# Existence and uniqueness of solutions of nonlinear evolution systems of n-th order partial differential equations in the complex plane 

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#### Abstract

We prove existence and uniqueness results for quasilinear systems of partial differential equations of the form $$
\partial_{t} \mathbf{h}+\left(-\partial_{y}\right)^{n} \mathbf{h}=\mathbf{g}_{2}\left(y, t,\left\{\partial_{y}^{j} \mathbf{h}\right\}_{j=0}^{n-1}\right) \partial_{y}^{n} \mathbf{h}+\mathbf{g}_{1}\left(y, t,\left\{\partial_{y}^{j} \mathbf{h}\right\}_{j=0}^{n-1}\right)+\mathbf{r}(y, t) ; \mathbf{h}(y, 0)=\mathbf{h}_{I}(y)
$$ for sufficiently large $y$ in a sector in the complex plane, under certain analyticity and decay conditions on $\mathbf{h}, \mathbf{h}_{I}, \mathbf{g}_{1}, \mathbf{g}_{2}$ and $\mathbf{r}$ such that the nonlinearity is formally small as $y \rightarrow \infty$. Due to the type of nonlinearity in the highest derivatives and the divergence of formal series solutions owed to the singularity of the system at infinity we needed to use new techniques (based on Borel-Laplace duality) to control the perturbative terms. Our methods of proof and norms used are of a constructive nature. The results given can be used to justify the formal asymptotic expansion of solutions in an appropriate sector and conversely to reconstruct actual solutions from formal series ones.


## 1 Introduction

The theory of partial differential equations, when one or more of the independent variables are in the complex plane is not very developed. The classic Cauchy-Kowalevski (C-K) theorem holds for a system of first-order equations (or those equivalent to it) when the quasi-linear equations have analytic coefficients and analytic initial data is specified on an analytic but non-characteristic curve. Then, the C-K theorem guarantees the local existence and uniqueness of analytic solutions. As is well known, its proof relies on the convergence of local power series expansions and, without the given hypotheses, the power series may have zero radius of convergence and the C-K method does not yield solutions.

[^0]More recently, Sammartino and Caflisch ([18], [19]) proved the existence of nonlinear Prandtl boundary layer equations for analytic initial data in a half-plane. This work involved inversion of the heat operator $\partial_{t}-\partial_{Y Y}$ and using the abstract Cauchy-Kowalewski theorem for the resulting integral equation.

While this method is likely to be generalizable to certain higher-order partial differential equations, it appears unsuitable for problems where the highest derivative terms appear in a non-linear manner. These terms cannot be controlled by inversion of a linear operator and estimates of the kernel, as used by Sammartino and Caflisch.

Nonlinearity of special forms in the highest derivative was considered in a general setting, encompassing nonlinear wave equations in the real domain, in a series of profound studies by Klainerman [10], [11], Shatah [20], Klainerman and Ponce [12], Ponce [17], Klainerman and Selberg [13], Shatah and Struwe [21] and others. The type of nonlinearity allowed in our paper and the complex plane setting do not appear to allow for an adaptation of those techniques, and we use a completely different and constructive method. The method is based on recently developed generalized Borel summation techniques [8], [5], [7]. Our results are presently restricted to one spatial dimension.

We show existence and uniqueness of solutions $\mathbf{h}(y, t)$ to the initial value problem for a general class of quasilinear system of partial differential equations of the form:
$\partial_{t} \mathbf{h}+\left(-\partial_{y}\right)^{n} \mathbf{h}=\mathbf{g}_{2}\left(y, t,\left\{\partial_{y}^{j} \mathbf{h}\right\}_{j=0}^{n-1}\right) \partial_{y}^{n} \mathbf{h}+\mathbf{g}_{1}\left(y, t,\left\{\partial_{y}^{j} \mathbf{h}\right\}_{j=0}^{n-1}\right)+\mathbf{r}_{1}(y, t) \quad$ with $\mathbf{h}(y, 0)=\mathbf{h}_{I}(y)$
where $\partial_{y}^{j}$ denotes the $j$-th derivative with respect to $y \in \mathbb{C}$. The function $\mathbf{h}$ takes values in $\mathbb{C}^{m_{1}}$. Suitable regularity and decay conditions are imposed for large $y$ in a sector in the complex plane (§2) and our results will hold in such a sector. These conditions make the terms in $\mathbf{h}$ on the right side of (1) formally small for large $y$. As will be discussed, existence cannot be expected, in general, to hold outside a specific sector.

By the transformation $y=z^{-1}$ it is seen that the problem is that of existence and uniqueness of solutions in a (sectorial) neighborhood of a point where the linear part of the principal symbol of a partial differential operator has a high order pole. Formal solutions as power series in $z$ are expected to have zero radius of convergence, cf. also $\S 6.2$ (iv). The setting (1) is the form usually encountered in applications.

As exemplified in $\S 6.2$ (v), classical spaces approaches to the problem may run into fundamental difficulties. For problems of this type, the essence of the methodology we have introduced recently in [7] in special cases of (1) has been to use Borel-Laplace duality to regularize the problem which is recast as an integral equation in the Borel plane. This regularization, akin to Borel summation, ${ }^{1}$ is instrumental in controlling the solution. The choice of appropriate Banach space after the Borel transform proves to be crucial, and after this choice, the contraction mapping argument itself follows from a sequence of relatively straightforward estimates in convolution Banach algebras. Borel summation methods have been also used recently in the context of the heat equation by Lutz, Miyake and Schäfke [14]. We illustrate in the Appendix, $\S 6.2$, the regularizing role of the Borel transform and discuss why it is instrumental in showing existence and

[^1]uniqueness. The method we use is constructive, in the sense that it permits recovering of actual solutions from formal ones presented as classically divergent power series in inverse powers of $y$, see also the notes in $\S 2$. The study is done in a sector in the complex plane whose width is crucial to existence and uniqueness of solutions with prescribed decay.

Our previous results [7] were limited to a class of partial differential equations that are first order in time, $t$, and third order in space, $y$. Further, those results, motivated by a set of applications, were restricted to scalar dependent variable with no nonlinearity in its derivatives.

Among the concrete equations amenable through rather straightforward transformations to the setting [7] and thus to the present more general one are the KdV equation, the equations $H_{t}=H^{3} H_{x x x}$ and $H_{t}+H_{x}=H^{3} H_{x x x}-H^{3} / 2$, both arising in Hele-Shaw dynamics, the equation $H_{t}=H^{1 / 3} H_{x x x}$ relevant to dendritic crystal growth, and many others. The last three of these PDEs were treated in detail in [7].

In the present paper we are generalizing the results of [7] to arbitrary order in the spatial variable. Further, the dependent function is allowed to be a vector $\mathbf{f}(y, t)$. The nonlinearity is that of a general quasi-linear equation. Indeed, it will be obvious that the result given here also generalizes easily to the case when the left side of (1) is replaced by $\partial_{t} \hat{\mathbf{f}}-A \partial_{y}^{n} \hat{\mathbf{f}}$, where $A$ is a constant matrix with positive eigenvalues, though we leave out this extra generality for the sake of relative simplicity in presentation.

Apart from the mathematical interest of understanding the question of existence and uniqueness of solutions to nonlinear PDEs in the complex plane near singular points ( $y=\infty$ in our context), sectorial existence of solutions to higher order nonlinear PDEs is important in many applications.

For instance, there are nonlinear PDEs for which the initial value problem, in the absence of a regularization is relatively simple; yet ill-posed in the sense of Hadamard for any Sobolev norm on the real domain. However, the analytically continued equations into the complex spatial domain are well-posed, even without a regularization term. There have been quite a few complex domain studies involving idealized equations modeling physical phenomena (see for instance, [15], [16], [1], [2], [3] and [4]) that follow Garabedian's [9] realization that an ill-posed elliptic initial value problem in the real spatial domain may become well-posed in the complex domain. In a particular physical context, it was suggested [23] that complex domain studies would be useful in understanding small regularization for some class of initial conditions. Formal and numerical computations show the usefulness of this approach in predicting singular effects [22]. However, many of the results for small but nonzero regularization were formal and relied fundamentally on the existence and uniqueness of analytic solutions to certain higher order nonlinear partial differential equations in a sector in the complex plane, with imposed far-field matching conditions. In [7] we addressed rigorously the existence questions in some third-order nonlinear PDEs.

In a more general context, one can expect that whenever regularization appears in the form of a small coefficient multiplying the highest spatial derivative, the resulting asymptotic equation in the neighborhood of initial complex singularities will satisfy higher order (not necessarily third order) nonlinear partial differential equation with sectorial far-field matching condition in the complex plane of the type discussed here. An example of such an application is the analysis of the local behavior of solutions to the well known Kuramato-Shivashinski equation: $u_{t}+u u_{x}+u_{x x}+\nu u_{x x x x}=0$ for small nonzero $\nu$, near a complex-singularity. The same type of existence questions is relevant in the small time asymptotics of higher order nonlinear PDEs near initial singularities, or the large $x$ behavior of solutions. Detailed discussions of other specific
equations, their sectorial regularity as well as singularity analysis, will be the object of a different paper.

## 2 Problem statement and main result

We study equation (1) where the inhomogeneous term $\mathbf{r}_{1}(y, t)$ is analytic in $y$ the domain

$$
\begin{equation*}
\mathcal{D}_{\rho_{0}}=\left\{(y, t) \in \mathbb{C} \times \mathbb{R}: \text { arg } y \in\left(-\frac{\pi}{2}-\frac{\pi}{2 n}, \frac{\pi}{2}+\frac{\pi}{2 n}\right),|y|>\rho_{0}>0,0 \leq t \leq T\right\} \tag{2}
\end{equation*}
$$

The restrictions on $\mathbf{g}_{1}, \mathbf{g}_{2}, \mathbf{r}_{1}$ and $\mathbf{h}_{\mathbf{I}}$ are better expressed once we transform the equation to a more convenient form, as is done shortly.

By taking derivatives of (1) with respect to $y$, $i$-times, $i$ ranging from 1 to $n-1$, it is possible to consider the extended $m_{1} \times n$ system of equations for $\mathbf{h}$ and its first $n-1$-derivatives. This system of vector equations is of the form (see Appendix for further details)

$$
\begin{equation*}
\partial_{t} \mathbf{f}+\left(-\partial_{y}\right)^{n} \mathbf{f}=\sum_{\mathbf{q} \succeq 0}^{\prime} \mathbf{b}_{\mathbf{q}}(y, t, \mathbf{f}) \prod_{l=1}^{m} \prod_{j=1}^{n}\left(\partial_{y}^{j} f_{l}\right)^{q_{l, j}}+\mathbf{r}(y, t) ; \quad \text { with } \mathbf{f}(y, 0)=\mathbf{f}_{I}(y) \tag{3}
\end{equation*}
$$

where $\sum^{\prime}$ means the sum over the multiindices $\mathbf{q}$ with

$$
\begin{equation*}
\sum_{l=1}^{m} \sum_{j=1}^{n} j q_{l, j} \leq n \tag{4}
\end{equation*}
$$

In (3), $\mathbf{f}$ is an $m=m_{1} \times n$ dimensional vector, $\mathbf{q}=\left\{q_{l, j}\right\}_{j=1, l=1}^{n, m}$ is a vector of integers and the notation $\mathbf{q} \succeq 0$ means $q_{l, j} \geq 0$ for all $l, j$. The inequality (4) implies in particular that none of the $q_{l, j}$ can exceed $n$ and that the summation on $\mathbf{q}$ involves only finitely many terms. The fact that (4) can always be ensured leads to important simplifications in the proofs. We denote

$$
\langle\mathbf{q}\rangle=\sum_{(l, j) \preceq(m, n)} q_{l, j}
$$

We assume that in $\mathcal{D}_{\rho_{0}}$ there exist constants $\alpha_{r} \geq 1$ (see also $\S 5.2$ ) and $A_{r}(T)$, with $\alpha_{r}$ independent of $T$ such that

$$
\begin{equation*}
\left|y^{\alpha_{r}} \mathbf{r}(y, t)\right|<A_{r}(T) \tag{5}
\end{equation*}
$$

In this paper the absolute value $|\cdot|$ of a vector is the max vector norm. Additionally, we require that $\mathbf{b}_{\mathbf{q}}$ is analytic in $\mathbf{f}$ and in its convergent representation

$$
\begin{equation*}
\mathbf{b}_{\mathbf{q}}(y, t ; \mathbf{f})=\sum_{\mathbf{k} \succeq 0} \mathbf{b}_{\mathbf{q}, \mathbf{k}}(y, t) \mathbf{f}^{\mathbf{k}} \text { where } \mathbf{k}=\left(k_{1}, k_{2}, \ldots, k_{m}\right) \text {, and } \mathbf{f}^{\mathbf{k}}=\prod_{i=1}^{m} f_{i}^{k_{i}} \tag{6}
\end{equation*}
$$

the functions $\mathbf{b}_{\mathbf{q}, \mathbf{k}}$ are analytic in $y$ in $\mathcal{D}_{\rho_{0}}$. We also take, without any loss of generality, $\mathbf{b}_{\mathbf{0}, \mathbf{0}}=\mathbf{0}$, altering accordingly (if needed) $\mathbf{r}(y, t)$.

We assume that in this domain, there exist positive constants $\beta, \alpha_{\mathbf{q}}$, and $A_{b}$, independent of $\mathbf{q}$ and $\mathbf{k}$ (with $\beta$ and $\alpha_{\mathbf{q}}$ independent of $T$ as well), such that

$$
\begin{equation*}
\left|y^{\alpha_{\mathbf{q}}+\langle\mathbf{k}\rangle \beta} \mathbf{b}_{\mathbf{q}, \mathbf{k}}\right|<A_{b}(T), \text { where } \quad\langle\mathbf{k}\rangle \equiv k_{1}+k_{2}+\cdots+k_{m} \tag{7}
\end{equation*}
$$

The series (6) converges in the domain $\mathcal{D}_{\phi, \rho}$, defined as (cf. also (2) and (7))

$$
\begin{equation*}
\mathcal{D}_{\phi, \rho}=\left\{(y, t): \arg y \in\left(-\frac{\pi}{2}-\phi, \frac{\pi}{2}+\phi\right),|y|>\rho>\rho_{0}, \text { where } 0<\phi<\frac{\pi}{2 n}, 0 \leq t \leq T\right\} \tag{8}
\end{equation*}
$$

if

$$
\begin{equation*}
|\mathbf{f}|<\rho^{\beta} \tag{9}
\end{equation*}
$$

Condition 1 The solution $\mathbf{f}(\cdot, t)$ sought for (3) is required to be analytic in $\mathcal{D}_{\phi, \rho ; y}$, where $\mathcal{D}_{\phi, \rho ; y}$ is the projection of $\mathcal{D}_{\phi, \rho}$ on the first component, for some $\rho>0$ (to be determined later). In the same domain, the solution and the initial condition $\mathbf{f}_{I}(y)$ must satisfy the conditions (cf. (5))

$$
\begin{equation*}
|y|^{\alpha_{r}}|\mathbf{f}(y, t)|<A_{f}(T) ; \quad|y|^{\alpha_{r}}\left|\mathbf{f}_{I}(y, t)\right|<A_{f}(T) \tag{10}
\end{equation*}
$$

for some $A_{f}(T)>0$ and $(y, t) \in \mathcal{D}_{\phi, \rho}$.
It is clear that for large $y$ such a solution $\mathbf{f}$ will indeed satisfy (9), the condition for the convergence of the infinite series in (3).

The general theorem proved in this paper is the following.
Theorem 2 For any $T>0$ and $0<\phi<\pi /(2 n)$, there exists $\tilde{\rho}$ such that the partial differential equation (3) has a unique solution $\mathbf{f}$ that is analytic in $y$ and is $O\left(y^{-1}\right)$ as $y \rightarrow \infty$ for $(y, t) \in$ $\mathcal{D}_{\tilde{\rho}, \phi}$. Furthermore, this solution satisfies $\mathbf{f}=O\left(y^{-\alpha_{r}}\right)$ as $y \rightarrow \infty$ in $\mathcal{D}_{\tilde{\rho}, \phi ; y}$.

Notes. 1. Existence necessitates the sector to be not too large; if the technique is used for reconstruction of actual solutions from formal ones as explained in 3. below, the sector must be large enough to ensure uniqueness; the width provided here ensures both.
2. The proof is made delicate by the presence of perturbation terms involving the highest derivative and the fact that, in the relevant limit $y \rightarrow \infty$ the equation is singular; the nature of this problem is sketched in $\S 6.2$.
3. As discussed earlier in [7] for special examples, this decaying behavior of the solution $\mathbf{f}$ is valid inside the specified sector and outside it one can expect infinitely many singularities with an accumulation point at infinity.
4. In $\S 5$ it is shown how our result and technique can be used to justify an asymptotic power series behavior of the solution, or to recover actual solutions from formal series expansions.

## 3 Inverse Laplace transform and equivalent integral equation

The inverse Laplace transform (ILT) $\mathbf{G}(p, t)$ of a function $\mathbf{g}(y, t)$ analytic in $y$ in $\mathcal{D}_{\phi, \rho ; y}$ (see Condition 1) and vanishing algebraically as $|y| \rightarrow \infty$ is given by:

$$
\begin{equation*}
\mathbf{G}(p, t)=\left[\mathcal{L}^{-1}\{\mathbf{g}\}\right](p, t) \equiv \frac{1}{2 \pi i} \int_{\mathcal{C}_{D}} e^{p y} \mathbf{g}(y, t) d y \tag{11}
\end{equation*}
$$

where $\mathcal{C}_{D}$ is a contour as in Fig. 1 (modulo homotopies), entirely within the domain $\mathcal{D}_{\phi, \rho ; y}$ and $p$ is restricted to the domain $\mathcal{S}_{\phi}$ where convergence of the integral is ensured, where

$$
\begin{equation*}
\mathcal{S}_{\phi} \equiv\{p: \arg p \in(-\phi, \phi), 0<|p|<\infty\} \tag{12}
\end{equation*}
$$

If $\mathbf{g}(y, t)=\mathbf{C} y^{-\alpha}$ for $\alpha>0$, then $\mathbf{G}(p, t)=\mathbf{C} p^{\alpha-1} / \Gamma(\alpha)$. From the following Lemma, it is clear that the same kind of behavior for the $\operatorname{ILT} \mathbf{G}(p, t)$ holds for small $p$ in $\mathcal{S}_{\phi}$, if $\mathbf{g}$ is $O\left(y^{-\alpha}\right)$ for large $y$.

Lemma 3 If $\mathbf{g}(y, t)$ is analytic in $y$ in $\mathcal{D}_{\phi, \rho ; y}$, and satisfies

$$
\begin{equation*}
\left|y^{\alpha} \mathbf{g}(y, t)\right|<A(T) \tag{13}
\end{equation*}
$$

for $\alpha \geq \alpha_{0}>0$, then for any $\delta \in(0, \phi)$ the ILT $\mathbf{G}=\mathcal{L}^{-1} \mathbf{g}$ exists in $\mathcal{S}_{\phi-\delta}$ and satisfies

$$
\begin{equation*}
|\mathbf{G}(p, t)|<C \frac{A(T)}{\Gamma(\alpha)}|p|^{\alpha-1} e^{2|p| \rho} \tag{14}
\end{equation*}
$$

for some $C=C\left(\delta, \alpha_{0}\right)$.
Proof. The proof is similar to that of Lemma 3.1 in [7]. We first consider the case when $2 \geq \alpha \geq \alpha_{0}$. Let $C_{\rho_{1}}$ be the contour $C_{D}$ in Fig. 1 that passes through the point $\rho_{1}+|p|^{-1}$, and given by $s=\rho_{1}+|p|^{-1}+\operatorname{ir} \exp (i \phi \operatorname{signum}(r))$ with $r \in(-\infty, \infty)$. Choosing $2 \rho>\rho_{1}>(2 / \sqrt{3}) \rho$, we have $|s|>\rho$ along the contour and therefore, with $\arg (p)=\theta \in(-\phi+\delta, \phi-\delta)$,

$$
|\mathbf{g}(s, t)|<A(T)|s|^{-\alpha} \quad \text { and } \quad\left|e^{s p}\right| \leq e^{\rho_{1}|p|+1} e^{-|r||p| \sin |\phi-\theta|}
$$

Thus

$$
\begin{align*}
& \left|\int_{C_{\rho_{1}}} e^{s p} \mathbf{g}(s, t) d s\right| \leq 2 A(T) e^{\rho_{1}|p|+1} \int_{0}^{\infty}\left|\rho_{1}+|p|^{-1}+i r e^{i \phi}\right|^{-\alpha} e^{-|p| r \sin \delta} d r \\
& \quad \leq \tilde{K} A(T) e^{\rho_{1}|p|}\left|\rho_{1}+|p|^{-1}\right|^{-\alpha} \int_{0}^{\infty} e^{-|p| r \sin \delta} d r \leq K \delta^{-1}|p|^{\alpha-1} e^{2 \rho|p|} \tag{15}
\end{align*}
$$

where $\tilde{K}$ and $K$ are constants independent of any parameter. Thus, the Lemma follows for $2 \geq \alpha \geq \alpha_{0}$, if we note that $\Gamma(\alpha)$ is bounded in this range of $\alpha$, with the bound only depending on $\alpha_{0}$.

For $\alpha>2$, there exists an integer $k>0$ so that $\alpha-k \in(1,2]$. Taking

$$
(k-1)!\mathbf{h}(y, t)=\int_{\infty}^{y} \mathbf{g}(z, t)(y-z)^{k-1} d z
$$

(clearly $\mathbf{h}$ is analytic in $y$, in $\mathcal{D}_{\phi, \rho}$ and $\mathbf{h}^{(k)}(y, t)=\mathbf{g}(y, t)$ ), we get

$$
\mathbf{h}(y, t)=\frac{(-y)^{k}}{(k-1)!} \int_{1}^{\infty} \mathbf{g}(y p, t)(p-1)^{k-1} d p=\frac{(-1)^{k} y^{k-\alpha}}{(k-1)!} \int_{1}^{\infty} \mathbf{A}(y p, t) p^{-\alpha}(p-1)^{k-1} d p
$$

with $|\mathbf{A}(y p, t)|<A(T)$, whence

$$
|\mathbf{h}(y, t)|<\frac{A(T) \Gamma(\alpha-k)}{|y|^{\alpha-k} \Gamma(\alpha)}
$$

From what has been already proved, with $\alpha-k$ playing the role of $\alpha$,

$$
\left|\mathcal{L}^{-1}\{\mathbf{h}\}(p, t)\right|<C(\delta) \frac{A(T)}{\Gamma(\alpha)}|p|^{\alpha-k-1} e^{2|p| \rho}
$$

Since $\mathbf{G}(p, t)=(-1)^{k} p^{k} \mathcal{L}^{-1}\{\mathbf{h}\}(p, t)$, by multiplying the above equation by $|p|^{k}$, the Lemma follows for $\alpha>2$ as well.
Comment 1: The constant $2 \rho$ in the exponential bound can be lowered to anything exceeding $\rho$, but (14) suffices for our purposes.

Comment 2: Corollary 4 below implies that for any $p \in \mathcal{S}_{\phi}$, the ILT exists for the functions $\mathbf{r}(y, t), \mathbf{b}_{\mathbf{q}, \mathbf{k}}(y, t)$, as well as for the solution $\mathbf{f}(y, t)$, whose existence is shown in the sequel.
Comment 3: Conversely, if $\mathbf{G}(p, t)$ is any integrable function satisfying the exponential bound in (14), it is clear that the Laplace Transform along a ray

$$
\begin{equation*}
\mathcal{L}_{\theta} \mathbf{G} \equiv \int_{0}^{\infty e^{i \theta}} d p e^{-p y} \mathbf{G}(p, t) \tag{16}
\end{equation*}
$$

exists and defines an analytic function of $y$ in the half-plane $\Re\left[e^{i \theta} y\right]>2 \rho$ for $\theta \in(-\phi, \phi)$.
Comment 4: The next corollary shows that there exist bounds for $\mathbf{B}_{\mathbf{q}, \mathbf{k}}=\mathcal{L}^{-1}\left\{\mathbf{b}_{\mathbf{q}, \mathbf{k}}\right\}$ and $\mathbf{R}=\mathcal{L}^{-1}\{\mathbf{r}\}$ independent of $\arg p$ in $\mathcal{S}_{\phi}$, because of the assumed analyticity and decay properties in the region $\mathcal{D}_{\rho_{0}}$, which contains $\mathcal{D}_{\phi, \rho}$.

Corollary 4 The ILT of the coefficient functions $\mathbf{b}_{\mathbf{q}, \mathbf{k}}$ ( $c f$. (6)) and the inhomogeneous term $\mathbf{r}(y, t)$ satisfy the following upper bounds for any $p \in \mathcal{S}_{\phi}$

$$
\begin{gather*}
\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}(p, t)\right|<\frac{C_{1}\left(\phi, \alpha_{\mathbf{q}}\right)}{\Gamma\left(\alpha_{\mathbf{q}}+\beta|\mathbf{k}|\right)} A_{b}(T)|p|^{\beta|\mathbf{k}|+\alpha_{\mathbf{q}}-1} e^{2 \rho_{0}|p|}  \tag{17}\\
|\mathbf{R}(p, t)|<\frac{C_{2}(\phi)}{\Gamma\left(\alpha_{r}\right)} A_{r}(T)|p|^{\alpha_{r}-1} e^{2 \rho_{0}|p|} \tag{18}
\end{gather*}
$$

Proof. The proof is similar to that of Corollary 3.2 in[7]. From the conditions assumed we see that $\mathbf{b}_{\mathbf{q}, \mathbf{k}}$ is analytic in $y \in \mathcal{D}_{\phi_{1}, \rho_{0} ; y}$ for any $\phi_{1}$ satisfying $(2 n)^{-1} \pi>\phi_{1}>\phi>0$. So Lemma 3 can be applied for $\mathbf{g}(y, t)=\mathbf{b}_{\mathbf{q}, \mathbf{k}}$, with $\phi_{1}=\phi+\left((2 n)^{-1} \pi-\phi\right) / 2$ replacing $\phi$, and with $\delta$ replaced by $\phi_{1}-\phi=\left((2 n)^{-1} \pi-\phi\right) / 2$, and the same applies to $\mathbf{R}(p, t)$, leading to (17) and (18). In the latter case, since $\alpha_{r} \geq 1, \alpha_{0}$ in Lemma 3 can be chosen to be 1 . Thus, one can choose $C_{2}$ to be independent of $\alpha_{r}$, as indicated in (18).

The formal inverse Laplace transform (Borel transform) of (3) with respect to $y$ is (see also (6))

$$
\begin{equation*}
\partial_{t} \mathbf{F}+p^{n} \mathbf{F}=\sum_{\mathbf{q} \succeq 0}^{\prime} \sum_{\mathbf{k} \succeq 0} \mathbf{B}_{\mathbf{q}, \mathbf{k}} * \mathbf{F}^{* \mathbf{k}} * \prod_{l=1}^{*} \prod_{j=1}^{*}\left((-p)^{j} F_{l}\right)^{* q_{l, j}}+\mathbf{R}(p, t) \tag{19}
\end{equation*}
$$

where the symbol $*$ stands for convolution

$$
\begin{equation*}
(f * g)(p):=\int_{0}^{p} f(s) g(p-s) d s \tag{20}
\end{equation*}
$$

${ }^{*} \prod$ is a convolution product (see also [5]) and $\mathbf{F}=\mathcal{L}^{-1} \mathbf{f}$. After inverting the differential operator on the left side of (19) with respect to $t$, we obtain the integral equation

$$
\begin{align*}
& \mathbf{F}(p, t)=\mathcal{N}(\mathbf{F}) \equiv \\
& \int_{0}^{t} e^{-p^{n}(t-\tau)} \sum_{\mathbf{q} \succeq 0}^{\prime} \sum_{\mathbf{k} \succeq 0} \mathbf{B}_{\mathbf{q}, \mathbf{k}}(p, \tau) * \mathbf{F}^{* \mathbf{k}}(p, \tau) *{ }^{*} \prod_{l=1}^{m}{ }^{*} \prod_{j=1}^{n}\left((-p)^{j} F_{l}(p, \tau)\right)^{* q_{l, j}} d \tau+\mathbf{F}_{0}(p, t) \tag{21}
\end{align*}
$$

where

$$
\begin{equation*}
\mathbf{F}_{0}(p, t)=e^{-p^{n} t} \mathbf{F}_{I}(p)+\int_{0}^{t} e^{-p^{n}(t-\tau)} \mathbf{R}(p, \tau) d \tau \quad \text { and } \mathbf{F}_{I}=\mathcal{L}^{-1}\left\{\mathbf{f}_{I}\right\} \tag{22}
\end{equation*}
$$

Our strategy is to reduce the problem of existence and uniqueness of a solution of (3) to the problem of existence and uniqueness of a solution of (21), under appropriate conditions.

## 4 Solution to the integral equation (21)

To establish the existence and uniqueness of solutions to the integral equation, we need to introduce an appropriate function class for both the solution and the coefficient functions.
Definition 5 Denoting by $\overline{\mathcal{S}_{\phi}}$ the closure of $\mathcal{S}_{\phi}$ defined in (12), $\partial \mathcal{S}_{\phi}=\overline{\mathcal{S}_{\phi}} \backslash \mathcal{S}_{\phi}$ and $\mathcal{K}=\overline{\mathcal{S}_{\phi}} \times$ $[0, T]$, we define for $\nu>0$ (later to be taken appropriately large) the norm $\|\cdot\|_{\nu}$ as

$$
\begin{equation*}
\|\mathbf{G}\|_{\nu}=M_{0} \sup _{(p, t) \in \mathcal{K}}\left(1+|p|^{2}\right) e^{-\nu|p|}|\mathbf{G}(p, t)| \tag{23}
\end{equation*}
$$

where $M_{0}$ is a constant (approximately 3.76) defined as

$$
\begin{equation*}
M_{0}=\sup _{s \geq 0}\left\{\frac{2\left(1+s^{2}\right)\left(\ln \left(1+s^{2}\right)+s \arctan s\right)}{s\left(s^{2}+4\right)}\right\} \tag{24}
\end{equation*}
$$

Note: For fixed $\mathbf{F},\|\mathbf{F}\|_{\nu}$ is nonincreasing in $\nu$.

Definition 6 We now define the following class of functions:

$$
\mathcal{A}_{\phi}=\left\{\mathbf{F}: \mathbf{F}(\cdot, t) \text { analytic in } \mathcal{S}_{\phi} \text { and continuous in } \overline{\mathcal{S}_{\phi}} \text { for } t \in[0, T] \text { s.t. }\|\mathbf{F}\|_{\nu}<\infty\right\}
$$

It is clear that $\mathcal{A}_{\phi}$ forms a Banach space.

Comment 5: Note that given $\mathbf{G} \in \mathcal{A}_{\phi}, \mathbf{g}(y, t)=\mathcal{L}_{\theta}\{\mathbf{G}\}$ exists for appropriately chosen $\theta$ when $\rho$ is large enough so that $\rho \cos (\theta+\arg y)>\nu$, and that $\mathbf{g}(y, t)$ is analytic in $y$ and $|y \mathbf{g}(y, t)|$ bounded for $y$ on any fixed ray in $\mathcal{D}_{\phi, \rho ; y}$.

Lemma 7 For $\nu>4 \rho_{0}+\alpha_{r}, \mathbf{F}_{I}$ in (22) satisfies

$$
\left\|\mathbf{F}_{I}\right\|_{\nu}<C(\phi) A_{f_{I}}(\nu / 2)^{-\alpha_{n}+1}
$$

while $\mathbf{R}$ satisfies the inequality

$$
\|\mathbf{R}\|_{\nu}<C(\phi) A_{r}(T)(\nu / 2)^{-\alpha_{n}+1}
$$

and therefore

$$
\begin{equation*}
\left\|\mathbf{F}_{0}\right\|_{\nu}<C(\phi)\left(T A_{r}+A_{f_{I}}\right)(\nu / 2)^{-\alpha_{r}+1} \tag{25}
\end{equation*}
$$

Proof. This proof is similar to that of Lemma 4.4 in [7]. First note the bounds on $\mathbf{R}$ in Corollary 4. We also note that $\alpha_{r} \geq 1$ and that for $\nu>4 \rho_{0}+\alpha_{r}$ we have

$$
\sup _{p} \frac{|p|^{\alpha_{r} \pm 1}}{\Gamma\left(\alpha_{r}\right)} e^{-\left(\nu-2 \rho_{0}\right)|p|} \leq \frac{\left(\alpha_{r} \pm 1\right)^{\alpha_{r} \pm 1}}{\Gamma\left(\alpha_{r}\right)} e^{-\alpha_{r} \mp 1}\left(\nu-2 \rho_{0}\right)^{-\alpha_{r} \mp 1} \leq K \alpha_{r}^{1 / 2 \pm 1}(\nu / 2)^{-\alpha_{r} \mp 1}
$$

where $K$ is independent of any parameter. The latter inequality follows from Stirling's formula for $\Gamma\left(\alpha_{r}\right)$ for large $\alpha_{r}$.

Using the definition of the $\nu$-norm and the two equations above, the inequality for $\|\mathbf{R}\|_{\nu}$ follows. Since $\mathbf{f}_{I}(y)$ is required to satisfy the same bounds as $\mathbf{r}(y, t)$, a similar inequality holds for $\left\|\mathbf{F}_{I}\right\|_{\nu}$. Now, from the relation (22),

$$
\left|\mathbf{F}_{0}(p, t)\right|<\left|\mathbf{F}_{I}(p)\right|+T \sup _{0 \leq t \leq T}|R(p, t)|
$$

Therefore, (25) follows.
Comment 6: Not all Laplace-transformable analytic functions in $\mathcal{D}_{\phi, \rho ; y}$ belong to $\mathcal{A}_{\phi}$. In our assumptions, the coefficients need not be bounded near $p=0$ and hence do not belong in $\mathcal{A}_{\phi}$. It is then useful to introduce the following function class:

## Definition 8

$$
\mathcal{H} \equiv\left\{\mathbf{H}: \mathbf{H}(p, t) \text { analytic in } \mathcal{S}_{\phi},|\mathbf{H}(p, t)|<C|p|^{\alpha-1} e^{\rho|p|}\right\}
$$

for some positive constants $C$ and $\alpha$ and $\rho$ which may depend on $\mathbf{H}$.

Lemma 9 If $\mathbf{H} \in \mathcal{H}$ and $\mathbf{F} \in \mathcal{A}_{\phi}$, then for $\nu>\rho+1$, for any $j, \mathbf{H} * F_{j}$ belongs to $\mathcal{A}_{\phi}$, and satisfies the following inequality ${ }^{2}$ :

$$
\begin{equation*}
\left\|\mathbf{H} * F_{j}\right\| \leq\left\||\mathbf{H}| *\left|F_{j}\right|\right\|_{\nu} \leq C \Gamma(\alpha)(\nu-\rho)^{-\alpha}\|\mathbf{F}\|_{\nu} \tag{26}
\end{equation*}
$$

[^2]The proof is a vector adaptation of that of Lemma 4.6 in [7].
Proof. From the elementary properties of convolution, it is clear that $\mathbf{H} * F_{j}$ is analytic in $\mathcal{S}_{\phi}$ and is continuous on $\overline{\mathcal{S}_{\phi}}$. Let $\theta=\arg p$. We have

$$
\left|\mathbf{H} * F_{j}(p)\right| \leq||\mathbf{H}| *| F_{j}|(p)| \leq \int_{0}^{|p|}\left|\mathbf{H}\left(s e^{i \theta}\right)\right|\left|F_{j}\left(p-s e^{i \theta}\right)\right| d s
$$

But

$$
\left|\mathbf{H}\left(s e^{i \theta}\right)\right| \leq C s^{\alpha-1} e^{|s| \rho}
$$

and

$$
\begin{equation*}
\int_{0}^{|p|} s^{\alpha-1} e^{|s| \rho}\left|F_{j}\left(p-s e^{i \theta}\right)\right| d s \leq\left\|F_{j}\right\|_{\nu} e^{\nu|p|}|p|^{\alpha} \int_{0}^{1} \frac{s^{\alpha-1} e^{-(\nu-\rho)|p| s}}{M_{0}\left(1+|p|^{2}(1-s)^{2}\right)} d s \tag{27}
\end{equation*}
$$

If $|p|$ is large, noting that $\nu-\rho \geq 1$, we obtain from Watson's lemma,

$$
\begin{equation*}
\int_{0}^{|p|} s^{\alpha-1} e^{|s| \rho}\left|F_{j}\left(p-s e^{i \theta}\right)\right| d s \leq K \Gamma(\alpha)\left\|F_{j}\right\|_{\nu} \frac{e^{\nu|p|}}{M_{0}\left(1+|p|^{2}\right)}|\nu-\rho|^{-\alpha} \tag{28}
\end{equation*}
$$

Now, for any other $|p|$, we obtain from (27),

$$
\int_{0}^{|p|} s^{\alpha-1} e^{|s| \rho}\left|F_{j}\left(p-s e^{i \theta}\right)\right| d s \leq K|\nu-\rho|^{-\alpha}\left\|F_{j}\right\|_{\nu} \frac{e^{\nu|p|} \Gamma(\alpha)}{M_{0}}
$$

Thus (28) must hold in general as it subsumes the above relation when $|p|$ is not large. From this relation, (26) follows by applying the definition of $\|\cdot\|_{\nu}$.

Corollary 10 For $\mathbf{F} \in \mathcal{A}_{\phi}$, and $\nu>4 \rho_{0}+1$, we have $\mathbf{B}_{\mathbf{q}, \mathbf{k}} * F_{l} \in \mathcal{A}_{\phi}$ and

$$
\left\|\mathbf{B}_{\mathbf{q}, \mathbf{k}} * F_{l}\right\|_{\nu} \leq\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{F}|\right\|_{\nu} \leq K C_{1}\left(\phi, \alpha_{\mathbf{q}}\right)(\nu / 2)^{-\langle\mathbf{k}\rangle \beta-\alpha_{\mathbf{q}}} A_{b}(T)\|\mathbf{F}\|_{\nu}
$$

Proof. The proof follows simply by using Lemma 9 , with $\mathbf{H}$ replaced by $\mathbf{B}_{\mathbf{q}, \mathbf{k}}$ and using the relations in Corollary 4.

Lemma 11 For $\mathbf{F} \in \mathcal{A}_{\phi}$, with $\nu>4 \rho_{0}+1$, for any $j, l$,

$$
\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}} *\left(p^{j} F_{l}\right)\right| \leq \frac{K C_{1}|p|^{j} e^{\nu|p|} A_{b}(T)}{M_{0}\left(1+|p|^{2}\right)}\|\mathbf{F}\|_{\nu}\left(\frac{\nu}{2}\right)^{-\beta\langle\mathbf{k}\rangle-\alpha_{\mathbf{q}}}
$$

Proof.
From the definition (20), it readily follows that

$$
\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}} *\left(p^{j} F_{l}\right)\right| \leq|p|^{j}\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *\left|F_{l}\right|
$$

The rest follows from Corollary (10), and definition of $\|\cdot\|_{\nu}$.

Lemma 12 For $\mathbf{F}, \mathbf{G} \in \mathcal{A}_{\phi}$ and $j \geq 0$

$$
\begin{equation*}
\left|\left(p^{j} F_{l_{1}}\right) * G_{l_{2}}\right| \leq|p|^{j}| | \mathbf{F}|*| \mathbf{G}| | \tag{29}
\end{equation*}
$$

Proof. Let $p=|p| e^{i \theta}$. Then,

$$
\begin{equation*}
\left|\left(p^{j} F_{l_{1}}\right) * G_{l_{2}}\right|=\left|\int_{0}^{p} \tilde{s}^{j} F_{l_{1}}(\tilde{s}) G_{l_{2}}(p-\tilde{s}) d \tilde{s}\right| \leq|p|^{j} \int_{0}^{|p|} d s\left|\mathbf{F}\left(s e^{i \theta}\right)\right|\left|\mathbf{G}\left(p-s e^{i \theta}\right)\right| \tag{30}
\end{equation*}
$$

from which the lemma follows.

Corollary 13 If $\mathbf{F} \in \mathcal{A}_{\phi}$, then

$$
\begin{equation*}
\left|\prod_{l=1}^{m} \prod_{j=1}^{n}\left(p^{j} F_{l}\right)^{* q_{l, j}}\right| \leq\left.|p|^{\sum j q_{l, j}}\left|\prod_{l=1}^{m} \prod_{j=1}^{n}\right| \mathbf{F}\right|^{* q_{l, j}} \mid \tag{31}
\end{equation*}
$$

where $\sum j q_{l, j}$ extends over all $(l, j) \preceq(m, n)$.
Proof. This follows simply from repeated application of Lemma 12.

Lemma 14 For $\mathbf{F}, \mathbf{G} \in \mathcal{A}_{\phi}$,

$$
||\mathbf{F}| *| \mathbf{G}\left|\left\lvert\, \leq \frac{e^{\nu|p|}}{M_{0}\left(1+|p|^{2}\right)}\|\mathbf{F}\|_{\nu}\|\mathbf{G}\|_{\nu}\right.\right.
$$

Proof.

$$
\begin{equation*}
||\mathbf{F}| *| \mathbf{G}\left|\left|=\left|\int_{0}^{p}\right| \mathbf{F}(\tilde{s})\right|\right| \mathbf{G}(p-\tilde{s})|d \tilde{s}| \leq \int_{0}^{|p|} d s\left|\mathbf{F}\left(s e^{i \theta}\right)\right|\left|\mathbf{G}\left(p-s e^{i \theta}\right)\right| \tag{32}
\end{equation*}
$$

Using the definition of $\|\cdot\|_{\nu}$, the above expression is bounded by

$$
\frac{e^{\nu|p|}}{M_{0}^{2}}\|\mathbf{F}\|_{\nu}\|\mathbf{G}\|_{\nu} \int_{0}^{|p|} \frac{d s}{\left(1+s^{2}\right)\left[1+(|p|-s)^{2}\right]} \leq \frac{|p|^{j} e^{\nu|p|}}{M_{0}\left(1+|p|^{2}\right)}\|\mathbf{F}\|_{\nu}\|\mathbf{G}\|_{\nu}
$$

The last inequality follows from the definition (24) of $M_{0}$ since

$$
\int_{0}^{|p|} \frac{1}{\left(1+s^{2}\right)\left[1+(|p|-s)^{2}\right]}=2 \frac{\ln \left(|p|^{2}+1\right)+|p| \tan ^{-1}|p|}{|p|\left(|p|^{2}+4\right)}
$$

Corollary 15 For $\mathbf{F}, \mathbf{G} \in \mathcal{A}_{\phi}$, then

$$
\||\mathbf{F}| *|\mathbf{G}|\|_{\nu} \leq\|\mathbf{F}\|_{\nu}\|\mathbf{G}\|_{\nu}
$$

Proof. The proof follows readily from Lemma 14 and definition of $\|\cdot\|_{\nu}$.

Lemma 16 For $\nu>4 \rho_{0}+1$,

$$
\begin{equation*}
\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}} * \mathbf{F}^{* k} * \prod_{l=1}^{m} \prod_{j=1}^{n}\left(p^{j} F_{l}\right)^{* q_{l, j}}\right| \leq \frac{e^{\nu|p|}|p|^{\sum j q_{l, j}}}{M_{0}\left(1+|p|^{2}\right)}\|\mathbf{F}\|_{\nu}^{\langle\mathbf{q}\rangle+\langle\mathbf{k}\rangle-1}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{F}|\right\|_{\nu} \tag{33}
\end{equation*}
$$

if $(\mathbf{q}, \mathbf{k}) \neq(\mathbf{0}, \mathbf{0})$ and is zero if $(\mathbf{q}, \mathbf{k})=(\mathbf{0}, \mathbf{0})$.
Proof. For $(\mathbf{k}, \mathbf{q})=(\mathbf{0}, \mathbf{0})$ we have $\mathbf{B}_{\mathbf{q}, \mathbf{k}}=0$ (see comments after eq. (6)). If $\mathbf{k} \neq \mathbf{0}$, we can use Corollary 13 to argue that the left hand side of (33) is bounded by

$$
\left.|p|^{\sum j q_{l, j}}| | \mathbf{B}_{\mathbf{q}, \mathbf{k}}|*| \mathbf{F}|*| \mathbf{F}\right|^{*(\langle\mathbf{k}\rangle-1)} * \prod_{l=1}^{m} \prod_{j=1}^{n}|\mathbf{F}|^{* q_{l, j}} \mid
$$

Using Corollaries 10 and 15 , the proof if $\mathbf{k} \neq 0$ follows. Similar steps work for the case $\mathbf{k}=\mathbf{0}$ and $\mathbf{q} \neq \mathbf{0}$, except that $\mathbf{B}_{\mathbf{q}, \mathbf{k}}$ is convolved with $p^{j_{1}} F_{l_{1}}$ for some ( $j_{1}, l_{1}$ ), for which the corresponding $q_{l_{1}, j_{1}} \neq 0$, and we now use Lemma 12 and the definition of $\|\cdot\|_{\nu}$.

Corollary 17 For $\nu>4 \rho_{0}+1$,

$$
\begin{equation*}
\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}} * \mathbf{F}^{* k} * \prod_{l=1}^{*} \prod_{j=1}^{n}\left(p^{j} F_{l}\right)^{* q_{l, j}}\right| \leq \frac{K C_{1} A_{b}(T) e^{\nu|p|}|p|^{\sum j q_{l, j}}}{M_{0}\left(1+|p|^{2}\right)}\left(\frac{\nu}{2}\right)^{-\langle\mathbf{k}\rangle \beta-\alpha_{\mathbf{q}}}\|\mathbf{F}\|_{\nu}^{\langle\mathbf{q}\rangle+\langle\mathbf{k}\rangle} \tag{34}
\end{equation*}
$$

Proof. The proof follows immediately from Corollary 10 and Lemma 16.

Lemma 18 For $\nu>4 \rho_{0}+1$ we have

$$
\begin{array}{rl}
\mid \int_{0}^{t} e^{-p^{n}(t-\tau)} \mathbf{B}_{\mathbf{q}, \mathbf{k}} * \mathbf{F}^{* k} * \prod_{l=1}^{m} \prod_{j=1}^{n}\left(p^{j} F_{l}\right)^{* q_{l, j}} & d \tau \mid \\
\leq \frac{C A_{b}(T) e^{\nu|p|}}{M_{0}\left(1+|p|^{2}\right)}\left(\frac{\nu}{2}\right)^{-\langle\mathbf{k}\rangle \beta-\alpha_{\mathbf{q}}}\|\mathbf{F}\|_{\nu}^{\langle\mathbf{q}\rangle+\langle\mathbf{k}\rangle}\left(\frac{\nu}{2}\right)^{-\beta\langle\mathbf{k}\rangle-\alpha_{\mathbf{q}}} T^{\left(n-\sum j q_{l, j}\right) / n} \tag{35}
\end{array}
$$

where the constant $C$ is independent of $T$, but depends on $\phi$.
Proof. The proof follows from Lemmas 11 and 16 and the fact that for $0 \leq l \leq n$,

$$
\begin{equation*}
|p|^{l} \int_{0}^{t} e^{-|p|^{n} \cos (n \theta)(t-\tau)} d \tau \leq \frac{T^{(n-l) / n}}{\cos ^{l / n}(n \phi)} \sup _{\gamma} \frac{1-e^{-\gamma^{n}}}{\gamma^{n-l}} \tag{36}
\end{equation*}
$$

Definition 19 For $\mathbf{F}$ and $\mathbf{h}$ in $\mathcal{A}_{\phi}$, and $\mathbf{B}_{\mathbf{q}, \mathbf{k}} \in \mathcal{H}$, as above, define $\mathbf{h}_{0}=\mathbf{0}$ and for $k \geq 1$,

$$
\begin{equation*}
\mathbf{h}_{\mathbf{k}} \equiv \mathbf{B}_{\mathbf{q}, \mathbf{k}} *\left[(\mathbf{F}+\mathbf{h})^{* \mathbf{k}}-\mathbf{F}^{* \mathbf{k}}\right] \tag{37}
\end{equation*}
$$

Lemma 20 For $\nu>4 \rho_{0}+1$, and for $\mathbf{k} \neq 0$,

$$
\begin{equation*}
\left\|\mathbf{h}_{\mathbf{k}}\right\|_{\nu} \leq\langle\mathbf{k}\rangle\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{\langle\mathbf{k}\rangle-1}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{h}|\right\|_{\nu} \tag{38}
\end{equation*}
$$

and is zero for $\mathbf{k}=0$.
Proof. The case of $\mathbf{k}=0$ follows from definition of $\mathbf{h}_{0}$. The general expression above for $\mathbf{k} \neq 0$ is proved by induction. The case of $\langle\mathbf{k}\rangle=1$ is obvious from (37). Assume formula (38) holds for all $\langle\mathbf{k}\rangle \leq l$. Then all multiindices of length $l+1$ can be expressed as $\mathbf{k}+\mathbf{e}_{i}$, where $\mathbf{e}_{i}$ is an the $m$ dimensional unit vector in the $i$-th direction for some $i$, and $\mathbf{k}$ has length $l$.
$\left\|\mathbf{h}_{\mathbf{k}+\mathbf{e}_{i}}\right\|_{\nu}=\left\|\mathbf{B}_{\mathbf{q}, \mathbf{k}} *\left(F_{i}+h_{i}\right) *(\mathbf{F}+\mathbf{h})^{* \mathbf{k}}-\mathbf{B}_{\mathbf{q}, \mathbf{k}} * F_{i} * \mathbf{F}^{* \mathbf{k}}\right\|_{\nu}=\left\|\mathbf{B}_{\mathbf{q}, \mathbf{k}} * h_{i} *(\mathbf{F}+\mathbf{h})^{* \mathbf{k}}+F_{i} * \mathbf{h}_{\mathbf{k}}\right\|_{\nu}$
Using (38) for $\langle\mathbf{k}\rangle=l$, we get

$$
\begin{gathered}
\leq\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{h}|\right\|_{\nu}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{l}+l\|\mathbf{F}\|_{\nu}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{l-1}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{h}|\right\|_{\nu} \\
\leq(l+1)\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{l}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{h}|\right\|_{\nu}
\end{gathered}
$$

Thus (38) holds for $\langle\mathbf{k}\rangle=l+1$.

Definition 21 For $\mathbf{F} \in \mathcal{A}_{\phi}$ and $\mathbf{h} \in \mathcal{A}_{\phi}$, and $\mathbf{B}_{\mathbf{q}, \mathbf{k}}$ as above define $\mathbf{g}_{\mathbf{0}}=\mathbf{0}$, and for $\langle\mathbf{q}\rangle \geq 1$,

$$
\begin{equation*}
\mathbf{g}_{\mathbf{q}} \equiv \mathbf{B}_{\mathbf{q}, \mathbf{k}} * \prod_{l=1}^{*} \prod_{j=1}^{*}\left(p^{j}\left[F_{l}+h_{l}\right]\right)^{* q_{l, j}}-\mathbf{B}_{\mathbf{q}, \mathbf{k}} * \prod_{l=1}^{*} \prod_{j=1}^{*}\left(p^{j} F_{l}\right)^{* q_{l, j}} \tag{39}
\end{equation*}
$$

Lemma 22 For $\nu>4 \rho_{0}+1, \mathbf{g}_{0}=0$ and for $\langle\mathbf{q}\rangle \geq 1$

$$
\begin{equation*}
\left|\mathbf{g}_{\mathbf{q}}\right| \leq \frac{e^{\nu|p|}|p|^{\sum j q_{l, j}}\langle\mathbf{q}\rangle}{M_{0}\left(1+|p|^{2}\right)}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{\langle\mathbf{q}\rangle-1}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{h}|\right\|_{\nu} \tag{40}
\end{equation*}
$$

and is zero for $\mathbf{q}=0$.
Proof. The case for $\mathbf{q}=0$ follows by definition of $\mathbf{g}_{\boldsymbol{0}}$. We prove the other cases by induction. The case $\langle\mathbf{q}\rangle=1$ is clear from (39), since only linear terms in $\mathbf{F}$ are involved. Assume the
inequality (40) holds for a particular $\mathbf{q}$. We now show that it holds when $\mathbf{q}$ is replaced by $\mathbf{q}+\mathbf{e}$, where $\mathbf{e}$ is a $m \times n$ dimensional unit vector, say in the $\left(l_{1}, j_{1}\right)$ direction. So,

$$
\begin{align*}
& \left|\mathbf{g}_{\mathbf{q}+\mathbf{e}}\right| \\
& =\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}} *\left[p^{j_{1}}\left(F_{l_{1}}+h_{l_{1}}\right)\right] * \prod_{l=1}^{*} \prod_{j=1}^{n}\left[p^{j}(\mathbf{F}+\mathbf{h})\right]^{* q_{l, j}}-\mathbf{B}_{\mathbf{q}, \mathbf{k}} *\left[p^{j_{1}} F_{l_{1}}\right] * \prod_{l=1}^{*} \prod_{j=1}^{n}\left[p^{j} \mathbf{F}\right]^{* q_{l, j}}\right| \\
& \leq\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}} *\left(p^{j_{1}} h_{l_{1}}\right)\right| *\left|\prod_{l=1}^{*} \prod_{j=1}^{*}\left[p^{j}(\mathbf{F}+\mathbf{h})\right]^{* q_{l, j}}\right|+\left|\left(p^{j_{1}} F_{l_{1}}\right) * \mathbf{g}_{\mathbf{q}}\right| \tag{41}
\end{align*}
$$

Using Lemma 16 and equation (40), we get the following upper bound

$$
\begin{aligned}
& \left|\mathbf{g}_{\mathbf{q}+\mathbf{e}}\right| \leq \frac{|p|^{j_{1}+\sum j q_{l, j}} e^{\nu|p|}}{M_{0}\left(1+|p|^{2}\right)}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{\sum q_{l, j}}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{h}|\right\|_{\nu} \\
& +\frac{|p|^{j_{1}+\sum j q_{l, j}}\langle\mathbf{q}\rangle e^{\nu|p|}}{M_{0}\left(1+|p|^{2}\right)}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{\langle\mathbf{q}\rangle-1}\|\mathbf{F}\|_{\nu}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{h}|\right\|_{\nu} \\
& \leq \frac{|p|^{\sum j\left(q_{l, j}+e_{l, j}\right)}(\langle\mathbf{q}+\mathbf{e}\rangle) e^{\nu|p|}}{M_{0}\left(1+|p|^{2}\right)}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{\langle\mathbf{q}\rangle}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{h}|\right\|_{\nu}
\end{aligned}
$$

Therefore (40) holds when $\mathbf{q}$ is replaced by $\mathbf{q}+\mathbf{e}$ and the induction step is proved.

Lemma 23 For $\mathbf{F}$ and $\mathbf{h}$ in $\mathcal{A}_{\phi}, \nu>4 \rho_{0}+1$,

$$
\begin{gather*}
\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}} *(\mathbf{F}+\mathbf{h})^{* \mathbf{k}} * \prod_{l=1}^{m} \prod_{j=1}^{n}\left(p^{j}\left(F_{l}+h_{l}\right)\right)^{* q_{l, j}}-\mathbf{B}_{\mathbf{q}, \mathbf{k}} * \mathbf{F}^{* \mathbf{k}} *{ }^{*} \prod_{l=1}^{m} \prod_{j=1}^{n}\left(p^{j} F_{l}\right)^{* q_{l, j}}\right| \\
\leq \frac{|p|^{\sum j q_{l, j}}(\langle\mathbf{q}\rangle+\langle\mathbf{k}\rangle) e^{\nu|p|}}{M_{0}\left(1+|p|^{2}\right)}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{\langle\mathbf{k}\rangle+\langle\mathbf{q}\rangle-1}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| * \mid \mathbf{h}\right\| \|_{\nu} \tag{42}
\end{gather*}
$$

if $(\mathbf{q}, \mathbf{k}) \neq(\mathbf{0}, \mathbf{0})$ and is zero otherwise.

Proof. It is clear from (37) that the left side of (42) is simply

$$
\left|\mathbf{h}_{\mathbf{k}} * \prod_{l=1}^{*} \prod_{j=1}^{n}\left(p^{j}\left(F_{l}+h_{l}\right)\right)^{* q_{l, j}}+\mathbf{F}^{* \mathbf{k}} * \mathbf{g}_{\mathbf{q}}\right|
$$

However, from Corollary 13, Lemmas 14 and 20,

$$
\left|\mathbf{h}_{\mathbf{k}} *^{*} \prod_{l=1}^{m} \prod_{j=1}^{n}\left(p^{j}\left(F_{l}+h_{l}\right)\right)^{* q_{l, j}}\right| \leq \frac{|p|^{\sum j q_{l, j}}\langle\mathbf{k}\rangle e^{\nu|p|}}{M_{0}\left(1+|p|^{2}\right)}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{\langle\mathbf{k}\rangle+\langle\mathbf{q}\rangle-1}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{h}|\right\|_{\nu}
$$

and from Corollary 13, Lemmas 14 and 22,

$$
\left|\mathbf{F}^{* \mathbf{k}} * \mathbf{g}_{\mathbf{q}}\right| \leq \frac{|p|^{\sum j q_{l, j}}\langle\mathbf{q}\rangle e^{\nu|p|}}{M_{0}\left(1+|p|^{2}\right)}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{\langle\mathbf{k}\rangle+\langle\mathbf{q}\rangle-1}\left\|\left|\mathbf{B}_{\mathbf{q}, \mathbf{k}}\right| *|\mathbf{h}|\right\|_{\nu}
$$

Combining these two inequalities, the proof of the lemma follows.

Lemma 24 For $\nu>4 \rho_{0}+1$ we have

$$
\begin{align*}
& \| \int_{0}^{t} e^{-p^{n}(t-\tau)}\left[\mathbf{B}_{\mathbf{q}, \mathbf{k}} *(\mathbf{F}+\mathbf{h})^{* \mathbf{k}} *{ }^{*} \prod_{l=1}^{m} \prod_{j=1}^{n}\left(p^{j}\left(F_{l}+h_{l}\right)\right)^{* q_{l, j}}\right. \\
& \left.\quad-\mathbf{B}_{\mathbf{q}, \mathbf{k}} * \mathbf{F}^{* \mathbf{k}} * \prod_{l=1}^{*} \prod_{j=1}^{m}\left(p^{j} F_{l}\right)^{* q_{l, j}}\right] d \tau \|_{\nu}^{n} \\
& \quad \leq A_{b}(T) C(\phi)(\langle\mathbf{q}\rangle+\langle\mathbf{k}\rangle)\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{\langle\mathbf{k}\rangle+\langle\mathbf{q}\rangle-1} T^{\left(n-\sum j q_{l, j}\right) / n}\left(\frac{\nu}{2}\right)^{-\beta\langle\mathbf{k}\rangle-\alpha_{\mathbf{q}}}\|\mathbf{h}\|_{\nu} \tag{43}
\end{align*}
$$

Proof. The proof follows from Corollary 10 and Lemma 23 and the definition of $\|\cdot\|_{\nu}$ along with the bound (36).

Lemma 25 For $\mathbf{F} \in \mathcal{A}_{\phi}$, and $\nu>4 \rho_{0}+\alpha_{r}$ large enough so that $\left(\frac{\nu}{2}\right)^{-\beta}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)<1$ (see Note after Definition (5)), $\mathcal{N}(\mathbf{F})$ defined in (21) satisfies the following bounds

$$
\begin{align*}
& \|\mathcal{N}(\mathbf{F})\|_{\nu} \leq\left\|\mathbf{F}_{0}\right\|_{\nu}+C(\phi) A_{b}(T) \sum_{\mathbf{q} \succeq 0}^{\prime} \sum_{\mathbf{k} \succeq 0} T^{\left(n-\sum j q_{l, j}\right) / n}\left(\frac{2^{\beta}\|\mathbf{F}\|_{\nu}}{\nu^{\beta}}\right)^{\langle\mathbf{k}\rangle}\left(\frac{\nu}{2}\right)^{-\alpha_{\mathbf{q}}}\|\mathbf{F}\|_{\nu}^{\langle\mathbf{q}\rangle}  \tag{44}\\
& \|\mathcal{N}(\mathbf{F}+\mathbf{h})-\mathcal{N}(\mathbf{F})\|_{\nu} \\
& \leq C(\phi) A_{b}(T)\|\mathbf{h}\|_{\nu} \sum_{\mathbf{q} \succeq 0}^{\prime} \sum_{\mathbf{k} \succeq 0} T^{\left(n-\sum j q_{l, j}\right) / n}\left(\frac{\nu}{2}\right)^{-\beta\langle\mathbf{k}\rangle-\alpha_{\mathbf{q}}}(\langle\mathbf{q}\rangle+\langle\mathbf{k}\rangle)\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)^{\langle\mathbf{q}\rangle+\langle\mathbf{k}\rangle-1} \tag{45}
\end{align*}
$$

Proof. The proofs are immediate from the expression (21) of $\mathcal{N}(\mathbf{F})$ and Lemmas 18, 20 and 24. The condition $\left(\frac{\nu}{2}\right)^{-\beta}\left(\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu}\right)<1$ guarantees the convergence of the infinite series involving summation in $\mathbf{k}$. Note that with respect to the multi-index $\mathbf{q}$ we only have a finite sum, due to (4). Note also that the summation over $\mathbf{k}$ can also be bounded by a more explicit function, if so desired.

Comment 7: Lemma 25 is the key to showing the existence and uniqueness of a solution in $\mathcal{A}_{\phi}$ to (21), since it provides the conditions for the nonlinear operator $\mathcal{N}$ to map a ball into itself as well the necessary contractivity condition.

Lemma 26 If there exists some $b>1$ so that

$$
\begin{equation*}
(\nu / 2)^{-\beta} b\left\|\mathbf{F}_{0}\right\|_{\nu}<1 \tag{46}
\end{equation*}
$$

and

$$
\begin{equation*}
C(\phi) A_{b}(T) \sum_{\mathbf{q} \succeq 0}^{\prime} \sum_{\mathbf{k} \succeq 0} T^{\left(n-\sum j q_{l, j}\right) / n}\left(\frac{\nu}{2}\right)^{-\langle\mathbf{k}\rangle \beta-\alpha_{\mathbf{q}}}\left\|b \mathbf{F}_{0}\right\|_{\nu}^{\langle\mathbf{k}\rangle+\langle\mathbf{q}\rangle}<1-\frac{1}{b} \tag{47}
\end{equation*}
$$

then the nonlinear mapping $\mathcal{N}$, as defined in (21), maps a ball of radius $b\left\|\mathbf{F}_{0}\right\|_{\nu}$ into itself. Further, if

$$
\begin{equation*}
C(\phi) A_{b}(T) \sum_{\mathbf{q} \succeq 0}^{\prime} \sum_{\mathbf{k} \succeq 0} T^{\left(n-\sum j q_{l, j}\right) / n}(\langle\mathbf{q}\rangle+\langle\mathbf{k}\rangle)\left(\frac{\nu}{2}\right)^{-\langle\mathbf{k}\rangle \beta-\alpha_{\mathbf{q}}}(3 b)^{\langle\mathbf{k}\rangle+\langle\mathbf{q}\rangle-1}\left\|\mathbf{F}_{0}\right\|_{\nu}^{\langle\mathbf{k}\rangle+\langle\mathbf{q}\rangle-1}<1 \tag{48}
\end{equation*}
$$

then $\mathcal{N}$ is a contraction there.
Proof. This is a simple application of Lemma 25, if we note that $\|\mathbf{F}\|_{\nu}^{k}<b^{k}\left\|\mathbf{F}_{0}\right\|_{\nu}^{k}$ and the fact that for both $\mathbf{F}$ and $\mathbf{F}+\mathbf{h}$ in the ball of radius $b\|\mathbf{F}\|_{0},\|\mathbf{F}\|_{\nu}+\|\mathbf{h}\|_{\nu} \leq 3 b\left\|\mathbf{F}_{0}\right\|_{\nu}$.

Lemma 27 For any given $T>0$ and $\phi$ in the interval $\left(0,(2 n)^{-1} \pi\right)$, for all sufficiently large $\nu$, there exists a unique $\mathbf{F} \in \mathcal{A}_{\phi}$ that satisfies the integral equation (21).

Proof. We choose $b=2$ for definiteness. It is clear from the bounds on $\left\|F_{0}\right\|_{\nu}$ in Lemma 7 that for given $T$, since $\alpha_{r} \geq 1$, we have $b(\nu / 2)^{-\beta}\left\|\mathbf{F}_{0}\right\|_{\nu}<1$ for all sufficiently large $\nu$. Further, it is clear by inspection that all conditions (46), (47) and (48) are satisfied for all sufficiently large $\nu$. The lemma now follows from the contractive mapping theorem.

### 4.1 Behavior of ${ }^{s} F$ near $p=0$

Proposition 28 For some $K_{1}>0$ and small $p$ we have $\left|{ }^{s} \mathbf{F}\right|<K_{1}|p|^{\alpha_{r}-1}$ and thus $|\mathbf{s}|<$ $K_{2}|y|^{-\alpha_{r}}$ for some $K_{2}>0$ in $\mathcal{D}_{\phi, \rho}$ as $|y| \rightarrow \infty$.

Proof. The idea of the proof here is to think of the solution ${ }^{s} \mathbf{F}$ to (21) as a solution to a linear equation of the form

$$
\begin{equation*}
{ }^{s} \mathbf{F}=\mathcal{G}\left({ }^{s} \mathbf{F}\right)+\mathbf{F}_{0} \quad \text { or } \quad{ }^{s} \mathbf{F}=(1-\mathcal{G})^{-1} \mathbf{F}_{0} \tag{49}
\end{equation*}
$$

Here, the suitably chosen operator $\mathcal{G}$, while depending on ${ }^{\mathcal{S}} \mathbf{F}$, is thought of a known quantity (in effect, ${ }^{s} \mathbf{F}$ is now known). The expression of a suitable $\mathcal{G}$ is, however, somewhat involved, and this is given in the following.

Convergence in $\|\cdot\|_{\nu}$ implies uniform convergence on compact subsets of $\mathcal{K}$ and we can interchange summation and integration in (21). With ${ }^{5} \mathbf{F}$ the unique solution of (21) we define

$$
\mathbf{G}_{i}=\sum^{\dagger} \mathbf{B}_{\mathbf{q}, \mathbf{k}} *{ }^{s} \mathbf{F}^{* \mathbf{k}_{i c}} * \prod_{l=1}^{m} \prod_{j=1}^{n}\left[(-p)^{j}\left({ }^{s} F_{l}\right)\right]^{* q_{l, j}}
$$

where $\sum^{\dagger}$ stands for the sum over $\mathbf{k} \neq 0, k_{i} \neq 0$ and $k_{i^{\prime}}=0$ for $i^{\prime}<i$, the notation ${ }^{s} \mathbf{F}^{* \mathbf{k}_{i c}}$ stands for $\left({ }^{s} F_{1}\right)^{* k_{1}} *\left({ }^{s} F_{2}\right)^{* k_{2}} * \cdots *\left({ }^{s} F_{m}\right)^{* k_{m}}$, except that the term ${ }^{s} F_{i}^{* k_{i}}$ is replaced by ${ }^{s} F_{i}^{*\left(k_{i}-1\right)}$ in this convolution.

$$
\hat{\mathbf{G}}_{j^{\prime}, l^{\prime}}=\sum^{\ddagger} \mathbf{B}_{\mathbf{q}, 0} * \prod_{l=1}^{*} \prod_{j=1}^{\ddagger}\left[(-p)^{j}\left({ }^{s} F_{l}\right)\right]^{* q_{l, j}}
$$

where $\sum^{\ddagger}$ stands for the sum over $\mathbf{q} \succ 0, q_{l^{\prime}, j^{\prime}} \neq 0$, and $q_{l, j}=0$ for $l+j<l^{\prime}+j^{\prime}$ and for $l+j=l^{\prime}+j^{\prime}$, when $l<l^{\prime}$. Also, $\prod^{\ddagger}$ indicates that the exponent of the term $l=l^{\prime}, j=j^{\prime}$ is changed from $q_{l^{\prime}, j^{\prime}}$ to $q_{l^{\prime}, j^{\prime}}-1$. Now we define a linear operator $\mathcal{G}$ by

$$
\mathcal{G} \mathbf{Q}=\int_{0}^{t} e^{-p^{n}(t-\tau)}\left[\sum_{i=1}^{m} \mathbf{G}_{i} * Q_{i}+\sum_{j^{\prime}=1}^{n} \sum_{l^{\prime}=1}^{m} \hat{\mathbf{G}}_{j^{\prime}, l^{\prime}} *\left((-p)^{j} Q_{l^{\prime}}\right)\right]
$$

Then, by carefully comparing with (21), one finds that ${ }^{s} \mathbf{F}$ satisfies (49).
For $a>0$ small enough define $\overline{\mathcal{S}}_{a}=\overline{\mathcal{S}} \cap\{p:|p| \leq a\}$. Since ${ }^{\boldsymbol{S}} \mathbf{F}$ is continuous in $\overline{\mathcal{S}}$ we have $\lim _{a \downarrow 0}\|\mathcal{G}\|=0$, where the norm is taken over $C\left(\overline{\mathcal{S}}_{a}\right)$.

By (5), (10), (22) and Lemma 3, we see that $\left\|\mathbf{F}_{0}\right\|_{\infty} \leq K_{3}|a|^{\alpha_{r}-1}$ in $\overline{\mathcal{S}}_{a}$ for some $K_{3}>0$ independent of $a$. Then, as $a \downarrow 0$, we have

$$
\left.\left.\max _{\overline{\mathcal{S}}_{a}}\right|^{s} \mathbf{F}(p, t)\left|=\left\|{ }^{s} \mathbf{F}\right\| \leq(1-\|\mathcal{G}\|)^{-1} \max _{\overline{\mathcal{S}}_{a}}\left\|\mathbf{F}_{0}\right\| \leq 2 K_{3}\right| a\right|^{\alpha_{r}-1}
$$

and thus for small $p$ we have $|\mathbf{F}(p, t)| \leq 2 K_{3}|p|^{\alpha_{r}-1}$ and the proposition follows.

### 4.2 Completion of proof of Theorem 2.

Lemma 3 shows that if $\mathbf{f}$ is any solution of (3) satisfying Condition 1 , then $\mathcal{L}^{-1}\{\mathbf{f}\} \in \mathcal{A}_{\phi-\delta}$ for $0<\delta<\phi$ for $\nu$ sufficiently large. For large $y$, the series (6) converges uniformly and thus $\mathbf{F}=\mathcal{L}^{-1}\{\mathbf{f}\}$ satisfies (21), which by Lemma 27 has a unique solution in $\mathcal{A}_{\phi}$ for any $\phi \in$ ( $\left.0,(2 n)^{-1} \pi\right)$. Conversely, if ${ }^{s} \mathbf{F} \in \mathcal{A}_{\tilde{\phi}}$ is the solution of (21) for $\nu>\nu_{1}$, then from Comment $5,{ }^{s} \mathbf{f}=\mathcal{L}^{s} \mathbf{F}$ is analytic in $y$ in $\mathcal{D}_{\phi, \rho}$ for $0<\phi<\tilde{\phi}<(2 n)^{-1} \pi$, for sufficiently large $\rho$, where in addition from Proposition 28, ${ }^{s} \mathbf{f}=O\left(y^{-\alpha_{r}}\right)$. This implies that the series in (3) converges uniformly and by the properties of Laplace transforms, ${ }^{s} \mathbf{f}$ solves (1) and satisfies condition (1).

## 5 Formal expansions and their rigorous justification

5.1 In a heuristic calculation of a formal solution to (3) relying on the smallness of $\mathbf{f}$ and our assumptions on the nonlinearity, the most important terms for large $y$ (giving the "dominant balance") are $\mathbf{f}_{t}$ on the left side of (3) and $\mathbf{r}(y, t)$ on the right side. This suggests that, to leading order, $\mathbf{f}(y, t) \sim \mathbf{f}_{I}(y)+\int_{0}^{t} \mathbf{r}(y, t) d t$ (since the functions $\mathbf{f}_{I}(y)$ and $\mathbf{r}(y, t)$ decay at a rate $y^{-\alpha_{r}}$, which for $\alpha_{r}>1$ is much less than $y^{-1}$, other terms in the differential equation (3) should not contribute). We then decompose $\mathbf{f}(y, t)=A_{1}(t) y^{-\alpha_{r}}+\tilde{\mathbf{f}}$ and substitute into (3); the equation
for $\tilde{\mathbf{f}}$ will generally satisfy an equation of the same form as $\mathbf{f}$ in (3), but with $\alpha_{r}$ replaced by a larger number; this procedure generates in principle a formal asymptotic series for the solution, as is the case in the examples in [7].
5.2 When the formal procedure in the preceding section gives the leading order behavior $\mathbf{f}(y, t) \sim$ $A_{1}(t) y^{-\alpha_{r}}$, the validity of this asymptotic relation is rigorously shown in $\mathcal{D}_{\phi, \tilde{\rho}}$ as follows. We write as before $\mathbf{f}(y, t)=A_{1}(t) y^{-\alpha_{r}}+\tilde{\mathbf{f}}$ in (3). With $\alpha_{r}$ in the equation for $\tilde{\mathbf{f}}$ larger than the $\alpha_{r}$ of $\mathbf{f}$, our theorem guarantees that $\tilde{\mathbf{f}}$ is indeed $o\left(y^{-\alpha_{r}}\right)$. We can recursively use this procedure on $\tilde{\mathbf{f}}$, and so on, to justify the asymptotic expansion of $\mathbf{f}$ for large $y$ in $\mathcal{D}_{\phi, \rho ; y}$.

These arguments also show that the assumption $\alpha_{r} \geq 1$ is not crucial, since the terms in the asymptotic series for $\mathbf{f}$ for which the exponent does not satisfy this condition can be subtracted out explicitly, as there are generally finitely many of them.
5.3 Conversely, if the inverse Laplace transform ${ }^{s} \mathbf{F}$ of the solution ${ }^{\mathbf{s} f}$ has a convergent Puiseux series about $p=0$, which is for instance the case in the examples treated in [7], then by Watson's Lemma ${ }^{\text {s }} \mathbf{f}$ will have for large $y$ an asymptotic series in inverse powers of $y$, and the representation ${ }^{\mathbf{s} \mathbf{f}}=\mathcal{L}\left\{{ }^{\mathbf{s} \mathbf{F}}\right\}$ makes ${ }^{\mathbf{s} \mathbf{f}}$ (by definition) the Borel sum of its asymptotic series. (The large width of the sector is then needed to guarantee the uniqueness).

## 6 Appendix

### 6.1 Derivation of equation (3) from (1)

We start from (1), where $\mathbf{h}$ is an $m_{1}$ dimensional vector field depending on $y$ and $t$. We define the $m=m_{1} \times n$-dimensional vector field as $\mathbf{f}=\left(\mathbf{h}, \partial_{y} \mathbf{h}, \partial_{y}{ }^{2} \mathbf{h}, \ldots, \partial_{y}{ }^{(n-1)} \mathbf{h}\right)$. Then it is clear from (50) that $\mathbf{g}_{1}$ and $\mathbf{g}_{2}$ only depend on $\mathbf{f}$. So, for showing that (1) implies (3) it is enough to show that for $1 \leq n^{\prime} \leq n$,

$$
\partial^{n^{\prime}-1}\left[\mathbf{g}_{1}(y, t, \mathbf{f})+\mathbf{g}_{2}(y, t, \mathbf{f}) \partial_{y} \mathbf{f}\right]
$$

is of the form on the right hand side of (3). We do so in two steps.
Lemma 29 For any $n^{\prime} \geq 1$,

$$
\begin{equation*}
\partial_{y}^{n^{\prime}-1} \mathbf{g}_{1}(y, t, \mathbf{f}(y, t))=\sum_{\mathbf{q} \succeq 0}^{\ddagger} \mathbf{b}_{\mathbf{q}}(y, t, \mathbf{f}) \prod_{l=1}^{m} \prod_{j=1}^{n^{\prime}-1}\left(\partial_{y}^{j} f_{l}\right)^{q_{l, j}} \tag{50}
\end{equation*}
$$

for some $\mathbf{b}_{\mathbf{q}, \mathbf{k}}$, depending on $n^{\prime}, g_{1}$, and its first $n^{\prime}-1$ derivatives with respect to its arguments, and where $\sum^{\ddagger}$ means the sum with the further restriction

$$
\sum_{l, j} j q_{l, j} \leq n^{\prime}-1
$$

Proof. The proof is by induction. We have, with obvious notation,

$$
\partial_{y} \mathbf{g}_{1}(y, t, \mathbf{f}(y, t))=\mathbf{g}_{1, y}+\mathbf{g}_{1, \mathbf{f}} \cdot \partial_{y} \mathbf{f}
$$

which is of the form (50). Assume (50) holds for $n^{\prime}=k \geq 1$, i.e.

$$
\partial_{y}^{k-1} \mathbf{g}_{1}(y, t, \mathbf{f})=\sum_{\mathbf{q} \succeq 0}^{\ddagger} \mathbf{b}_{\mathbf{q}}(y, t, \mathbf{f}) \prod_{l=1}^{m} \prod_{j=1}^{k-1}\left(\partial_{y}^{j} f_{l}\right)^{q_{l, j}}
$$

Taking a $y$ derivative, we get

$$
\begin{aligned}
& \partial_{y}^{k} \mathbf{g}_{\mathbf{1}}(y, t, \mathbf{f})=\sum_{\mathbf{q} \succeq 0}\left(\sum_{i=1}^{m} \frac{\partial}{\partial f_{i}} \mathbf{b}_{\mathbf{q}}(y, t, \mathbf{f}) \partial_{y} f_{i}+\partial_{1} \mathbf{b}_{\mathbf{q}}(y, t, \mathbf{f})\right) \prod_{l=1}^{m} \prod_{j=1}^{k-1}\left(\partial_{y}^{j} f_{l}\right)^{q_{l, j}} \\
& \quad+\sum_{\mathbf{q} \succeq 0} \mathbf{b}_{\mathbf{q}}(y, t, \mathbf{f}) \sum_{l^{\prime}=1}^{m} \sum_{j^{\prime}=1}^{k-1} q_{l^{\prime}, j^{\prime}}\left(\partial_{y}^{j^{\prime}} f_{l^{\prime}}\right)^{q_{l^{\prime}, j^{\prime}}-1}\left(\partial_{y}^{j^{\prime}+1} f_{l^{\prime}}\right) \prod_{l=1}^{m \backslash l^{\prime}} \prod_{j=1}^{k-1 \backslash j^{\prime}}\left(\partial_{y}^{j} f_{l}\right)^{q_{l, j}}
\end{aligned}
$$

where the notation $\partial_{1} \mathbf{b}_{\mathbf{q}}$ denotes partial with respect to the first argument of $\mathbf{b}$ and $\prod_{l=1}^{m \backslash l^{\prime}}$ stands for the product with $l=l^{\prime}$ term excluded. It is easily seen from the above expressions that the product of the number of derivatives times the power, when added up at most equals

$$
j^{\prime}+1+j^{\prime}\left(q_{l^{\prime}, j^{\prime}}-1\right)+\sum_{j \neq j^{\prime}} \sum_{l \neq l^{\prime}} j q_{l, j}=\sum_{l, j} j q_{l, j}+1 \leq k-1+1=k
$$

Thus, (50) holds for $n^{\prime}=k+1$, with a different $\mathbf{b}$. The induction step is proved.
Lemma 30 For any $n^{\prime}$ between 1 and $n$,

$$
\begin{equation*}
\partial^{n^{\prime}-1}\left[\mathbf{g}_{\mathbf{2}}(y, t, \mathbf{f}) \partial_{y} \mathbf{f}\right]=\sum_{\mathbf{q} \succeq 0}^{\ddagger} \mathbf{b}_{\mathbf{q}}(y, t, \mathbf{f}) \prod_{l=1}^{m} \prod_{j=1}^{n^{\prime}}\left(\partial_{y}^{j} f_{l}\right)^{q_{l, j}} \tag{51}
\end{equation*}
$$

for some $\mathbf{b}_{\mathbf{q}, \mathbf{k}}$, depending on $n^{\prime}, \mathbf{g}_{2}$ and its first $n^{\prime}-1$ derivatives with respect to each argument, and with further restriction

$$
\sum_{l, j} j q_{l, j} \leq n^{\prime}
$$

Proof. It is clear that the $n^{\prime}-1$ derivative of $\mathbf{g}_{2} \partial_{y} \mathbf{f}$ is a linear combination of

$$
\partial_{y}^{n^{\prime \prime}}\left[\mathbf{g}_{2}\right] \partial_{y}^{n^{\prime}-n^{\prime \prime}} \mathbf{f}
$$

for $n^{\prime \prime}$ ranging from 0 to $n^{\prime}-1$. Since the derivatives of $\mathbf{g}_{2}$ are of the form (50), just like $\mathbf{g}_{1}$, it follows that in the above expression the sum of the product of the number of derivatives and the power to which they are raised will be

$$
\sum_{j=1}^{n^{\prime \prime}} \sum_{l=1}^{m} j q_{l, j}+n^{\prime}-n^{\prime \prime} \leq n^{\prime \prime}+n^{\prime}-n^{\prime \prime}=n^{\prime}
$$

Hence the lemma follows.

### 6.2 Simple illustrations of regularization by Borel summation

(i) A vast literature has emerged recently in the field of Borel summation starting with the fundamental contributions of Écalle, see e.g. [8] and it is impossible to give a quick account of the breadth of this field. See for example [6] for more references. Yet, in the context of relatively general PDEs, very little is known. In this section we discuss informally and using admittedly rather trivial examples, the regularizing power of Borel summation techniques.
(ii) Singular perturbations give rise to nonanalytic behavior and divergent series in both ODEs and PDEs. We start by looking at a simple ODE. Infinity is an irregular singular point of the equation $f^{\prime}-f=1 / x$ as can be easily seen after the transformation $x=1 / z$ in which variable the coefficients of the first order equation have, after normalization, a double pole at $x=0$. This is manifested in the factorial divergence of the formal power series solution $\tilde{f}=\sum_{k=0}^{\infty} \frac{(-1)^{k} k!}{x^{k+1}}$

The Borel transform is by definition the formal (term-by-term) inverse Laplace transform and gives $\mathcal{B}(\tilde{f})(p)=\sum_{k=0}^{\infty}(-1)^{k} p^{k}$ which is a convergent series. In effect, after formal inverse Laplace transform, the equation becomes $(p+1) F(p)+1=0\left(F(p):=\mathcal{L}^{-1} f\right)$ which is regular at $p=0$. It is easily seen that $F$ is Laplace transformable and $\mathcal{L} F=f$ is a solution of the original equation.
(iii) To go one step further, consider the heat equation

$$
\begin{equation*}
h_{t}-h_{y y}=0 \tag{52}
\end{equation*}
$$

first for small $t$. (In the present paper we are interested, rather, in finite $t$ and large $y$, situation sketched in (iv) below.) Power series solutions $\tilde{h}=\sum_{k=0}^{\infty} H_{k}(y) t^{k}$, even with $H_{0}$ real-analytic, generally have zero radius of convergence. Indeed, the recurrence relation for the $H_{k}$ is $k H_{k}=$ $H_{k-1}^{\prime \prime}$ meaning that for $H_{0}$ analytic but not entire, $H_{k}$ will roughly behave like $k$ ! const ${ }^{k}$ for large $k$. The factorial divergence suggests as a remedy Borel summation. This is customarily performed with respect to a large variable; substituting $T=1 / t$ in (52) yields

$$
\begin{equation*}
\frac{\partial^{2} h}{\partial y^{2}}+T^{2} \frac{\partial h}{\partial T}+\frac{1}{2} h=0 \tag{53}
\end{equation*}
$$

Inverse Laplace transform of (53) with respect to $T$ yields

$$
\begin{equation*}
H_{y y}-p H_{p p}-\frac{3}{2} H_{p}=0 \tag{54}
\end{equation*}
$$

which is more regular at $p=0$ (dual to $T=\infty$ ). Indeed, substituting $H(p, y)=p^{-1 / 2} F(2 i \sqrt{p}, y)$ transforms (54) into the Laplace equation

$$
F_{y y}+F_{s s}=0
$$

which is elliptic, and a Cauchy-Kowalewski approach leads to analytic solutions. In this sense the problem has thus been regularized. (In fact, it is easy to obtain the classical solution of the heat equation by starting with the general solution of the Laplace equation and inverting the steps.) A thorough study of the Borel summability in the context of the heat equation is found in [14].
(iv) In this paper we focus on (linear or) nonlinear equations, whose singular point is $y=\infty$. Eq. (52) can be used to illustrate the behavior of formal power series solutions in a neighborhood of such a singular point. The substitution

$$
h(t, y)=\sum_{k=1}^{\infty} F_{k}(t) y^{-k}
$$

yields the recurrence $F_{1}^{\prime}=F_{2}^{\prime}=0$ and for $m>2, F_{m}^{\prime}=(m-1)(m-2) F_{m-2}$. Thus, for some constant $c$,

$$
F_{2 n}=4^{n} \Gamma(n+1 / 2) n^{-1} c t^{n-1}
$$

and thus, unless $c=0$ ( or $t=0$ ), this formal power series solution has zero radius of convergence. As seen in this paper in a much more general setting, the divergence is dealt with successfully by Borel (over)summation.
(v) Even when the perturbation is relatively weak in the limit $y \rightarrow \infty$, discarding small parts of the principal symbol of a partial differential operator can be especially problematic. The gist of this difficulty is visible even in a very simple, linear example. The exact solution of the PDE

$$
\begin{equation*}
f_{t}+\left(1+n y^{-n-1}\right)^{-1} f_{y}=0 \tag{55}
\end{equation*}
$$

is $\Phi\left(t-y+y^{-n}\right)$, with $\Phi$ arbitrary, while the equation in which the term $n y^{-n-1}$ is discarded has the solution $\Phi(t-y)$. If $\Phi(v)=\exp (-\exp (-v))$ both solutions are very small, but no matter how large $n$ is, their ratio does not approach one as $y \rightarrow+\infty$. The doubly iterated exponential in the solution is the key to understanding this type of difficulty, and it translates in the presence of solutions as simple transseries of level two [8]. (Roughly it is the number of iterated exponentials that gives the level of the transseries.) While in an ODE with polynomial coefficients this level is bounded (consequently, for instance, in the $\operatorname{ODE}\left(1+y^{-2}\right) f^{\prime}+f$ the term in $y^{-2}$ does not contribute to the leading order asymptotic behavior of the solutions), in PDEs the level can be arbitrary. To exclude in PDEs the presence of too small terms, complex analytic techniques are very useful. Requiring that the solution is a Laplace transform implies algebraic decay in a sector of opening at least $\pi$, ruling out, e.g., the iterated exponential. Then, the ILT with respect to $y$ of the equation (55) with $n=2, F_{t}-p F=\sin (p) * F$ (with $F=\mathcal{L}^{-1} f$ ) can be easily controlled by inverting the ordinary differential operator on the left side and using a contractive mapping argument in a norm $\|\cdot\|_{\nu}$ as introduced in the present paper. This amounts to showing that the term in $y^{-2}$ is negligible-in suitable norms.

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Figure 1: Contour $C_{D}$ in the $p$-plane.


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[^1]:    ${ }^{1}$ More exactly, for the technical reason of simplifying the principal symbol of the transformed operator, we are performing Borel oversummation, meaning that the power of the factorial divided out of the coefficients of the divergent series is higher than the minimum necessary to ensure convergence.

[^2]:    ${ }^{2}$ In the following equation, $\|\cdot\|_{\nu}$ is extended naturally to functions which are only continuous in $\mathcal{K}$.

