linear combination. We emphasize that on this short time interval the action $k_0 \mapsto k_0 \pm 1$ is large (of size $Cc\xi$), while all other actions are small perturbations of the identity.

Lemma 11. *Suppose that at time $-\frac{d}{\xi}$ it holds that $u(k) = \delta_{k_0}$. Then at time $\frac{d}{\xi}$, u satisfies*

1. $|u(k_0) - 1| \leq \frac{1}{1-c_1} \frac{c_2}{1-c_2}$,
2. $|u(k_0 \pm 1) \mp 2c\xi \arctan(\xi t)|_{-d/\xi}^{d/\xi} \leq \frac{2c}{1-4c^2d} 2c\xi \arctan(\xi t)|_{t_0}^{t_1}$,
3. $|u(k)| \leq |\frac{c}{\xi}|^{|k-k_0|-1} \frac{1}{1-c}$ else,

where $c_1 = \frac{4cd}{\xi} = \frac{4}{\xi}$ and $c_2 = 8c^2d \arctan(\xi t)_{t_1}^2 = 8c \arctan(d) \approx 4\pi c$.

Proof. *Ad (i)* Let $\gamma = (k_0, \dots, k_0)$ be a path starting and ending in k_0 , then we may roughly estimate

$$\iint_{t_0 \leq \tau_1 \leq \dots \leq \tau_l} \prod_i a(\gamma_i, \tau_i) \operatorname{sgn}(\gamma_{i+1} - \gamma_i) \, d\tau_i$$

by

$$\prod \|a(\gamma_i, \tau_i)\|_{L^1([-d/\xi, d/\xi])} = \left(\frac{2cd}{\xi}\right)^{j_1} \left(c\xi \arctan(\xi t)\Big|_{-d/\xi}^{d/\xi}\right)^{j_2},$$

where j_1 corresponds to the number of non-resonance and j_2 to the number of resonances $\gamma_i = k_0$ (recall that the last entry of γ does not appear in the integral). Then in order to start end in k_0 it needs to hold that $j_1 \geq j_2$, since we have to come back to k_0 before leaving it again. Thus, the contribution of the path γ is controlled by

$$\left(\frac{2cd}{\xi}\right)^{j_1-j_2} \left(2c^2d \arctan(\xi t)\Big|_{-d/\xi}^{d/\xi}\right)^{j_2}.$$

Estimating the number of paths of given length $j_1 + j_2$ by $2^{j_1+j_2}$ the sum over the integrals of all such paths can be controlled by

$$\frac{1}{1 - \frac{4cd}{\xi}} \frac{8c^2d \arctan(\xi t)\Big|_{-d/\xi}^{d/\xi}}{1 - 8c^2d \arctan(\xi t)\Big|_{-d/\xi}^{d/\xi}},$$

where we used that $j_2 \geq 1$.

Ad (ii) Let us first consider the special paths $\gamma = (k_0, k_0 \pm 1)$, which yield an integral

$$\mp \int_{-d/\xi}^{d/\xi} a(k_0) d\tau = \mp c\xi \arctan(\xi t)\Big|_{-d/\xi}^{d/\xi} \gg 1.$$

For any other paths γ starting in k_0 and ending in $k_0 \pm 1$, we may again roughly bound the integral by

$$\left(\frac{2cd}{\xi}\right)^{j_1} \left(c\xi \arctan(\xi t)\Big|_{-d/\xi}^{d/\xi}\right)^{j_2},$$

where now $j_1 \geq j_2 - 1 \geq 0$ and we already treated $j_2 = 1, j_1 = 0$ separately. We may thus express this bound as

$$\left(\frac{2cd}{\xi}\right)^{\hat{j}_1} \left(2c^2d \arctan(\xi t)\Big|_{-d/\xi}^{d/\xi}\right)^{\hat{j}_2} \left(c\xi \arctan(\xi t)\Big|_{-d/\xi}^{d/\xi}\right),$$

where $(\hat{j}_1, \hat{j}_2) = (0, 0)$ is excluded. Again estimating the number of such paths from above by $2^{\hat{j}_1+\hat{j}_2+1}$ and summing the geometric series, we obtain the desired result.

Ad (iii) Let $k \notin \{k_0 - 1, k_0, k_0 + 1\}$. Then given a path $\gamma = (k_0, \dots, k)$ there is a last time where $\gamma_i = k_0$, after which the remainder path is non-resonant at least $|k - k_0|$ times. Grouping all paths with the same remainder, we first estimate the contribution by the segments up to the last resonance as in (i) by

$$\frac{1}{1 - \frac{4cd}{\xi}} \frac{8c^2d \arctan(\xi t)\Big|_{-d/\xi}^{d/\xi}}{1 - 8c^2d \arctan(\xi t)\Big|_{-d/\xi}^{d/\xi}}$$

and then estimate the sum over all possible remainders by

$$\sum_{j \geq |k-k_0|} \left(\frac{4cd}{\xi}\right)^j = \frac{1}{1 - \frac{4cd}{\xi}} \left(\frac{4cd}{\xi}\right)^{|k-k_0|},$$

where we again introduced a factor 2^j to account for the number of all remainders of a given length. \square

Lemma 12. *Let $l \neq k_0$ and suppose that at time $-\frac{d}{\xi}$ it holds that $u(k) = \delta_{kl}$. Then u satisfies*

$$u(k)|_{-\frac{d}{\xi}}^{\frac{d}{\xi}} \leq \left(\frac{4cd}{\xi}\right)^{|k-l|} + \left(\frac{4cd}{\xi}\right)^{|k-k_0|+|l-k_0|-1} \left(8c^2d \arctan(t\xi)|_{-\frac{d}{\xi}}^{\frac{d}{\xi}}\right) \times \frac{1}{(1 - \frac{4cd}{\xi})^2} \frac{1}{1 - 8c^2d \arctan(t\xi)|_{-\frac{d}{\xi}}^{\frac{d}{\xi}}}.$$

Here, with slight abuse of notation $|a - b|$ denotes the minimal path length between a and b , i.e. $|a - a| = 2$.

Proof. Let $k \in \mathbb{Z}$ be given. Consider first the case of a purely non-resonant path γ , which can be estimated by

$$\left(\frac{2cd}{\xi}\right)^{|\gamma|}.$$

Any such path has length at least $|k - l|$ (with $|l - l| = 2$). The first summand in the above estimate is hence obtained by a geometric series.

Next consider a path with at least one resonance. We group all paths that share the segment from first to last resonance and observe that at least $|l - k_0|$ non-resonances are needed to reach k_0 for the first time and at least $|k - k_0| - 1$ non-resonances to reach k from k_0 . We then may again control the sum over all path segments from k_0 to k_0 as in the previous lemma and control the first and last segment by a geometric series. \square

Using the result of the preceding lemmas, we can now prove Proposition 6.

Proof of Proposition 6. By linearity we may decompose our initial data $u(t_1)$ as

$$u(t_1) = u(t_1)|_{k=k_0} + u(t_1)|_{k \in \{k_0-1, k_0+1\}} + u^r(t_1).$$

Then by Lemma 12 and the bound on $u^r(t_1)$ established in Proposition 5, we may estimate its contribution to $u(t_2)$ by

$$\left(\left(\frac{4}{\xi}\right)^{|l|} * u^r(t_1)\right)(l) + Cc \left(\frac{4}{\xi}\right)^{|l-k_0|} \left(\left(\frac{4}{\xi}\right)^{|-k_0|-1} * u^r(t_1)\right)(l),$$

where $C = \frac{16 \arctan(d)}{(1 - \frac{4}{\xi})^2 (1 - 16c \arctan(d))}$. This contribution is thus very smooth and small and can be considered a perturbation.

Concerning the contributions by k_0 and $k_0 \pm 1$, we note that Lemmas 11 and 12 combined with the bounds on $u(t_1)$ established in Proposition 5 control $u(k)$ for $k \in \{k_0 - 1, k_0, k_0 + 1\}$ in the desired way and show that

$$\left(\frac{u(k_0+1)-u(k_0-1)}{u(k_0)^2} \right) \Big|_{t=t_2} \approx \begin{pmatrix} 1 & 2Cc\xi \\ \frac{1}{\xi} & 1 \end{pmatrix} \left(\frac{u(k_0+1)-u(k_0-1)}{u(k_0)^2} \right) \Big|_{t=t_1}.$$

We thus conclude as in the proof of Lemma 9. \square

5.2.3. The Interval I_3 It remains to control the evolution of the full equation on the right interval I_3 , where we are gain focused on the power law dependence in terms of ξ . Combining the growth and decay behavior on this interval with the descriptions derived in the previous subsections, we can then conclude our proof of Theorem 8 and thus obtain a description of the growth of our solution on the full resonant interval $I_1 \cup I_2 \cup I_3$. Subsequently, in Section 5.3, we combine several of these intervals (for different k) to construct chains of resonance and associated solutions exhibiting norm inflation. Furthermore, we combine countably many chains to construct solutions with non-trivial asymptotic behavior.

Similarly to Section 5.1.3 we observe that the right interval is basically a mirror image of the left interval I_1 and can hence be treated in a similar way as in Section 5.2.1 except that the time direction is in a sense reversed. Thus, we additionally have to control inverses of the solution operator.

Proposition 7. (Right interval, inhomogeneous case) *Let u be as in Theorem 8. Then it holds that at time t_1 , all modes satisfy*

$$u \leq C\xi^\gamma,$$

and $u(k_0 \pm 1)$ is bounded below by $\frac{C}{2}\xi^\gamma$.

Proof. We again make a bootstrap ansatz similar to the one in Proposition 5, but now take into account the different powers of ξ on different modes:

$$u(k_0 + 1) - u(k_0 - 1) \approx 2\theta c\xi^\gamma t^{1-\gamma_2} \tag{64}$$

$$u(k_0) \approx \theta|\xi|^\gamma t^{\gamma_1}, \tag{65}$$

$$u(k_0 + 1) + u(k_0 - 1) \Big|_{d/\xi}^t \leq \xi^\gamma (e^{4.1c(t-d/\xi)} - 1 + t^{\gamma_1} - d/\xi^{\gamma_1}), \tag{66}$$

$$u(k) \leq \xi^\gamma (e^{4.1c(t-d/\xi)} - 1 + t^{\gamma_1} - d/\xi^{\gamma_1}) \text{ else.} \tag{67}$$

Ad (67) and (66): Here, we again use that $u(k)$, $\frac{u(k_0+1)+u(k_0-1)}{2} \leq \xi^\gamma$ at time d/ξ and thus estimate

$$\int_{d/\xi}^t \partial_t u \leq 2c\xi^\gamma e^{4.1c(t+1)} < \xi^\gamma (e^{4.1c(t+1)} - 1),$$

where we roughly estimate $t^{\gamma_1} - d/\xi^{\gamma_1} \leq 1 \leq e^{4.1c(t-d/\xi)}$ again.

Ad (64) and (65): We make the same variation of constants ansatz as in the proof of Proposition 5 and again for simplicity of notation write the power law approximate solutions instead. Thus consider

$$\begin{pmatrix} u_1 \\ u_2 \end{pmatrix} = S(t) \begin{pmatrix} \alpha(t) \\ \beta(t) \end{pmatrix} \approx \alpha(t) \begin{pmatrix} t^{\gamma_1-1} \\ t^{\gamma_1} \frac{\gamma_2}{c} \end{pmatrix} + \beta(t) \begin{pmatrix} t^{\gamma_2-1} \\ t^{\gamma_2} \frac{\gamma_1}{c} \end{pmatrix}.$$

Here, as in the proof of 6 in the following we present our argument in terms the approximate coefficients for simplicity of notation.

We first compute $\alpha, \beta|_{t=d/\xi}$ by solving

$$\begin{pmatrix} \alpha \\ \beta \end{pmatrix} \Big|_{t=\frac{d}{\xi}} = \begin{pmatrix} t^{\gamma_2} \frac{\gamma_1}{c} & -t^{\gamma_2-1} \\ -t^{\gamma_1} \frac{\gamma_2}{c} & t^{\gamma_1-1} \end{pmatrix} \Big|_{t=\frac{d}{\xi}} \begin{pmatrix} u_1 \\ u_2 \end{pmatrix} \Big|_{t=\frac{d}{\xi}}.$$

Plugging in the relations between $\xi^{-\gamma_2} u_1(t_2)$ and $\xi^{1-\gamma_2} u_2(t_2)$, we see that this contribution satisfies the claimed estimates.

It hence remains again to study the perturbations due to the inhomogeneity. As in the proof of Proposition 5 we compute

$$\partial_t \begin{pmatrix} \alpha \\ \beta \end{pmatrix} = -c \begin{pmatrix} t^{\gamma_1-1} & t^{\gamma_2-1} \\ -t^{\gamma_1} \frac{\gamma_2}{c} & t^{\gamma_2} \frac{\gamma_1}{c} \end{pmatrix}^{-1} f.$$

Plugging in our bounds (66) and (67) by $C\xi^\gamma$ into f and integrating it follows that

$$\left| \begin{pmatrix} \alpha \\ \beta \end{pmatrix} \Big|_{t=\frac{d}{\xi}} \right| \ll c\xi^\gamma.$$

Thus, the value of (α, β) at time $\frac{d}{\xi}$ is dominant and satisfies the estimates. \square

This concludes the proof of Theorem 8, which established growth of certain modes of our solution around the time $\frac{\eta}{k}$ with a rate proportional to $(\frac{\xi}{k^2})^\gamma$. As sketched in Section 3 in the following section we connect this growth on subsequent time intervals (for $k, k - 1, k - 2, \dots$) to form an echo chain, which exhibits norm inflation. Furthermore, since we do not require any smallness condition on η unlike prior works we can consider ξ arbitrarily large. Furthermore, using the linearity of the equation and considering families of frequency which are well separated, we can construct solutions which exhibit infinitely many echo chains of increasing lengths. In particular, these solutions do not become trivial after a final resonant time and can hence exhibit non-trivial asymptotic behavior like blow-up.

5.3. Modified Scattering and Inviscid Damping

In the preceding sections we have shown that the linearized Euler equations (7) exhibit resonances due to the Orr mechanism. In particular, we showed that the logarithmic smallness constraint imposed in prior works is not only a technical restriction but that for large frequencies the asymptotics differ greatly, which manifests in a modified exponent. Furthermore, removing this constraint allows us to consider initial data which is *not* compactly supported in Fourier space.

As our main result we thus show that the linearized Euler equations not only exhibit full echo chains and the associated Gevrey norm inflation but that for data in a critical Gevrey regularity class we obtain blow-up and damping at the same time. That is, we construct smooth initial data $\omega_0 \in \mathcal{G}_{C, \frac{1}{2}}$ which exhibits infinitely many resonance chains of arbitrarily long length. For this solution any Sobolev norm $\|\omega(t)\|_{H^s}$, $s \geq 0$ diverges to infinity as time tends to infinity but despite this the associated velocity perturbation converges. Linear inviscid damping persists despite the failure of scattering and nearby traveling wave-like solutions are the underlying cause of Gevrey blow-up! We stress that one of the main challenges we had to overcome in this work is that one has to be able to consider infinitely many echo chains of increasing length. In particular, one has to allow for arbitrarily large frequencies, while prior works had imposed a logarithmic smallness constraint.

Theorem 9. (Modified Scattering) *Let $0 < c < 0.2$ be given. Then there exists $C_0 = C_0(c)$ such that if $\omega_0 \in \mathcal{G}_{C, \frac{1}{2}}$ with $C > C_0$ then $\omega(t) \in \mathcal{G}_{C-C_0, \frac{1}{2}}$ globally in time and $u(t)$ converges in $\mathcal{G}_{C-C_0, \frac{1}{2}}$ as $t \rightarrow \infty$.*

On the other hand, for every $C < C_0/2$ and every $s \in \mathbb{R}$, there exists $\omega_0 \in \mathcal{G}_{C, \frac{1}{2}}$ such that $\sup_{t \in [0, \infty)} \|\omega(t)\|_{H^\sigma} = \infty$ for any $\sigma \geq s$ but such that $\omega(t)$ converges in H^{s-} as $t \rightarrow \infty$.

Furthermore, for $s \geq 0$ the corresponding velocity field converges strongly in L^2 to a shear flow as $t \rightarrow \infty$. Linear inviscid damping holds despite the divergence of $\omega(t)$ in higher regularity.

Proof. We proceed similarly as in the case of Theorem 6. Consider a frequency $\eta \gg 1$ and let $1 \ll k \leq \sqrt{\eta}$ to be fixed later. Then by the local well-posedness established in Section 2 we may prescribe smooth initial data $\omega_0^{\eta, k}$ such that at the time $t = \frac{\eta}{k}$, ω is given by $e^{i\eta y + ikx}$. Then, we iteratively apply Theorem 8 to obtain that after time $t = 2\eta$, the mode $e^{i\eta y + ix}$ is the largest (within a factor) and of size

$$C^k \left(\frac{\eta^k}{(k!)^2} \right)^\gamma,$$

where we used that $\xi = \frac{\eta}{k^2}$ changes in each step. We may choose $k = k_\eta$ to maximize this product, which leads to factor $\exp(\tilde{C}\gamma\sqrt{\eta}) =: g(\eta)$. Furthermore, by Proposition 1 after this time $t = 2\eta$, the evolution is asymptotically stable and a small perturbation of the identity.

Let now $\psi \in H^s(\mathbb{R})$ be given and consider the initial datum:

$$\omega_0 = \int_\eta \frac{1}{g(\eta)} \tilde{\psi}(\eta) \omega_0^{\eta, k_\eta}.$$

Then by the definition of $g(\eta)$ and the properties of the evolution, $\omega(t)$ will asymptotically to leading order be given by

$$\omega_\infty = e^{ix} \int_\eta \tilde{\psi}(\eta) e^{i\eta y} d\eta = e^{ix} \psi(y).$$

We can thus prescribe final data. We further observe that $\frac{1}{g(\eta)} \leq \exp(-C\sqrt{\eta})$ and thus $\omega_0 \in \mathcal{G}_{C, \frac{1}{2}}$. Our proof thus concludes by choosing $\psi \in H^s \setminus H^\sigma$ appropriately. \square

6. Discussion

In view of applications to the nonlinear dynamics we note that we have several competing (de)stabilizing effects, whose interaction makes this a very challenging problem:

- On the one hand the norm inflation mechanism of Section 5 and similarly discussed in [5] shows that the vorticity may exhibit instability unless it is initially small in a sufficiently strong Gevrey class.
- On the other hand resonant times are well-separated and for any given η there are no resonances after time η . In particular, in the present problem fixing any finite radius R as a frequency cut-off $\mathcal{F}^{-1}\chi_{B_R}\mathcal{F}\omega$ and its corresponding velocity field do converge irrespective of the regularity of the initial data.
- Any instability will thus have to sustain a sequence of infinitely many separate echo chains for a sequence of times tending to infinity to ensure that the flow is not asymptotically stable after all (see Sections 2 and 5.3).
- While such a sequence of echo chains can be constructed in our model due to its decoupling structure (see Section 5.3), in the full nonlinear problem the conservation of enstrophy limits the possible relative growth. That is the conservation law imposes a hard ceiling for instability in that the L^2 energy remains bounded uniformly. Hence, it might be that the linear(!) instability mechanism of echoes is only applicable for finite times, after which the enstrophy limits further growth and the asymptotic stability of Section 2 takes over.

Here the modified scattering results of Section 5.3 and [16] provide a first indication that this may result in non-trivial but asymptotically stable behavior.

We further stress that, while stability of the linearized problem in Sobolev [12, 14, 15] and Gevrey spaces [8] is fundamental to attack the nonlinear problem, this article shows that it is further essential to understand the linearization around non-shear low frequency perturbations, which appear naturally in the nonlinear problem.

In the present work we have, for simplicity of calculation and presentation, considered a single-mode perturbation $c \cos(x)$. We expect analogous results to also hold for more general finite sums of small frequency perturbations, though involving quite involved calculations. Our choice $\cos(x)$ is motivated by its simplicity and the fact that, as an eigenfunction of the Laplacian, it is a stationary solution of the Euler equations in Lagrangian coordinates with respect to Couette flow.

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A. Special Functions and a Proof of Theorem 4

As mentioned in Section 4, the problem

$$\partial_t^2 u + \frac{c^2}{\xi - 2 + t^2} u = 0, t \in (-1, 1)$$

allows for an explicit solution in terms of special functions, which we discuss in the following. Since these explicit solutions are not available for perturbed coefficients or stable under further perturbation, in Sections 4 and 5 we have instead opted to study (and control) the evolution on three intervals I_1, I_2, I_3 and, in particular, establish the growth and decay of u and $\partial_t u$ on each interval. The interaction of these growths and decays then results in the same power law growth as we observe in the following for the special functions:

Proof of Theorem 4. Denoting $\xi = \frac{k^2}{\eta}$ for simplicity of notation, one may obtain the following explicit solution for boundary conditions $u(-1) = 1, \partial_t u(-1) = 0$ (e.g. using Mathematica)

$$\begin{aligned} & \left(-3c^2 t {}_2F_1 \left(\frac{3}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{3}{4}; \frac{3}{2}; -\frac{1}{\xi^2} \right) \right. \\ & 2F_1 \left(\frac{1}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{1}{4}; \frac{3}{2}; -\frac{t^2}{\xi^2} \right) \\ & - c^2 {}_2F_1 \left(\frac{5}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{5}{4}; \frac{5}{2}; -\frac{1}{\xi^2} \right) \\ & 2F_1 \left(-\frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}, \frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}; \frac{1}{2}; -\frac{t^2}{\xi^2} \right) \\ & + 3\xi^2 {}_2F_1 \left(\frac{1}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{1}{4}; \frac{3}{2}; -\frac{1}{\xi^2} \right) \\ & \left. 2F_1 \left(-\frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}, \frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}; \frac{1}{2}; -\frac{t^2}{\xi^2} \right) \right) \end{aligned}$$

$$\begin{aligned}
 &/(3c^2 {}_2F_1\left(\frac{1}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{1}{4}; \frac{3}{2}; -\frac{1}{\xi^2}\right) \\
 &{}_2F_1\left(\frac{3}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{3}{4}; \frac{3}{2}; -\frac{1}{\xi^2}\right) \\
 &- c^2 {}_2F_1\left(-\frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}, \frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}; \frac{1}{2}; -\frac{1}{\xi^2}\right) \\
 &{}_2F_1\left(\frac{5}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{5}{4}; \frac{5}{2}; -\frac{1}{\xi^2}\right) \\
 &+ 3\xi^2 {}_2F_1\left(-\frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}, \frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}; \frac{1}{2}; -\frac{1}{\xi^2}\right) \\
 &{}_2F_1\left(\frac{1}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{1}{4}; \frac{3}{2}; -\frac{1}{\xi^2}\right) \tag{68}
 \end{aligned}$$

Here ${}_2F_1$ denotes a hypergeometric function. We may then evaluate this formula at $t = 1$ and use the series expansion of ${}_2F_1(a, b, c, x)$ at $x = \infty$:

$$\begin{aligned}
 &x^{-a-b}\left(x^b\left(\frac{(-1)^{-a}\Gamma(b-a)\Gamma(c)}{\Gamma(b)\Gamma(c-a)} + \frac{(-1)^{-a}a(a-c+1)\Gamma(b-a)\Gamma(c)}{(a-b+1)\Gamma(b)\Gamma(c-a)x} + O\left(\left(\frac{1}{x}\right)^2\right)\right)\right) \\
 &+ x^a\left(\frac{(-1)^{-b}\Gamma(a-b)\Gamma(c)}{\Gamma(a)\Gamma(c-b)} + \frac{(-1)^{-b}b(b-c+1)\Gamma(a-b)\Gamma(c)}{(-a+b+1)\Gamma(a)\Gamma(c-b)x} + O\left(\left(\frac{1}{x}\right)^2\right)\right). \tag{69} \\
 &\tag{70}
 \end{aligned}$$

This then, for example, yields that, for $u(-1) = 1, u'(-1) = 0,$

$$\begin{aligned}
 u(1) = &(-3c^2 {}_2F_1\left(\frac{1}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{1}{4}; \frac{3}{2}; -\frac{1}{r^2}\right) \\
 &{}_2F_1\left(\frac{3}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{3}{4}; \frac{3}{2}; -\frac{1}{r^2}\right) \\
 &- c^2 {}_2F_1\left(-\frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}, \frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}; \frac{1}{2}; -\frac{1}{r^2}\right) \\
 &{}_2F_1\left(\frac{5}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{5}{4}; \frac{5}{2}; -\frac{1}{r^2}\right) \\
 &+ 3r^2 {}_2F_1\left(-\frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}, \frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}; \frac{1}{2}; -\frac{1}{r^2}\right) \\
 &{}_2F_1\left(\frac{1}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{1}{4}; \frac{3}{2}; -\frac{1}{r^2}\right) \\
 &/(3c^2 {}_2F_1\left(\frac{1}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{1}{4}; \frac{3}{2}; -\frac{1}{r^2}\right) \\
 &{}_2F_1\left(\frac{3}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{3}{4}; \frac{3}{2}; -\frac{1}{r^2}\right) \\
 &- c^2 {}_2F_1\left(-\frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}, \frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}; \frac{1}{2}; -\frac{1}{r^2}\right) \\
 &{}_2F_1\left(\frac{5}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{5}{4}; \frac{5}{2}; -\frac{1}{r^2}\right)
 \end{aligned}$$

$$\begin{aligned}
 &+ 3r^2 {}_2F_1\left(-\frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}, \frac{1}{4}\sqrt{1-4c^2} - \frac{1}{4}; \frac{1}{2}; -\frac{1}{r^2}\right) \\
 &{}_2F_1\left(\frac{1}{4} - \frac{1}{4}\sqrt{1-4c^2}, \frac{1}{4}\sqrt{1-4c^2} + \frac{1}{4}; \frac{3}{2}; -\frac{1}{r^2}\right)
 \end{aligned}$$

can be approximated as

$$\frac{r^{-\sqrt{1-4c^2}} r^2 \left(-\frac{3\left(2^{\sqrt{1-4c^2}-2} c^2 (\sqrt{1-4c^2}+1) \Gamma\left(\frac{1}{2}\sqrt{1-4c^2}\right)^2\right)}{\Gamma\left(\frac{1}{2}(\sqrt{1-4c^2}+3)\right)^2} \right) + o}{3r^2 + o},$$

where we denoted $r = \xi^{-1} \ll 1$ and o refers to terms decaying to higher order in r . We note in particular that the powers r^2 cancel. Approximating the value of the Γ functions by their value in $c = 0$, we thus obtain

$$\xi^\gamma \pi c^2.$$

Similar calculations for $\partial_t u$ and other initial data lead to the following coefficient matrix:

$$\begin{pmatrix} u(1) \\ \partial_t u(1) \end{pmatrix} \approx \xi^\gamma \begin{pmatrix} \pi c^2 & -\pi c^2 \\ \pi c^2 & -\pi c^2 \end{pmatrix} \begin{pmatrix} u(-1) \\ \partial_t u(-1) \end{pmatrix}$$

□

While the above calculations and asymptotics are explicit, they are also rather opaque. The splitting of the evolution into intervals I_1, I_2, I_3 studied in Lemmas 3, 4 and Proposition 3 instead provides a much clearer view of the underlying mechanism and yields the same leading asymptotics in terms of ξ .

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