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SPACE-TIME ANALYTICITY OF WEAK SOLUTIONS TO SEMILINEAR PARABOLIC SYSTEMS WITH VARIABLE COEFFICIENTS

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Dedicated to the memory of Professor Alan C. Lazer

ABSTRACT. We study analytic smooth solutions of a general, strongly parabolic semilinear Cauchy problem of $2m$ -th order in $\mathbb{R}^N \times (0, T)$ with analytic coefficients (in space and time variables) and analytic initial data (in space variables). They are expressed in terms of holomorphic continuation of global (weak) solutions to the system valued in a suitable Besov interpolation space of $B^{s;p,p}$ -type at every time moment $t \in [0, T]$. Given $0 < T' < T \leq \infty$, it is proved that any $B^{s;p,p}$ -type solution $u : \mathbb{R}^N \times (0, T) \rightarrow \mathbb{C}^M$ with analytic initial data possesses a bounded holomorphic continuation $u(x + iy, \sigma + i\tau)$ into a complex domain in $\mathbb{C}^N \times \mathbb{C}$ defined by $(x, \sigma) \in \mathbb{R}^N \times (T', T)$, $|y| < A'$ and $|\tau| < B'$, where $A', B' > 0$ are constants depending upon T' . The proof uses the extension of a weak solution to a $B^{s;p,p}$ -type solution in a domain in $\mathbb{C}^N \times \mathbb{C}$, such that this extension satisfies the Cauchy-Riemann equations. The holomorphic extension is obtained with a help from holomorphic semigroups and maximal regularity theory for parabolic problems in Besov interpolation spaces of $B^{s;p,p}$ -type imbedded (densely and continuously) into an L^p -type Lebesgue space. Applications include *risk models* for European options in Mathematical Finance.

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1. INTRODUCTION

In this article we investigate the analyticity (in space and time variables) of strict L^p -type solutions $\mathbf{u} = (u_1, \dots, u_M) : \mathbb{R}^N \times (0, T) \rightarrow \mathbb{C}^M$ (or \mathbb{C}^M) of the classical Cauchy problem for a strongly parabolic system of M (coupled) semilinear partial differential equations of order $2m$ ($m \geq 1$ – an integer) with analytic coefficients and with analytic initial data \mathbf{u}_0 belonging to the real interpolation space $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$, such that the function $\mathbf{u} : [0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ is continuous. Here, $\mathbf{B}^{s;p,p}(\mathbb{R}^N) = [B^{s;p,p}(\mathbb{R}^N)]^M$ where $B^{s;p,p}(\mathbb{R}^N)$ denotes the Besov space

$$B^{s;p,p}(\mathbb{R}^N) := (L^p(\mathbb{R}^N), W^{2m,p}(\mathbb{R}^N))_{s/(2m),p} = (L^p(\mathbb{R}^N), W^{2m,p}(\mathbb{R}^N))_{1-(1/p),p}$$

with $1 < p < \infty$, $p > 2 + \frac{N}{m}$, and $s = 2m(1 - \frac{1}{p}) \in (0, 2m)$. This space is defined by real interpolation, e.g., in Adams and Fournier [1, Chapt. 7], §7.6–§7.23, pp. 208–221, Lunardi [65, Chapt. 1], §1.2.2, pp. 20–25, or in Triebel [84, Chapt. 1], §1.2–§1.8, pp. 18–55. Since the Besov space $B^{s;p,p}(\mathbb{R}^N)$ is not imbedded into the Hilbert space $L^2(\mathbb{R}^N)$ whenever $2 < p < \infty$, we find it convenient to consider *strict* L^p -type solutions $\mathbf{u} : \mathbb{R}^N \times (0, T) \rightarrow \mathbb{C}^M$ having the *maximal regularity property* (cf. Ashyralyev and Sobolevskii [9, Chapt. 3, pp. 21–36] and Prüss [74]) rather than weak L^2 -type solutions treated in Takáč [82] for the corresponding linear partial differential equation, but with arbitrary *nonsmooth* initial data $\mathbf{u}_0 \in \mathbf{L}^2(\mathbb{R}^N)$. Consequently, we will be able to apply the classical theory of linear and semilinear evolutionary problems of parabolic type in a Besov space as presented, e.g., in Amann [6, Chapt. III, §4, pp. 128–191], Clément and Li [20], Lunardi [65, Chapt. 7, pp. 257–289], Köhne, Prüss, and Wilke [55], and Tanabe [81, Chapt. 5–6, pp. 117–229]. Our Cauchy problem has the following general form for a semilinear $2m^{\text{th}}$ -order parabolic problem,

$$\begin{aligned} \frac{\partial \mathbf{u}}{\partial t} + \mathbf{P}\left(x, t, \frac{1}{i} \frac{\partial}{\partial x}\right) \mathbf{u} &= \mathbf{f}\left(x, t; \left(\frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta}\right)_{|\beta| \leq m}\right) \quad \text{for } (x, t) \in \mathbb{R}^N \times (0, T); \\ \mathbf{u}(x, 0) &= \mathbf{u}_0(x) \quad \text{for } x \in \mathbb{R}^N. \end{aligned} \quad (1.1)$$

Here, $\partial/\partial x = (\partial/\partial x_1, \dots, \partial/\partial x_N)$ stands for the spatial gradient and $\xi \mapsto \mathbf{P}(x, t, \xi)$ is a polynomial of order $2m$ in the variable $\xi = (\xi_1, \dots, \xi_N) \in \mathbb{R}^N$ (or \mathbb{C}^N); its coefficients are $M \times M$ matrices (real or complex) which are assumed to be real analytic (jointly) in both variables $x \in \mathbb{R}^N$ and $t \in (0, T)$. Also the nonlinearity $(x, t; X) \mapsto \mathbf{f}(x, t; X)$ (a reaction function valued in \mathbb{R}^M or \mathbb{C}^M) is assumed to

be analytic in all variables $x \in \mathbb{R}^N$, $t \in (0, T)$, and $X = (X_\beta)_{|\beta| \leq m} \in \mathbb{R}^{M\tilde{N}}$ (or $\mathbb{C}^{M\tilde{N}}$), where we have substituted $X_\beta = \frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta} \in \mathbb{R}^M$ (or \mathbb{C}^M) for the (mixed) partial derivative of \mathbf{u} with a multi-index $\beta = (\beta_1, \dots, \beta_N) \in (\mathbb{Z}_+)^N$ of order $|\beta| = \beta_1 + \dots + \beta_N$, $|\beta| \leq m$. Here, $\mathbb{Z}_+ = \{0, 1, 2, \dots\}$ and the Euclidean dimension of the m -jet X equals to $M\tilde{N}$ with

$$\tilde{N} = \sum_{k=0}^m \sum_{|\beta|=k} \binom{k}{\beta} = \sum_{k=0}^m N^k \quad \text{where} \quad \binom{k}{\beta} := \frac{k!}{\beta_1! \beta_2! \dots \beta_N!}. \quad (1.2)$$

As usual, \mathbb{R}^N and \mathbb{C}^N , respectively, denote the N -dimensional real and complex Euclidean spaces, $i = \sqrt{-1}$, and $M, N \in \mathbb{N}$ where $\mathbb{N} = \{1, 2, 3, \dots\}$. We have identified $X_\beta = \frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta} \equiv \mathbf{u}(x, t)$ for $\beta = (0, 0, \dots, 0)$ of order $|\beta| = 0$.

As already indicated, we impose certain standard *strong ellipticity* and *analyticity* hypotheses on the coefficients of the partial differential operator $\mathbf{P}(x, t, \frac{1}{i} \frac{\partial}{\partial x})$ and on the reaction function $\mathbf{f}(x, t; X)$ as well. Assuming that $\mathbf{u}_0 \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ ($p > 2 + \frac{N}{m}$) possesses a complex analytic extension to a strip $\mathfrak{X}^{(\kappa_0)}$ of constant width in $\mathbb{C}^N = \mathbb{R}^N + i\mathbb{R}^N$ and the first-order partial derivatives

$$\frac{\partial}{\partial t} \mathbf{f}(x, t; X) \quad \text{and} \quad \frac{\partial}{\partial X_\beta} \mathbf{f}(x, t; X), \quad \text{for } |\beta| \leq m,$$

are locally uniformly bounded for $(x, t; X) \in \mathbb{R}^N \times (0, T) \times \mathbb{R}^{M\tilde{N}}$, in this work we show that the (unique) strict (L^p -type) solution $\mathbf{u} = \mathbf{u}(x, t)$ of problem (1.1) is *real analytic* in $(x, t) \in \mathbb{R}^N \times (0, T)$. Notice that the latter condition (local boundedness of all first-order partial derivatives $\partial \mathbf{f} / \partial X_\beta$) is equivalent with $X \mapsto \mathbf{f}(x, t; X)$ being locally uniformly Lipschitz continuous.

This analyticity claim is motivated by the standard formula for the solution of the Cauchy problem for the heat equation in \mathbb{R}^N (with the Laplace operator Δ , i.e., $\mathbf{P}(x, t, \frac{1}{i} \frac{\partial}{\partial x}) = -\Delta$, $\mathbf{f}(x, t; X) = \mathbf{0}$, and $M = 1$); see e.g. John [50], Chapt. 7, Sect. 1, eq. (1.11), p. 209. The heat equation case has been significantly generalized in Takáč et al. [83, Theorem 2.1, p. 429], where only the leading coefficients of the operator $\mathbf{P}(x, t, \frac{1}{i} \frac{\partial}{\partial x})$ are assumed to be constant, but it is required that $\mathbf{u}_0 \in \mathbf{L}^\infty(\mathbb{R}^N) = [L^\infty(\mathbb{R}^N)]^M$. In our present work, the *analyticity* hypothesis on the initial data \mathbf{u}_0 resembles more to a nonlocal version of the classical Cauchy-Kowalewski theorem (John [50], Chapt. 3, Sect. 3(d), pp. 73–77). We will show that, under this analyticity hypothesis on $\mathbf{u}(\cdot, 0) = \mathbf{u}_0$, if a solution $\mathbf{u} : \mathbb{R}^N \times [0, T) \rightarrow \mathbb{C}^M$ exists, then it must be analytic in $\mathbb{R}^N \times (0, T)$. We are able to specify also the *domain of analyticity* in terms of a *complex analytic extension*. The restriction on the initial data $\mathbf{u}_0 \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$, with the conditions $p > 2 + \frac{N}{m}$ and $s = 2m(1 - \frac{1}{p}) \in (0, 2m)$, allows us to take advantage of (the continuity of) the Sobolev(-Besov) imbedding $B^{s;p,p}(\mathbb{R}^N) \hookrightarrow C^m(\mathbb{R}^N) \cap W^{m,\infty}(\mathbb{R}^N)$; see, e.g., Adams and Fournier [1, Chapt. 7], Theorem 7.34(c), p. 231. This more restrictive condition on the initial data \mathbf{u}_0 enables us to work with an m -jet $X = (X_\beta)_{|\beta| \leq m} \in \mathbb{C}^{M\tilde{N}}$ whose components $X_\beta = \frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta} \in \mathbb{C}^M$ are bounded continuous functions of $(x, t) \in \mathbb{R}^N \times [0, T)$; thus, each $X_\beta(\cdot, t)$ ($|\beta| \leq m$) belongs to $\mathbf{L}^\infty(\mathbb{R}^N)$ at every time $t \in [0, T)$. Consequently, we can apply the Banach fixed point theorem to problem (1.1) in a way similar to [83, Theorem 2.1, p. 429]. For instance, in a typical second-order parabolic problem (i.e., (1.1) with $m = 1$) we can allow for a reaction function $\mathbf{f}(x, t; \mathbf{u}, \frac{\partial \mathbf{u}}{\partial x})$ depending

on \mathbf{u} and its gradient $\partial\mathbf{u}/\partial x$ ($\equiv iD_x\mathbf{u}$), besides the independent variables $x \in \mathbb{R}^N$ and $t \in (0, T)$.

The main contribution of our present article is that we are able to *remove* the hypothesis that the leading coefficients must be constant, in analogy with Takáč [82, Theorem 3.3, p. 59] where the corresponding linear system is treated. In contrast to [83, Proposition A.4, p. 446], this means that we *cannot* calculate the Green function for the Cauchy problem with the leading coefficients only,

$$\begin{aligned} \frac{\partial\mathbf{u}}{\partial t} + (-1)^m \sum_{|\alpha|=2m} \mathbf{P}^{(\alpha)}(x, t) \frac{\partial^{|\alpha|}\mathbf{u}}{\partial x^\alpha} &= \mathbf{0} \quad \text{for } (x, t) \in \mathbb{R}^N \times (0, T); \\ \mathbf{u}(x, 0) &= \mathbf{u}_0(x) \quad \text{for } x \in \mathbb{R}^N, \end{aligned} \quad (1.3)$$

and then simply take advantage of the variation-of-constants formula [83, eq. (3.22), p. 437] to obtain the solution of the original problem (1.1). Fortunately, the methods from [82], based on a priori L^2 -type estimates combined with the Cauchy-Riemann equations, are applicable also to our semilinear system (1.1) provided that already the initial data \mathbf{u}_0 are analytic. Here, each $\mathbf{P}^{(\alpha)}(x, t)$ is an $M \times M$ matrix and recall that $\partial^{|\alpha|}\mathbf{u}/\partial x^\alpha = \frac{\partial^{|\alpha|}\mathbf{u}}{\partial x_1^{\alpha_1} \dots \partial x_N^{\alpha_N}}$ denotes the (mixed) partial derivative of $\mathbf{u} : \mathbb{R}^N \times (0, T) \rightarrow \mathbb{C}^M$ with a multi-index $\alpha = (\alpha_1, \dots, \alpha_N) \in (\mathbb{Z}_+)^N$ of order $|\alpha| = \alpha_1 + \dots + \alpha_N$. This means that, for the semilinear parabolic Cauchy problem (1.1), we do not improve the regularity properties of (in general) nonsmooth initial data to *analytic* regularity as time passes by (for $t \in (0, T)$). We show only that the analytic regularity of the initial data \mathbf{u}_0 (at $t = 0$) is *preserved* for all times $t \in (0, T)$. In contrast, analytic regularity of the initial data is not assumed in [82, 83].

As in [82, 83], our method is based on the simple fact that a function $u : \mathbb{R}^N \times (0, T) \rightarrow \mathbb{R}$ (or \mathbb{C}) is real analytic if and only if it has a *holomorphic* (i.e., complex analytic) extension $\tilde{u} : \Omega \rightarrow \mathbb{C}$ to some complex domain Ω such that $\mathbb{R}^N \times (0, T) \subset \Omega \subset \mathbb{C}^N \times \mathbb{C}$, i.e., $u = \tilde{u}|_{\mathbb{R}^N \times (0, T)}$, the restriction of \tilde{u} to $\mathbb{R}^N \times (0, T)$. If the domain Ω is fixed then the holomorphic extension \tilde{u} of u to Ω is always unique, see e.g. John [50], Chapt. 3, Sect. 3(c), pp. 70–72. Thus, in order to show that the weak solution $\mathbf{u} = \mathbf{u}(x, t)$ of problem (1.1) is real analytic in $\mathbb{R}^N \times (0, T)$, it suffices to construct a holomorphic extension $\tilde{\mathbf{u}}$ of \mathbf{u} to some complex domain Ω ($\mathbb{R}^N \times (0, T) \subset \Omega \subset \mathbb{C}^N \times \mathbb{C}$). Because of the uniqueness (of a holomorphic extension), we often drop the tilde “ $\tilde{}$ ” in the notation for the (unique) holomorphic extension. Analogous ideas (holomorphic extension, uniqueness, and Bergman and Szegő spaces of holomorphic functions) were used earlier in Hayashi [35, 36, 37, 38].

Instead of using the Green function method (cf. [83]), we establish the existence of solutions to the Cauchy problem (1.1) in a complex parabolic domain $\mathfrak{X}^{(r)} \times [0, T)$ in $\mathbb{C}^N \times \mathbb{C}$ with initial data \mathbf{u}_0 from a space of holomorphic functions whose domain $\mathfrak{X}^{(r)} = \mathbb{R}^N + iQ^{(r)}$ is a *tube* in \mathbb{C}^N with *base* $Q^{(r)} = (-r, r)^N$, for some $0 < r < \infty$, see Takáč [82, (21), p. 58]. The (complex) analyticity in space is then verified by means of the Cauchy-Riemann equations, whereas the (complex) analyticity in time is obtained from the properties of holomorphic semigroups in the Besov space $\mathbf{B}^{s;p,p}(\mathbb{R}^N) = [B^{s;p,p}(\mathbb{R}^N)]^M$. Our use of the Cauchy-Riemann equations already at the initial time $t = 0$ requires that \mathbf{u}_0 be (complex) *analytic* in $\mathfrak{X}^{(r)}$.

To provide a quick, nontechnical hint to our approach, we now give an illustrative weaker version of our main result, Theorem 3.4 in Section 3, for a single equation

in one space dimension ($M = N = 1$),

$$\begin{aligned} \frac{\partial u}{\partial t} &= a(x, t) \frac{\partial^2 u}{\partial x^2} + b(x, t) \frac{\partial u}{\partial x} + c(x, t)u + f(x, t; u, \frac{\partial u}{\partial x}) \quad \text{for } (x, t) \in \mathbb{R}^1 \times (0, T); \\ u(x, 0) &= u_0(x) \quad \text{for } x \in \mathbb{R}^1. \end{aligned} \tag{1.4}$$

We begin with the complexifications of the spatial and temporal variables, $x \in \mathbb{R}^1$ and $t \in (0, T)$, respectively: Given any real numbers $0 < r < \infty$ and $0 < T' \leq T < \infty$, we introduce the complex domains

$$\begin{aligned} \mathfrak{X}^{(r)} &:= \{z = x + iy \in \mathbb{C} : |y| < r\} = \mathbb{R} + i(-r, r), \\ \Delta_\vartheta &:= \{t = \varrho e^{i\theta} \in \mathbb{C} : \varrho > 0 \text{ and } \theta \in (-\vartheta, \vartheta)\}, \quad \vartheta = \arctan(r/T'), \end{aligned} \tag{1.5}$$

$$\begin{aligned} \Delta_\vartheta^{(T')} &:= \Delta_\vartheta \cap \{t \in \mathbb{C} : 0 < \Re t < T'\} \\ &= \{t = \varrho e^{i\theta} \in \mathbb{C} : |\theta| < \vartheta \text{ and } 0 < \varrho < T'/\cos\theta\}, \end{aligned} \tag{1.6}$$

$$\begin{aligned} \Delta_\vartheta^{T',T} &:= \Delta_\vartheta^{(T)} \cap \{t \in \mathbb{C} : |\Im t| < T' \cdot \tan \vartheta\} = \cup_{0 \leq \xi \leq T-T'} (\xi + \Delta_\vartheta^{(T')}) \\ &= \cup_{0 \leq \xi \leq T-T'} \{\xi + t' \in \mathbb{C} : t' \in \Delta_\vartheta^{(T')}\} = [0, T - T'] + \Delta_\vartheta^{(T')} \end{aligned} \tag{1.7}$$

with the angle $\vartheta \in (0, \pi/2)$ given by $\tan \vartheta = r/T'$. Of course, if $T = T'$ then $\Delta_\vartheta^{T',T} = \Delta_\vartheta^{(T')}$ is an open triangle. Clearly, we have

$$\Delta_\vartheta^{T',T} = \cup_{0 < r \leq T'} \mathfrak{F}_{r \cdot \cot \vartheta, T}^{(r)} = \cup_{0 < r \leq T'} [(r \cdot \cot \vartheta, T) + i(-r, r)],$$

where

$$\mathfrak{F}_{T',T}^{(r)} := \{t = \sigma + i\tau \in \mathbb{C} : T' < \sigma < T \text{ and } |\tau| < r\} = (T', T) + i(-r, r). \tag{1.8}$$

We set $\mathfrak{F}_{0,T}^{(r)} = (0, T) + i(-r, r)$ if $T' = 0$. The closures in \mathbb{C} of $\mathfrak{X}^{(r)}$, Δ_ϑ , $\Delta_\vartheta^{(T')}$, $\Delta_\vartheta^{T',T}$, and $\mathfrak{F}_{T',T}^{(r)}$ are denoted by $\bar{\mathfrak{X}}^{(r)}$, $\bar{\Delta}_\vartheta$, $\bar{\Delta}_\vartheta^{(T')}$, $\bar{\Delta}_\vartheta^{T',T}$, and $\bar{\mathfrak{F}}_{T',T}^{(r)}$, respectively.

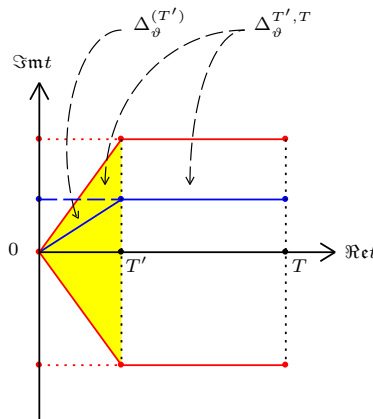


FIGURE 1. Triangle Δ_ϑ starting at the origin has been defined in (1.5). Its shifts to the right create the region $\Delta_\vartheta^{(T')}$ in (1.6)

The Banach space of all continuous ($B^{s;p,p}(\mathbb{R})$ -valued) functions $u : [0, T] \rightarrow B^{s;p,p}(\mathbb{R})$ is denoted by $C([0, T] \rightarrow B^{s;p,p}(\mathbb{R}))$; it is endowed with the natural supremum norm

$$\|u\|_{L^\infty(0,T)} := \sup_{t \in [0,T]} \|u(\cdot, t)\|_{B^{s;p,p}(\mathbb{R})} < \infty.$$

Theorem 1.1 ($M = N = 1$). *Let $p > 2 + \frac{1}{m}$, $s = 2m(1 - \frac{1}{p})$, and $0 < T < \infty$. Assume that there are constants $A, B > 0$ such that all coefficients a, b , and c and the partial derivative $\partial a / \partial x$ are bounded, continuously differentiable functions in the Cartesian product $\bar{\mathfrak{X}}^{(A)} \times \bar{\mathfrak{T}}_{0,T}^{(B)}$, with $\Re ca \geq \text{const} > 0$, and all a, b , and c are holomorphic in $\mathfrak{X}^{(A)} \times \mathfrak{T}_{0,T}^{(B)}$. Furthermore, let us assume that the first-order time derivatives of all functions a, b, c , and $\partial a / \partial x$ are bounded in $\bar{\mathfrak{X}}^{(A)} \times \bar{\mathfrak{T}}_{0,T}^{(B)}$. Finally, assume that f is holomorphic in $\mathfrak{X}^{(A)} \times \mathfrak{T}_{0,T}^{(B)} \times \mathbb{C}^2$, $f = f(x, t; u, \eta)$ where $\eta = \partial u / \partial x$, with all functions $f, \partial f / \partial t, \partial f / \partial u$, and $\partial f / \partial \eta$ being locally bounded in $\mathfrak{X}^{(A)} \times \mathfrak{T}_{0,T}^{(B)} \times \mathbb{C}^2$, and it satisfies*

$$\int_{-\infty}^{\infty} |f(x + iy, t; 0, 0)|^p dx \leq K^p \quad (K = \text{const} < \infty)$$

for all $y \in [-A, A]$ and $t \in \bar{\mathfrak{T}}_{0,T}^{(B)}$.

(i) *Given any $u_0 \in B^{s;p,p}(\mathbb{R})$, the Cauchy problem (1.4) possesses a unique weak solution $u \in C([0, T_0] \rightarrow B^{s;p,p}(\mathbb{R}))$ defined on a (possibly shorter) time interval $[0, T_0] \subset [0, T]$ of some positive length $T_0 \in (0, T]$.*

(ii) *Furthermore, if $u_0 : \mathbb{R}^1 \rightarrow \mathbb{C}$ possesses a (unique) holomorphic extension to a complex strip $\mathfrak{X}^{(r_0)} \subset \mathbb{C}$, $0 < r_0 \leq A$, denoted by $u_0 : \mathfrak{X}^{(r_0)} \rightarrow \mathbb{C}$ again, such that*

$$\mathfrak{N}^{(r_0)}(u_0) := \sup_{y \in [-r_0, r_0]} \|u_0(\cdot + iy)\|_{B^{s;p,p}(\mathbb{R})} < \infty$$

holds (cf. (3.10) below), then also any (global) weak solution $u \in C([0, T] \rightarrow B^{s;p,p}(\mathbb{R}))$ can be (uniquely) extended to a holomorphic function in $\mathfrak{X}^{(A')} \times \Delta_{\vartheta}^{T',T}$, denoted again by $u(x + iy, t)$, where all numbers $A' \in (0, A]$, $T' \in (0, T]$, and $\vartheta \in (0, \pi/2)$ are sufficiently small, $T' \cdot \tan \vartheta \leq B$, and $u(\cdot + iy, \cdot) : t \mapsto u(\cdot + iy, t) : \bar{\Delta}_{\vartheta}^{T',T} \rightarrow B^{s;p,p}(\mathbb{R})$ is continuous for every fixed $y \in [-r_0, r_0]$ together with

$$\sup_{t \in \Delta_{\vartheta}^{T',T}} \mathfrak{N}^{(A')}(u(\cdot + iy, t)) = \sup_{(y,t) \in [-r_0, r_0] \times \Delta_{\vartheta}^{T',T}} \|u(\cdot + iy, t)\|_{B^{s;p,p}(\mathbb{R}^N)} < \infty.$$

In particular, the extension u is holomorphic in $\mathfrak{X}^{(A')} \times \mathfrak{T}_{T',T}^{(B')}$ with $B' = T' \cdot \tan \vartheta \leq B$, where $T' > 0$ and $\vartheta > 0$ are small enough.

We remark that $\mathfrak{T}_{T',T}^{(B')} \subset \Delta_{\vartheta}^{T',T} \subset \mathfrak{T}_{0,T}^{(B)}$, by $B' = T' \cdot \tan \vartheta \leq B$. If $0 < T_0 < T$ in Part (i) of this theorem, then we have to replace T by $T = T_0$ in part (ii). Notice that the condition that both (continuous) partial derivatives $\partial f / \partial u$ and $\partial f / \partial \eta$ are locally bounded in $\mathfrak{X}^{(A)} \times \mathfrak{T}_{0,T}^{(B)} \times \mathbb{C}^2$ is equivalent with $(u, \eta) \mapsto f(x, t; u, \eta)$ being locally Lipschitz continuous.

Remark 1.2. It follows easily from Theorem 1.1 that every weak solution $u \in C([0, T] \rightarrow B^{s;p,p}(\mathbb{R}))$ to the Cauchy problem (1.4) is classical in the sense that it is of class C^∞ over the open set $\mathbb{R} \times (0, T)$ and satisfies (1.4) pointwise and the initial condition $u(\cdot, 0) = u_0 \in B^{s;p,p}(\mathbb{R})$ in the $B^{s;p,p}(\mathbb{R})$ -limit, i.e., $\|u(\cdot, t) - u_0\|_{B^{s;p,p}(\mathbb{R})} \rightarrow$

0 as $t \rightarrow 0+$. For the special case of the reaction function $f : (u, \eta) \mapsto f(x, t; u, \eta)$ being linear, a *weak* solution $u \in C([0, T] \rightarrow L^2(\mathbb{R}))$ to the (linear) Cauchy problem (1.4) is defined e.g. in Evans [26], Chapt. 7, §1.1, p. 352, or Lions [62], Chapt. IV, §1, p. 44, or [63], Chapt. III, (1.11), p. 102. In that case the initial condition $u(\cdot, 0) = u_0 \in L^2(\mathbb{R})$ holds in the sense of the $L^2(\mathbb{R})$ -limit, $\|u(\cdot, t) - u_0\|_{L^2(\mathbb{R})} \rightarrow 0$ as $t \rightarrow 0+$. The main reason why we prefer to work with the notion of a *weak solution* as opposed to a *classical solution* of the Cauchy problem (1.4) is the fact that already a weak solution is *unique*. The uniqueness of a weak solution is an important technical argument in our proofs of Theorem 1.1 and Theorem 3.4 (Section 3).

In fact, we work sometimes also with the so called *mild solutions* to the Cauchy problem (1.4) that make sense in $C([0, T] \rightarrow L^2(\mathbb{R}))$; cf. Takáč [82, Sect. 3 and 4], even though we use them in the Besov space $B^{s;p,p}(\mathbb{R})$ in place of $L^2(\mathbb{R})$. Mild solutions do not require any additional regularity knowledge as they are defined by the well known variation-of-constants formula (Pazy [72, §5.7, p. 168]). Thus, they are even “weaker” than the weak solutions, but in our situation one can easily verify that every mild solution is also a weak solution to problem (1.4) and vice versa; see e.g. Ball [10] (or [72, Theorem on p. 259]).

The same remarks apply also to the more general Cauchy problem (1.1).

This article is organized as follows. We introduce some basic notation (mostly complex domains) in Section 2. Our main analyticity result, Theorem 3.4, supplemented by an additional explanation in Proposition 3.5, is stated in Section 3. Their proofs are gradually built up in Sections 4 through 8: First, the Cauchy problem in $\mathbb{R}^N \times (0, T)$ is treated as an abstract initial value problem in Section 4. There, an important abstract a priori $B^{s;p,p}$ -type estimate is established in Theorem 4.7. Analyticity in time for this abstract problem is proved in Section 5 (Theorem 5.3). Then, in Section 6, we treat analyticity in space for the semilinear parabolic Cauchy problem in $\mathbb{R}^N \times (0, T)$, see Proposition 6.5, provided the initial data are already analytic (in space). We show that the analyticity in space is preserved for all times in $[0, T]$. We combine the time and space analyticity results from Sections 5 and 6 into Theorem 7.1 in Section 7. This theorem is still only “local” in time. Our main analyticity result, Theorem 3.4, is proved in Section 8, together with Proposition 3.5 and, in particular, the “regularity” estimates (3.13) and (3.14). Section 9 treats an application to a *Risk Model in Mathematical Finance*. Finally, Section 10 contains some historical remarks and comments concerning the analyticity of solutions to linear and semilinear elliptic and parabolic systems and its applications to relevant classical problems.

2. NOTATION

We stick to the classical notation $\mathbb{N} = \{1, 2, 3, \dots\}$, $\mathbb{Z}_+ = \{0, 1, 2, 3, \dots\} = \mathbb{N} \cup \{0\}$, and $\mathbb{Z} = \{0, \pm 1, \pm 2, \pm 3, \dots\} = \mathbb{Z}_+ \cup (-\mathbb{Z}_+)$, together with $\mathbb{R} = (-\infty, \infty)$, $\mathbb{R}_+ = [0, \infty)$, and $\mathbb{C} = \mathbb{R} + i\mathbb{R} \cong \mathbb{R}^2$. Typically, we denote by $x = (x_1, x_2, \dots, x_N)$ and $y = (y_1, y_2, \dots, y_N)$ points in \mathbb{R}^N and by $z = (z_1, z_2, \dots, z_N)$ points in \mathbb{C}^N . We often write $\zeta = \xi + i\eta$ for $\zeta \in \mathbb{C}$ and $\xi, \eta \in \mathbb{R}$, i.e., $\Re \zeta = \xi$ and $\Im \zeta = \eta$ are the real and imaginary parts of $\zeta \in \mathbb{C}$, respectively. Similarly, $z = x + iy$ for $z \in \mathbb{C}^N$ and $x, y \in \mathbb{R}^N$, or equivalently $z_i = x_i + iy_i$ ($i = 1, 2, \dots, N$) for $z_i \in \mathbb{C}$ and $x_i, y_i \in \mathbb{R}$, i.e., $\Re z = x$ and $\Im z = y$. Hence, we identify $\mathbb{C}^N = \mathbb{R}^N \oplus i\mathbb{R}^N$ (or simply $\mathbb{C}^N = \mathbb{R}^N + i\mathbb{R}^N$) as vector spaces over the field \mathbb{R} and thus consider \mathbb{R}^N to be a (vector) subspace of \mathbb{C}^N . We use a bar ($\bar{}$) to denote the complex

conjugate $\bar{\zeta}$ of a number $\zeta \in \mathbb{C}$. The complex conjugate of a vector $z \in \mathbb{C}^N$ is denoted by $\bar{z} = (\bar{z}_1, \bar{z}_2, \dots, \bar{z}_N)$. Similarly, the complex conjugate function of a complex-valued function $f(z)$ (for $f : \mathbb{C}^N \rightarrow \mathbb{C}$, for instance) is denoted by $\bar{f}(z) \equiv f(\bar{z})$. Furthermore, we denote by $(z, w) = \sum_{i=1}^N z_i \bar{w}_i$ the standard Euclidean inner product of $z, w \in \mathbb{C}^N$ and by $|z| = \left(\sum_{i=1}^N |z_i|^2\right)^{1/2}$ the induced (Euclidean) norm of $z \in \mathbb{C}^N$. We will often use the sum (ℓ^1 -) and the maximum (ℓ^∞ -) norms of $z \in \mathbb{C}^N$, respectively:

$$|z|_1 = \sum_{i=1}^N |z_i| \quad \text{and} \quad |z|_\infty = \max_{1 \leq i \leq N} |z_i|.$$

Finally, we write $z \cdot w = \sum_{i=1}^N z_i w_i$ for the *bilinear* product of $z, w \in \mathbb{C}^N$, which is not to be confused with the inner product $(z, w) = \sum_{i=1}^N z_i \bar{w}_i$ if $w \notin \mathbb{R}^N$ (which is *sesquilinear*). The Euclidean (ℓ^2 -) norm of $z \in \mathbb{C}^N$ is abbreviated as $|z| \equiv |z|_2 = \sqrt{(z, z)} = \left(\sum_{i=1}^N |z_i|^2\right)^{1/2}$.

The vector space (over the field \mathbb{R}) of all real-valued (square) $M \times M$ matrices $\mathbf{A} = (a_{ij})_{i,j=1}^M$ is denoted by $\mathbb{R}^{M \times M}$. Similarly, the vector space (over the field \mathbb{C}) of all complex-valued $M \times M$ matrices \mathbf{A} is denoted by $\mathbb{C}^{M \times M}$.

Given $r \in (0, \infty)$, we denote by $Q^{(r)} = (-r, r)^N = \{y \in \mathbb{R}^N : |y|_\infty < r\}$ the N -dimensional open cube in \mathbb{R}^N (centered at the origin) with side lengths $2r$, and by $\bar{Q}^{(r)} = [-r, r]^N$ its closure.

To formulate our main hypotheses, given $r, T \in (0, \infty)$ and $T' \in [0, T)$, we introduce the following complex domains for the complexifications of the spatial and temporal variables, $x \in \mathbb{R}^N$ and $t \in (0, T)$, respectively:

$$\begin{aligned} \mathfrak{X}^{(r)} &:= \{z = x + iy \in \mathbb{C}^N : |y|_\infty < r\} = \mathbb{R}^N + iQ^{(r)}, \\ \mathfrak{T}_{T', T}^{(r)} &:= \{t = \sigma + i\tau \in \mathbb{C} : T' < \sigma < T \text{ and } |\tau| < r\}; \end{aligned} \tag{2.1}$$

see (1.8).

The former, $\mathfrak{X}^{(r)}$, is a *tube* (often called a *strip*) in \mathbb{C}^N with *base* $Q^{(r)}$ and the latter, $\mathfrak{T}_{T', T}^{(r)}$, is an open rectangle in the complex plane \mathbb{C} . Notice that $\mathfrak{T}_{T', T}^{(r)}$ contains the interval (T', T) . The remaining temporal domains, Δ_ϑ and $\Delta_\vartheta^{(T')}$, have been introduced in (1.5) and (1.6), respectively, with the angle $\vartheta \in (0, \pi/2)$ given by $\tan \vartheta = r/T'$.

Our techniques will use holomorphic semigroups in an open *sector* $\Delta_\vartheta \subset \mathbb{C}$ defined in (1.5), with a given angle $\vartheta \in (0, \pi/2)$, but often locally in time in an open *triangle* $\Delta_\vartheta^{(T')} \subset \mathbb{C}$ defined in (1.6), where $0 < T' < \infty$. Their respective closures in \mathbb{C} are denoted by $\bar{\Delta}_\vartheta$ and $\bar{\Delta}_\vartheta^{(T')}$; both contain the origin $0 \in \mathbb{C}$. Finally, for $0 < T' \leq T < \infty$ we recall the definition of the temporal domain $\Delta_\vartheta^{T', T}$ introduced in (1.7) with the angle $\vartheta \in (0, \pi/2)$ given by $\tan \vartheta = r/T'$. Its closure in \mathbb{C} is denoted by $\bar{\Delta}_\vartheta^{T', T}$. Clearly, $\xi + \Delta_\vartheta^{(T')} = \{\xi + t' \in \mathbb{C} : t' \in \Delta_\vartheta^{(T')}\}$. Of course, if $T = T'$ then $\Delta_\vartheta^{T', T} = \Delta_\vartheta^{(T')}$ is an open triangle.

Throughout this article we work with complex-valued functions; hence, all Banach and Hilbert spaces of functions we consider are complex (over the field \mathbb{C}). We work with the standard inner product in $L^2(\mathbb{R}^N)$ defined by $(u, v)_{L^2} := \int_{\mathbb{R}^N} u \bar{v} \, dx$ for $u, v \in L^2(\mathbb{R}^N)$. The induced norm is abbreviated by $\|u\|_{L^2} \equiv \|u\|_{L^2(\mathbb{R}^N)}$. We warn the reader that we identify the dual space $\mathcal{H}' = (L^2(\mathbb{R}^N))'$ of the complex

Hilbert space $\mathcal{H} = L^2(\mathbb{R}^N)$ with \mathcal{H} itself by means of the (complex) Riesz representation theorem which yields an *anti-linear* isomorphism of \mathcal{H} onto \mathcal{H}' (cf. Adams and Fournier [1, Chapt. 2], Theorem 2.44, p. 47).

The following notation is taken from Krantz [57, Chapt. 0]. Given a domain Ω in \mathbb{R}^r , we denote by $C^k(\Omega)$ ($k \in \mathbb{Z}_+$) the vector space of all k -times continuously differentiable functions $f : \Omega \rightarrow \mathbb{C}$ and by $C^k(\tilde{\Omega})$ the vector space of all $f : \tilde{\Omega} \rightarrow \mathbb{C}$ such that $f|_{\Omega} \in C^k(\Omega)$ and each partial derivative $\frac{\partial^{|\alpha|} f}{\partial x^\alpha}$ of f ($\alpha \in (\mathbb{Z}_+)^r$) of order $|\alpha| \leq k$ can be extended to a continuous function on $\tilde{\Omega}$. Of course, $f|_{\Omega}$ stands for the restriction of f to Ω and all partial derivatives are taken in the real variable sense ($x \in \mathbb{R}^r$). In case $\Omega \subset \mathbb{C}^r = \mathbb{R}^r \oplus i\mathbb{R}^r \cong \mathbb{R}^{2r}$ ($r \in \mathbb{N}$) is a complex domain, the (mixed) partial derivative $\frac{\partial^{|\alpha|+|\beta|} f}{\partial x^\alpha \partial y^\beta}$ of f ($\alpha, \beta \in (\mathbb{Z}_+)^r$) of order $|\alpha| + |\beta| \leq k$ is taken in the real variable sense, where $x, y \in \mathbb{R}^r$ in $(x, y) \cong z = x + iy \in \mathbb{C}^r$. The vector spaces $C^k(\Omega)$ and $C^k(\tilde{\Omega})$ are defined analogously, with $|\alpha| + |\beta| \leq k$. If Ω is bounded then $C^k(\tilde{\Omega})$ is a Banach space endowed with a maximum-type norm. If Ω is not bounded in general then we denote by $C_{\text{unif}}^0(\tilde{\Omega})$ the vector space of all uniformly continuous functions $f : \tilde{\Omega} \rightarrow \mathbb{C}$, and by $C_{\text{bdd}}^0(\Omega) = C^0(\Omega) \cap L^\infty(\Omega)$ the vector space of all bounded continuous functions $f : \Omega \rightarrow \mathbb{C}$. Alternatively, if Ω is not bounded (in \mathbb{R}^r or \mathbb{C}^r endowed with the Euclidean metric d) then we may consider a compactification $\tilde{\Omega}$ of Ω , i.e., a compact metric space $\tilde{\Omega}$ (endowed with a metric \tilde{d}) such that there is a homeomorphism $j : \Omega \rightarrow \tilde{\Omega}$ of Ω onto a dense subset $j(\Omega)$ of $\tilde{\Omega}$, such that for any pair of sequences $\{x_n\}_{n=1}^\infty, \{y_n\}_{n=1}^\infty \subset \Omega$ we have $d(x_n, y_n) \rightarrow 0$ if and only if $\tilde{d}(j(x_n), j(y_n)) \rightarrow 0$ as $n \rightarrow \infty$. Clearly, by identifying Ω with the subset $j(\Omega)$ of $\tilde{\Omega}$ we identify Ω with a dense subset of $\tilde{\Omega}$. Hence, we can identify $C^0(\tilde{\Omega})$ with a vector subspace of $C_{\text{bdd}}^0(\Omega) \cap C_{\text{unif}}^0(\Omega)$. As a simple example, we may take $\tilde{\Omega}$ to be the one-point compactification of a domain $\Omega \subset \mathbb{R}^r$ (or $\Omega \subset \mathbb{C}^r$). In particular, we have $\tilde{\mathbb{R}}^r = \mathbb{R}^r \cup \{\infty\}$ and $\tilde{\mathbb{C}}^r = \mathbb{C}^r \cup \{\infty\}$ endowed with the metric \tilde{d} defined in the next section (Section 3, (3.2)).

Finally, if $\Omega \subset \mathbb{C}^r$ is a complex domain, we denote by $A(\Omega)$ the Fréchet space of all holomorphic functions $f : \Omega \rightarrow \mathbb{C}$ endowed with the (complete metrizable) topology of uniform convergence on compact subsets of Ω . As usual, we abbreviate the *Cauchy-Riemann* partial differential operators

$$\frac{\partial}{\partial z_i} := \frac{1}{2} \left(\frac{\partial}{\partial x_i} - i \frac{\partial}{\partial y_i} \right) \quad \text{and} \quad \frac{\partial}{\partial \bar{z}_i} := \frac{1}{2} \left(\frac{\partial}{\partial x_i} + i \frac{\partial}{\partial y_i} \right). \quad (2.2)$$

Remark 2.1. We will often use the following classical fact; see e.g. John [50, Theorem, p. 70] or Krantz [57, Definition II, p. 3]: Let $\Omega \subset \mathbb{C}^N$ be a complex domain ($N \geq 1$). A continuously differentiable function $h : \Omega \rightarrow \mathbb{C}$ (in the real variable sense, partially with respect to $x_i, y_i \in \mathbb{R}; i = 1, 2, \dots, N$) is holomorphic (i.e., complex analytic) *if and only if* it verifies the *Cauchy-Riemann* equations in Ω , i.e., $\partial h / \partial \bar{z}_i = 0$ in Ω ; for $i = 1, 2, \dots, N$.

3. STATEMENT OF MAIN RESULTS

Let us abbreviate the differential operators (and derivatives)

$$D_x := \frac{1}{i} \frac{\partial}{\partial x} = \left(\frac{1}{i} \frac{\partial}{\partial x_i} \right)_{i=1}^N, \quad \partial_t := \frac{\partial}{\partial t},$$

$$D_x^\alpha := i^{-|\alpha|} \frac{\partial^{|\alpha|}}{\partial x^\alpha} = i^{-|\alpha|} \frac{\partial^{|\alpha|}}{\partial x_1^{\alpha_1} \dots \partial x_N^{\alpha_N}} \quad \text{for } \alpha = (\alpha_i)_{i=1}^N \in (\mathbb{Z}_+)^N.$$

We assume that the operator

$$\begin{aligned} \mathbf{P}(x, t, D_x) &= \sum_{|\alpha|, |\beta| \leq m} D_x^\alpha \left(\mathbf{P}^{\alpha\beta}(x, t) D_x^\beta \right) \\ &\equiv \sum_{|\alpha|, |\beta| \leq m} i^{-|\alpha| - |\beta|} \frac{\partial^{|\alpha|}}{\partial x^\alpha} \left(\mathbf{P}^{\alpha\beta}(x, t) \frac{\partial^{|\beta|}}{\partial x^\beta} \right), \end{aligned} \quad (3.1)$$

for $(x, t) \in \mathbb{R}^N \times (0, T)$, is a linear partial differential operator of order $2m$ in *divergence form* with the coefficients $i^{-|\alpha| - |\beta|} \mathbf{P}^{\alpha\beta}(x, t)$ indexed by $\alpha, \beta \in (\mathbb{Z}_+)^N$ with $|\alpha| \leq m$ and $|\beta| \leq m$, where each $\mathbf{P}^{\alpha\beta}(x, t) = (P_{jk}^{\alpha\beta})_{j,k=1}^M$ is an $M \times M$ matrix with real (or complex) entries $P_{jk}^{\alpha\beta} = P_{jk}^{\alpha\beta}(x, t)$. The reader is referred to Friedman [32, Part. 1, Sect. 12, pp. 32–37] or John [50, Chapt. 6, Sect. 2, pp. 190–195] for general facts about such operators.

Let us abbreviate the product domain $\Omega = \mathfrak{X}^{(r_0)} \times \Delta_{\vartheta_0}^{T_0, T} \subset \mathbb{C}^N \times \mathbb{C}$, with some $r_0 \in (0, \infty)$, $0 < T_0 \leq T < \infty$, and $\vartheta_0 \in (0, \pi/2)$; the closure of Ω in $\mathbb{C}^N \times \mathbb{C}$ is denoted by $\bar{\Omega}$. We introduce also a compactification $\tilde{\Omega}$ of Ω ,

$$\tilde{\Omega} = \tilde{\mathfrak{X}}^{(r_0)} \times \bar{\Delta}_{\vartheta_0}^{T_0, T} = (\tilde{\mathbb{R}}^N + i\bar{Q}^{(r)}) \times \bar{\Delta}_{\vartheta_0}^{T_0, T} \cong \tilde{\mathbb{R}}^N \times \bar{Q}^{(r)} \times \bar{\Delta}_{\vartheta_0}^{T_0, T},$$

where $\tilde{\mathbb{R}}^N = \mathbb{R}^N \cup \{\infty\}$ denotes the one-point compactification of \mathbb{R}^N endowed with the standard metric

$$\tilde{d}(x, y) = \tilde{d}(y, x) = \begin{cases} \frac{|x-y|}{1+|x-y|} & \text{if } x, y \in \mathbb{R}^N; \\ 1 & \text{if } x \in \mathbb{R}^N, y = \infty; \\ 0 & \text{if } x = y = \infty. \end{cases} \quad (3.2)$$

Hence, $\tilde{\mathfrak{X}}^{(r_0)} = \tilde{\mathbb{R}}^N + i\bar{Q}^{(r)} \cong \tilde{\mathbb{R}}^N \times \bar{Q}^{(r)}$ is a compactification of $\mathfrak{X}^{(r_0)} \subset \mathbb{C}^N$.

We assume that the operator \mathbf{P} and the function \mathbf{f} satisfy the following hypotheses in the product domain $\Omega = \mathfrak{X}^{(r_0)} \times \Delta_{\vartheta_0}^{T_0, T}$ defined in (1.7) and (2.1).

3.1. Hypothesis.

- (H1) For each pair $\alpha, \beta \in (\mathbb{Z}_+)^N$ with $|\alpha| \leq m$ and $|\beta| \leq m$, the entries $P_{jk}^{\alpha\beta} : \bar{\Omega} \rightarrow \mathbb{C}$ ($j, k = 1, 2, \dots, M$) of the coefficient $\mathbf{P}^{\alpha\beta} = (P_{jk}^{\alpha\beta})_{j,k=1}^M$ belong to $C^1(\bar{\Omega}) \cap L^\infty(\Omega) \cap A(\Omega)$. Moreover, we assume that also all partial derivatives $\frac{\partial^{|\alpha'|}}{\partial x^{\alpha'}} P_{jk}^{\alpha\beta}(x, t)$ of order $|\alpha'| \leq |\alpha|$ ($\alpha' \in (\mathbb{Z}_+)^N$) are in $C^1(\bar{\Omega})$. The entries $P_{jk}^{\alpha\beta}$ of the *leading coefficients* ($|\alpha| = |\beta| = m$) are assumed to belong also to $C_{\text{unif}}^0(\bar{\Omega})$ besides being in $C^1(\bar{\Omega}) \cap L^\infty(\Omega) \cap A(\Omega)$.

This is the case if the entries $P_{jk}^{\alpha\beta}$ of the leading coefficients ($|\alpha| = |\beta| = m$) belong also to $C^0(\bar{\Omega})$ besides being in $C^1(\bar{\Omega}) \cap L^\infty(\Omega) \cap A(\Omega)$. This claim follows directly from $C^0(\bar{\Omega}) \subset C_{\text{unif}}^0(\bar{\Omega})$ which means that any continuous function $f : \bar{\Omega} \rightarrow \mathbb{C}$ is uniformly continuous on $\bar{\Omega}$ with the restriction $f|_\Omega$ to Ω being uniformly continuous and, thus, $f|_\Omega : \Omega \rightarrow \mathbb{C}$ possesses a (unique) continuous extension $\tilde{f} : \bar{\Omega} \rightarrow \mathbb{C}$ to $\bar{\Omega}$ which turns out to be uniformly continuous, as well.

(H2) Operator \mathbf{P} is *strongly elliptic* in $\bar{\Omega}$, i.e., there exists a constant $c \in (0, \infty)$ such that the inequality

$$\Re \left(\sum_{j,k=1}^M \sum_{|\alpha|=|\beta|=m} P_{jk}^{\alpha\beta}(z, t) \xi^{\alpha+\beta} \eta_k \bar{\eta}_j \right) \geq c |\xi|^{2m} |\boldsymbol{\eta}|^2 \tag{3.3}$$

holds for all $(z, t) \in \bar{\Omega}$ and for all $\xi = (\xi_1, \dots, \xi_N) \in \mathbb{R}^N$ and $\boldsymbol{\eta} = (\eta_1, \dots, \eta_M) \in \mathbb{C}^M$, where $\xi^{\alpha+\beta} = \xi_1^{\alpha_1+\beta_1} \dots \xi_N^{\alpha_N+\beta_N}$ and $\alpha = (\alpha_1, \dots, \alpha_N)$ is in $(\mathbb{Z}_+)^N$, and $\beta = (\beta_1, \dots, \beta_N)$ is in $(\mathbb{Z}_+)^N$.

(H3) The components $f_j : \bar{\Omega} \times \mathbb{C}^{M\tilde{N}} \rightarrow \mathbb{C}$ ($j = 1, 2, \dots, M$) of the reaction function $\mathbf{f} = (f_1, \dots, f_M)$ belong to

$$C^1(\bar{\Omega} \times \mathbb{C}^{M\tilde{N}}) \cap A(\Omega \times \mathbb{C}^{M\tilde{N}}).$$

(Recall that the integer \tilde{N} is defined in (1.2).) Moreover, we assume that, for every bounded subset $\Sigma \subset \mathbb{C}^{M\tilde{N}}$, their first-order time derivatives $\frac{\partial}{\partial t} f_j(x, t; X)$ together with their first-order partial derivatives

$$\frac{\partial}{\partial X_{\beta,k}} f_j(x, t; X), \quad \text{for } |\beta| \leq m \text{ and } j, k = 1, 2, \dots, M,$$

with respect to the components $X_{\beta,k}$ of the vector $X_\beta = (X_{\beta,1}, \dots, X_{\beta,M}) = \frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta} \in \mathbb{C}^M$ (or \mathbb{C}^M) are *uniformly bounded* on the set $\Omega \times \Sigma$. Finally, we assume that the function $\mathbf{f} : \bar{\Omega} \times \mathbb{C}^{M\tilde{N}} \rightarrow \mathbb{C}^M$ satisfies

$$\int_{\mathbb{R}^N} |\mathbf{f}(x + iy, t; \vec{\mathbf{0}})|^p dx \leq K^p \quad \text{for all } y \in \bar{Q}^{(r_0)} \text{ and } t \in \bar{\Delta}_{\vartheta_0}^{T_0, T}, \tag{3.4}$$

where $K \in (0, \infty)$ is a constant and $\vec{\mathbf{0}} := (0)_{|\beta| \leq m} \equiv (0, \dots, 0) \in \mathbb{C}^{M\tilde{N}}$.

Remark 3.1. The local boundedness condition in (H3) on the first-order partial derivatives $\frac{\partial}{\partial t} f_j(x, t; X)$ and $\frac{\partial f_j}{\partial X_{\beta,k}}(x, t; X)$ on the set $\Omega \times \Sigma$ will be used later (cf. (6.8) in Section 6) in the following equivalent form:

For a bounded subset $\Sigma \subset \mathbb{C}^{M\tilde{N}}$ and indices $j = 1, 2, \dots, M$, there is a constant $C_j \equiv C_j(\Sigma) \in (0, \infty)$ such that the following inequalities,

$$\left| \frac{\partial f_j}{\partial t}(x, t; X) \right| \leq C_j(\Sigma) \quad \text{and} \quad \left| \frac{\partial f_j}{\partial X_{\beta,k}}(x, t; X) \right| \leq C_j(\Sigma) \tag{3.5}$$

for all $(x, t) \in \Omega$ and for all $X = (X_\beta)_{|\beta| \leq m} \in \Sigma$, hold for all $\beta = (\beta_1, \dots, \beta_N) \in (\mathbb{Z}_+)^N$ with $|\beta| \leq m$ and for all $k = 1, 2, \dots, M$.

A simple, but more restrictive alternative to formulate (H1)–(H3) is to replace $\Omega = \mathfrak{X}^{(r_0)} \times \Delta_{\vartheta_0}^{T_0, T}$ by a larger, but simpler product domain $\Omega_0 = \mathfrak{X}^{(r_0)} \times \mathfrak{T}_{0, T}^{(\tau_0)}$ (defined in (1.8) and (2.1)) with $\tau_0 = T_0 \cdot \tan \vartheta_0$; hence, $\Omega \subset \Omega_0$, thanks to $\Delta_{\vartheta_0}^{T_0, T} \subset \mathfrak{T}_{0, T}^{(\tau_0)}$.

Let us recall our abbreviation $X = (X_\beta)_{|\beta| \leq m} \in \mathbb{R}^{M\tilde{N}}$ (or $\mathbb{C}^{M\tilde{N}}$) with $\tilde{N} = \sum_{k=0}^m N^k$ from (1.2) and make it more precise as follows: When dealing with *complex* (partial) derivatives of the function f_j with respect to the variable $X_{\beta,k}$, we prefer to replace $X_{\beta,k}$ by the complex variable $Z_{\beta,k} = X_{\beta,k} + iY_{\beta,k} \in \mathbb{C}$ ($X_{\beta,k}, Y_{\beta,k} \in \mathbb{R}$); $j, k = 1, 2, \dots, M$, and write $\frac{\partial f_j}{\partial Z_{\beta,k}}(x, t; Z)$ in place of $\frac{\partial f_j}{\partial X_{\beta,k}}(x, t; X)$.

The *strong ellipticity inequality* (3.3) can be improved as follows; cf. Takáč [82, Remark 3.1, p. 57]:

Remark 3.2. In a smaller domain $\Omega' = \mathfrak{X}^{(r_0)} \times \Delta_{\vartheta'_0}^{T_0, T} \subset \Omega$, with some number $\vartheta'_0 \in (0, \vartheta_0]$ (small enough), inequality (3.3) holds in the following qualitatively stronger form, cf. Agmon [2], Theorem 7.12, ineq. (7.21) on p. 87:

$$\Re\left(e^{i\theta} \cdot \sum_{j,k=1}^M \sum_{|\alpha|=|\beta|=m} P_{jk}^{\alpha\beta}(z, t) \xi^{\alpha+\beta} \eta_k \bar{\eta}_j\right) \geq c' |\xi|^{2m} |\boldsymbol{\eta}|^2 \tag{3.3'}$$

for all $\theta \in [-\vartheta'_0, \vartheta'_0]$, for all $(z, t) \in \bar{\Omega}'$, and for all $\xi = (\xi_1, \dots, \xi_N) \in \mathbb{R}^N$ and $\boldsymbol{\eta} = (\eta_1, \dots, \eta_M) \in \mathbb{C}^M$, where $c' \in (0, c]$ is a constant. Recall that $0 < \vartheta'_0 \leq \vartheta_0$ and $\Omega' \subset \Omega$. Consequently, without loss of generality, we may remove the prime ($'$) from both ϑ'_0 and c' in (3.3') and assume that

$$\Re\left(e^{i\theta} \cdot \sum_{j,k=1}^M \sum_{|\alpha|=|\beta|=m} P_{jk}^{\alpha\beta}(z, t) \xi^{\alpha+\beta} \eta_k \bar{\eta}_j\right) \geq c |\xi|^{2m} |\boldsymbol{\eta}|^2 \tag{3.6}$$

for all $\theta \in [-\vartheta_0, \vartheta_0]$ and for all $(z, t) \in \bar{\Omega}$, $\xi = (\xi_1, \dots, \xi_N) \in \mathbb{R}^N$, and $\boldsymbol{\eta} = (\eta_1, \dots, \eta_M) \in \mathbb{C}^M$, where $c > 0$ is a constant. We prefer to use inequality (3.6) in place of (3.3).

The *Gårding inequality* (in the whole space \mathbb{R}^N) below is an important consequence of inequality (3.6); see e.g. Agmon [2, Theorem 7.6, p. 78]:

Corollary 3.3 (Gårding’s inequality). *Under Hypotheses (H1), (H2), there exist constants c_1 and c_2 , $c_1 > 0$ and $0 \leq c_2 < \infty$, such that*

$$\begin{aligned} & \Re\left[e^{i\theta} \cdot \sum_{|\alpha|=|\beta|=m} \int_{\mathbb{R}^N} D_x^{\bar{\alpha}} \mathbf{w} \cdot \mathbf{P}^{\alpha\beta}(x + iy, t) D_x^{\beta} \mathbf{w} \, dx\right] \\ & \geq c_1 \sum_{|\alpha|=m} \|D_x^{\alpha} \mathbf{w}\|_{L^2(\mathbb{R}^N)}^2 - c_2 \|\mathbf{w}\|_{L^2(\mathbb{R}^N)}^2 \end{aligned} \tag{3.7}$$

holds for all $\mathbf{w} \in W^{m,2}(\mathbb{R}^N)$ and for all $\theta \in [-\vartheta_0, \vartheta_0]$, $y \in \bar{Q}^{(r_0)}$, and $t \in \bar{\Delta}_{\vartheta_0}^{T_0, T}$.

Proof. The reader is referred to Agmon [2, Theorem 7.6, pp. 78–86] for a proof. We remark that the proof of Gårding’s inequality ([2, Lemma 7.9, p. 81]) requires the *uniform equicontinuity* of the leading coefficients $\mathbf{P}^{\alpha\beta}(x + iy, t)$ as functions of $x \in \mathbb{R}^N$ parametrized by $y \in \bar{Q}^{(r_0)}$ and $t \in \bar{\Delta}_{\vartheta_0}^{T_0, T}$, where $\mathbf{P}^{\alpha\beta}(x + iy, t) = (P_{jk}^{\alpha\beta})_{j,k=1}^M$ for $|\alpha| = |\beta| = m$. This is guaranteed by our Hypothesis (H1) that all $P_{jk}^{\alpha\beta}$ ($|\alpha| = |\beta| = m$) belong to $C_{\text{unif}}^0(\bar{\Omega})$ as functions of $(z, t) = (x + iy, t) \in \bar{\Omega}$. \square

To give a natural lower estimate on the domain of holomorphy (i.e., the domain of complex analyticity) of a weak solution \mathbf{u} to the Cauchy problem (1.1), we introduce a few more subsets of $\mathbb{C}^N \times \mathbb{C}$ (cf. [82, p. 58] or [83, p. 428]).

We use the subdomain $\Gamma_T^{(T')}(r', \vartheta') = \mathfrak{X}^{(r')} \times \Delta_{\vartheta'}^{T', T}$ of the (larger) domain

$$\Omega = \Gamma_T^{(T_0)}(r_0, \vartheta_0) := \mathfrak{X}^{(r_0)} \times \Delta_{\vartheta_0}^{T_0, T} \subset \mathbb{C}^N \times \mathbb{C} \tag{3.8}$$

defined above (H1)–(H3). The three constants $T' \in (0, T_0]$, $r' \in (0, r_0]$, and $\vartheta' \in (0, \vartheta_0]$ used below will be specified later (in Theorem 3.4).

We recall the tube $\mathfrak{X}^{(r_0)} = \mathbb{R}^N + iQ^{(r_0)}$ (called often a *strip*) in \mathbb{C}^N with base $Q^{(r_0)} \subset \mathbb{R}^N$ defined in (2.1) and recall from (1.6) also the definition of the set $\Delta_{\vartheta}^{(T')}$ in \mathbb{C} . In formula (1.7) we employ the time translation of $\Delta_{\vartheta}^{(T')}$ by r units, i.e., the set $r + \Delta_{\vartheta}^{(T')}$, to define the union $\Delta_{\vartheta}^{T',T}$ of such translations for $0 \leq r \leq T - T' < \infty$. It is easy to see that, for $0 < s \leq T < \infty$ and $0 < \vartheta_0 < \pi/2$, we have

$$\begin{aligned} &\Delta_{\vartheta_0}^{s,T} \\ &= \Delta_{\vartheta_0}^{(s)} \cup ([s, T] + i(-s \cdot \tan \vartheta_0, s \cdot \tan \vartheta_0)) \\ &= \{t = \sigma + i\tau \in \mathbb{C} : 0 < \sigma < T \text{ and } |\tau| < r \cdot \tan \vartheta_0 \text{ where } r = \min\{\sigma, s\}\}. \end{aligned} \tag{3.9}$$

Recall that the closure of $\Delta_{\vartheta_0}^{s,T}$ (and $\Delta_{\vartheta_0}^{(s)}$, respectively) in \mathbb{C} is denoted by $\bar{\Delta}_{\vartheta_0}^{s,T}$ (and $\bar{\Delta}_{\vartheta_0}^{(s)}$). Given any $r \in [0, T)$, we observe that the (real) time r section of $\Delta_{\vartheta_0}^{s,T}$ is given by

$$\{t \in \Delta_{\vartheta_0}^{s,T} : \Re t = r\} = r + i(-r' \cdot \tan \vartheta_0, r' \cdot \tan \vartheta_0) \subset \mathbb{C}$$

where $r' = \min\{r, s\}$. These sets in the complex plane \mathbb{C} have already been sketched in Figure 1 above and will be sketched more precisely in Figure 2 below.

The Cartesian product $\Gamma_T^{(s)}(r_0, \vartheta_0) = \mathfrak{X}^{(r_0)} \times \Delta_{\vartheta_0}^{s,T}$ defined in (3.8) is our most important complex analyticity domain in $\mathbb{C}^N \times \mathbb{C}$. Recall that $B^{s;p,p}(\mathbb{R}^N) = (L^p(\mathbb{R}^N), W^{2m,p}(\mathbb{R}^N))_{1-(1/p),p} \hookrightarrow L^p(\mathbb{R}^N)$ with $p > 2 + \frac{N}{m}$. Our main result reads as follows.

Theorem 3.4. *Let $m, M, N \geq 1$, $p > 2 + \frac{N}{m}$, $0 < T < \infty$, and assume that (H1)–(H3) are satisfied in the product domain $\Omega = \Gamma_T^{(T_0)}(r_0, \vartheta_0) = \mathfrak{X}^{(r_0)} \times \Delta_{\vartheta_0}^{T_0,T} \subset \mathbb{C}^N \times \mathbb{C}$ with some constants $0 < r_0 < \infty$, $0 < T_0 \leq T$, and $0 < \vartheta_0 < \pi/2$; cf. (1.7) and (2.1).*

- (i) *For $t_0 \in [0, T)$ and any initial value $\mathbf{u}_0 \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ at time $t = t_0$, there is a number $T_1 \in (t_0, T]$, depending on t_0 and \mathbf{u}_0 , such that the Cauchy problem (1.1) on the (local) time interval $[t_0, T_1] \subset [0, T]$ with the initial condition $\mathbf{u}(\cdot, t_0) = \mathbf{u}_0$ in $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$ possesses a unique weak solution $\mathbf{u} \in C([t_0, T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$.*
- (ii) *Furthermore, for any initial data $\mathbf{u}_0 \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ at time $t = 0$, any (global) weak solution $\mathbf{u} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ of the Cauchy problem (1.1), if it exists, is always unique and it possesses a unique (temporal) extension to the space-time domain $\mathbb{R}^N \times \bar{\Delta}_{\vartheta'}^{T',T}$, denoted again by \mathbf{u} , such that the $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$ -valued function $\mathbf{u} : \bar{\Delta}_{\vartheta'}^{T',T} \mapsto \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ is continuous in $\bar{\Delta}_{\vartheta'}^{T',T}$ and its restriction to $\Delta_{\vartheta'}^{T',T}$ is holomorphic, provided the numbers $T' \in (0, T_0]$ and $\vartheta' \in (0, \vartheta_0]$ are small enough. This temporal extension is a unique weak solution to the Cauchy problem (1.1) in $\mathbb{R}^N \times \bar{\Delta}_{\vartheta'}^{T',T}$.*
- (iii) *If, in addition to $\mathbf{u}_0 \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$, \mathbf{u}_0 possesses a (unique) holomorphic extension $\tilde{\mathbf{u}}_0 : \mathfrak{X}^{(\kappa_0)} \rightarrow \mathbb{C}^M$ from \mathbb{R}^N to the complex domain $\mathfrak{X}^{(\kappa_0)} \subset \mathbb{C}^N$, for some $\kappa_0 \in (0, r_0]$, such that the function*

$$\tilde{\mathbf{u}}_0(\cdot + iy) : x \mapsto \tilde{\mathbf{u}}_0(x + iy) : \mathbb{R}^N \rightarrow \mathbb{C}^M$$

belongs to $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$ for each $y \in Q^{(\kappa_0)}$ and

$$\mathfrak{N}^{(\kappa_0)}(\tilde{\mathbf{u}}_0) := \sup_{y \in Q^{(\kappa_0)}} \|\tilde{\mathbf{u}}_0(\cdot + iy)\|_{B^{s;p,p}(\mathbb{R}^N)} < \infty, \tag{3.10}$$

then also any (global) weak solution $\mathbf{u} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ of the Cauchy problem (1.1), if it exists, possesses a unique extension to the space-time domain $\mathfrak{X}^{(r')} \times \bar{\Delta}_{\vartheta'}^{T',T}$, denoted by $\tilde{\mathbf{u}}$, with the following properties, provided the numbers $T' \in (0, T_0]$, $r' \in (0, \kappa_0]$, and $\vartheta' \in (0, \vartheta_0]$ are small enough:

(iii.1) $\tilde{\mathbf{u}}(\cdot + iy, t) \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ holds for all $(y, t) \in Q^{(r')} \times \bar{\Delta}_{\vartheta'}^{T',T}$, together with

$$\sup_{t \in \bar{\Delta}_{\vartheta'}^{T',T}} \mathfrak{N}^{(r')}(\tilde{\mathbf{u}}(\cdot, t)) = \sup_{t \in \bar{\Delta}_{\vartheta'}^{T',T}} \sup_{y \in Q^{(\kappa_0)}} \|\tilde{\mathbf{u}}(\cdot + iy, t)\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)} < \infty, \quad (3.11)$$

(iii.2) the function

$$\tilde{\mathbf{u}} : (y, t) \mapsto \tilde{\mathbf{u}}(\cdot + iy, t) : Q^{(r')} \times \bar{\Delta}_{\vartheta'}^{T',T} \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$$

is continuous in the space-time variable $(y, t) \in Q^{(r')} \times \bar{\Delta}_{\vartheta'}^{T',T}$,

(iii.3) $\tilde{\mathbf{u}}$ is holomorphic in the complex domain

$$\Gamma_T^{(T')}(r', \vartheta') = \mathfrak{X}^{(r')} \times \Delta_{\vartheta'}^{T',T} \subset \mathfrak{X}^{(\kappa_0)} \times \Delta_{\vartheta_0}^{T_0,T} \subset \Omega = \Gamma_T^{(T_0)}(r_0, \vartheta_0) = \mathfrak{X}^{(r_0)} \times \Delta_{\vartheta_0}^{T_0,T}.$$

Finally, the extension $\tilde{\mathbf{u}}$ verifies the partial differential equation in the Cauchy problem (1.1) pointwise in $\Gamma_T^{(T')}(r', \vartheta')$, i.e., in the classical sense, and obeys the initial data as follows,

$$\|\tilde{\mathbf{u}}(\cdot + iy, t) - \tilde{\mathbf{u}}_0(\cdot + iy)\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)} \rightarrow 0 \quad \text{as } t \rightarrow 0, \quad t \in \Delta_{\vartheta'}^{(T')}, \quad (3.12)$$

for every $y \in Q^{(r')}$.

In Part (iii), properties (iii.1) and (iii.2) combined with the Sobolev(-Besov) imbedding $\mathbf{B}^{s;p,p}(\mathbb{R}^N) \hookrightarrow C^m(\mathbb{R}^N) \cap W^{m,\infty}(\mathbb{R}^N)$ guarantee the continuity and boundedness of the function $\tilde{\mathbf{u}} : \mathfrak{X}^{(r')} \times \bar{\Delta}_{\vartheta'}^{T',T} \rightarrow \mathbb{C}^M$.

Our condition $p > 2 + \frac{N}{m}$ is natural (and sharp) to guarantee the continuity of the Sobolev imbedding

$$\begin{aligned} \mathbf{B}^{s;p,p}(\mathbb{R}^N) &= [B^{s;p,p}(\mathbb{R}^N)]^M \\ &\hookrightarrow \mathbf{C}^m(\mathbb{R}^N) \cap \mathbf{W}^{m,\infty}(\mathbb{R}^N) = [C^m(\mathbb{R}^N) \cap W^{m,\infty}(\mathbb{R}^N)]^M \end{aligned}$$

which follows from the Sobolev inequalities and the Sobolev imbedding

$$B^{s-m;p,p}(\mathbb{R}^N) \hookrightarrow C_{\text{bdd}}^0(\mathbb{R}^N) := C(\mathbb{R}^N) \cap L^\infty(\mathbb{R}^N)$$

for $N/p < s - m < \infty$; see, e.g., Adams and Fournier [1, Chapt. 7], Theorem 7.34(c), p. 231. In our proposition below we abbreviate $\varsigma_1(s) = \min\{s, 1\}$ for $s \in \mathbb{R}_+$.

The above proposition follows directly from our proof of Part (iii) of Theorem 3.4. The temporal integration path in the double space-time integrals (3.13) and (3.14) is sketched in Figure 2.

Proposition 3.5. *In the situation of Theorem 3.4 above, there exist constants $c_0 > 0$ and $C_0 > 0$ depending solely on $\kappa_0, \vartheta_0, K, T', r', \vartheta'$ and the supremum norm*

$$\|\|\mathbf{u}\|\|_{L^\infty(0,T)} := \|\mathbf{u}\|_{C([0,T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))} = \sup_{0 \leq t \leq T} \|\mathbf{u}(\cdot, t)\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)} \quad (< \infty)$$

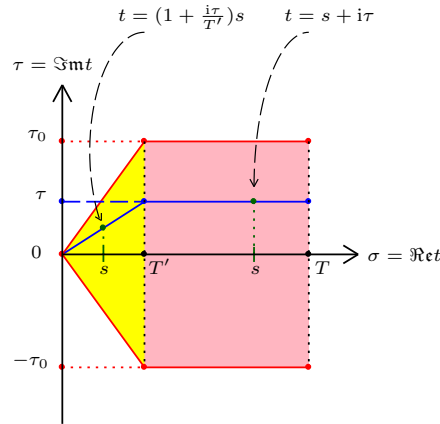


FIGURE 2. Here is a more detailed version of Figure 1 used in Proposition 3.5.

of the (global) weak solution $\mathbf{u} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$, such that the following estimate holds for all pairs $(y, t) \in Q^{(r')} \times \bar{\Delta}_{\vartheta'}^{T',T}$ with $t = \sigma + i\tau$ ($\sigma, \tau \in \mathbb{R}$):

$$\int_0^\sigma \int_{\mathbb{R}^N} |\partial_t \tilde{\mathbf{u}}(x + iy, s + i\varsigma_1(s/T')\tau)|^p dx ds + c_0 \sum_{|\alpha| \leq 2m} \int_0^\sigma \int_{\mathbb{R}^N} |D_x^\alpha \tilde{\mathbf{u}}(x + iy, s + i\varsigma_1(s/T')\tau)|^p dx ds \leq C_0. \tag{3.13}$$

Similarly, there are constants $c'_0 > 0$ and $C'_0 > 0$ depending solely on the constants c_0 and C_0 , such that the following estimate holds for all pairs $(y, t) \in Q^{(r')} \times \bar{\Delta}_{\vartheta'}^{T',T}$ with $t = \sigma + i\tau$ ($\sigma, \tau \in \mathbb{R}$):

$$\|\tilde{\mathbf{u}}(\cdot + iy, \sigma + i\varsigma_1(\sigma/T')\tau)\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)}^p + c'_0 \sum_{|\alpha| \leq 2m} \int_0^\sigma \int_{\mathbb{R}^N} |D_x^\alpha \tilde{\mathbf{u}}(x + iy, s + i\varsigma_1(s/T')\tau)|^p dx ds \leq C'_0. \tag{3.14}$$

Remark 3.6. (See Figure 2). Notice that, in (3.13) and (3.14), the temporal argument in the function $D_x^\alpha \mathbf{u}(x + iy, s + i\varsigma_1(s/T')\tau)$ reads $s + i\varsigma_1(s/T')\tau = (1 + i(\tau/T'))s$ whenever $0 \leq s < T' < \infty$, whereas $s + i\varsigma_1(s/T')\tau = s + i\tau$ holds whenever $0 < T' \leq s (\leq \sigma \leq T < \infty)$.

4. ABSTRACT CAUCHY PROBLEM IN AN INTERPOLATION SPACE

We assume that $E = (E_0, E_1)$ is a Banach couple, that is, E_0, E_1 are Banach spaces such that E_1 is densely and continuously imbedded into E_0 , i.e., $E_1 \hookrightarrow E_0$. We consider only complex Banach spaces over the field \mathbb{C} . Given a number $1 < p < \infty$, we denote by

$$E_{1-\frac{1}{p},p} \equiv (E_0, E_1)_{1-\frac{1}{p},p}$$

the real interpolation space between E_1 and E_0 obtained by the trace method as follows, with the parameter $\theta = 1 - \frac{1}{p} \in (0, 1)$. We define such an interpolation space for any $\theta \in (0, 1)$ below, cf. Lunardi [65, Chapt. 1], §1.2.2, pp. 20–25. The

reader is referred to Adams and Fournier [1, Chapt. 7], §7.6–§7.23, pp. 208–221, or Triebel [84, Chapt. 1], §1.8, pp. 41–55, for further details. The trace spaces were originally introduced in Lions [59, 60, 61].

Let X_θ^p denote the Banach space of all Bochner-measurable functions $u : \mathbb{R}_+ \rightarrow E_0$ endowed with the weighted Lebesgue norm

$$\|u\|_{X_\theta^p} := \left(\int_0^\infty \|t^{1-\theta} u(t)\|_{E_0}^p \frac{dt}{t} \right)^{1/p} \equiv \left(\int_0^\infty \|u(t)\|_{E_0}^p \frac{dt}{t^{1-(1-\theta)p}} \right)^{1/p} < \infty. \quad (4.1)$$

Notice that $X_{1-\frac{1}{p}}^p = L^p(\mathbb{R}_+ \rightarrow E_0)$. Analogously, we define the Banach space Y_θ^p of all functions $u \in X_\theta^p$ with the following properties: u can be identified (by equality a.e. in \mathbb{R}_+) with a Bochner-measurable function $u : \mathbb{R}_+ \rightarrow E_1$ satisfying

$$[u]_{Y_\theta^p} := \left(\int_0^\infty \|t^{1-\theta} u(t)\|_{E_1}^p \frac{dt}{t} \right)^{1/p} \equiv \left(\int_0^\infty \|u(t)\|_{E_1}^p \frac{dt}{t^{1-(1-\theta)p}} \right)^{1/p} < \infty, \quad (4.2)$$

and there is a function $v \in X_\theta^p$, denoted by $v = u'$ in the sequel, such that the equality

$$u(t_2) - u(t_1) = \int_{t_1}^{t_2} v(s) ds \quad \text{holds in } E_0 \text{ for all } 0 < t_1 \leq t_2 < \infty. \quad (4.3)$$

Applying Hölder's inequality to this equation, it is easy to show that every $u \in Y_\theta^p$ is θ -Hölder continuous on any compact interval $[0, T]$ as a function valued in E_0 , see, e.g., Lunardi [65, Chapt. 1], §1.2.2, p. 20. The Banach space Y_θ^p is endowed with the norm

$$\|u\|_{Y_\theta^p} := [u]_{Y_\theta^p} + \|u'\|_{X_\theta^p} < \infty. \quad (4.4)$$

For the special case $\theta = 1 - \frac{1}{p}$, a useful equivalent norm on $Y_{1-\frac{1}{p}}^p$ is defined by

$$\|u\|_{Y_{1-\frac{1}{p}}^p}^\# := \left(\int_0^\infty \|u(t)\|_{E_1}^p dt + \int_0^\infty \|u'(t)\|_{E_0}^p dt \right)^{1/p}. \quad (4.5)$$

Thus,

$$Y_{1-\frac{1}{p}}^p = L^p(\mathbb{R}_+ \rightarrow E_1) \cap W^{1,p}(\mathbb{R}_+ \rightarrow E_0)$$

is an abstract Sobolev space Amann [6, Chapt. III, §1.1], pp. 88–89).

Finally, the interpolation *trace space* is the vector space

$$E_{\theta,p} \equiv (E_0, E_1)_{\theta,p} := \{u(0) \in E_0 : u \in Y_\theta^p\}$$

of the initial values $x = u(0) \in E_0$ of all functions $u \in Y_\theta^p$ endowed with the trace norm

$$\|x\|_{E_{\theta,p}} := \inf \left\{ \|u\|_{Y_\theta^p} : x = u(0) \text{ for some } u \in Y_\theta^p \right\} \quad (4.6)$$

which makes the (linear) *trace mapping* $\tau : u \mapsto u(0) : Y_\theta^p \rightarrow E_{\theta,p}$ bounded (i.e., continuous), with the operator norm ≤ 1 . Equivalently to (4.6), we have $\|x\|_{E_{\theta,p}} \leq \|u\|_{Y_\theta^p}$ for every $u \in Y_\theta^p$ with $u(0) = x$ and there exists a sequence $\{u_n\}_{n=1}^\infty \subset Y_\theta^p$ such that $u_n(0) = x$ and $\|u_n\|_{Y_\theta^p} \rightarrow \|x\|_{E_{\theta,p}}$ as $n \rightarrow \infty$. It can be shown that there is a constant $c = c(\theta, p) > 0$, depending only on $E = (E_0, E_1)$, $\theta \in (0, 1)$, and $p \in (1, \infty)$, such that

$$\|x\|_{E_{\theta,p}} \leq c \|x\|_{E_0}^{1-\theta} \|x\|_{E_1}^\theta \quad \text{holds for all } x \in E_1; \quad (4.7)$$

see, e.g., Triebel [84, Chapt. 1], Theorem 1.3.3(g), p. 25, combined with Theorem 1.8.2, pp. 44–45. As an easy consequence of the definition of $E_{\theta,p}$ for $\theta = 1 - \frac{1}{p}$, i.e.,

$(1 - \theta)p = 1$, one can show that the abstract Sobolev space $Y_{1-\frac{1}{p}}^p$ is continuously imbedded into the Fréchet space $C(\mathbb{R}_+ \rightarrow E_{1-\frac{1}{p},p})$ of all continuous functions $u : \mathbb{R}_+ \rightarrow E_{1-\frac{1}{p},p}$ endowed with the (locally convex) topology of uniform convergence on every compact time interval $[t_1, t_2] \subset \mathbb{R}_+$,

$$Y_{1-\frac{1}{p}}^p = L^p(\mathbb{R}_+ \rightarrow E_1) \cap W^{1,p}(\mathbb{R}_+ \rightarrow E_0) \hookrightarrow C(\mathbb{R}_+ \rightarrow E_{1-\frac{1}{p},p}).$$

We complete our definition by setting $E_{\theta,p} := E_\theta$ if $\theta \in \{0, 1\}$.

In what follows we deal with applications of the interpolation trace space $E_{\theta,p}$ (with $\theta = 1 - \frac{1}{p}$) to abstract linear and nonlinear evolutionary problems of type

$$\begin{aligned} \frac{du}{dt} - A(t, u(t))u(t) &= f(t, u(t)) + g(t) \quad \text{for a.e. } t \in (0, T); \\ u(0) &= u_0 \in E_{1-\frac{1}{p},p}. \end{aligned} \quad (4.8)$$

Here, $u : (0, T) \rightarrow E_0$ is the unknown function valued in the Banach space E_0 and $0 < T \leq \infty$. A rigorous definition of a *strict solution* u of the initial value problem (4.8) will be given below, in Definition 4.4. Essentially, we follow Clément and Li [20], Section 1, pp. 17–18. A closely related approach is carried out also in Köhne, Prüss, and Wilke [55].

We denote by $\mathcal{L}(E_1 \rightarrow E_0)$ the Banach space of all bounded (i.e., continuous) linear operators $B : E_1 \rightarrow E_0$ endowed with the standard operator norm $\|B\|_{\mathcal{L}(E_1 \rightarrow E_0)}$. Let us denote by $I : E_1 \hookrightarrow E_0$ the continuous imbedding of E_1 into E_0 ; hence, $I \in \mathcal{L}(E_1 \rightarrow E_0)$. We identify I with the identity mapping in the whole of E_0 and abbreviate $\mathcal{L}(E_0) \equiv \mathcal{L}(E_0 \rightarrow E_0)$.

If, for some complex number $\lambda \in \mathbb{C}$, the operator $\lambda I - B \in \mathcal{L}(E_1 \rightarrow E_0)$ is invertible with an inverse denoted by $(\lambda I - B)^{-1} : E_0 \rightarrow E_1 \hookrightarrow E_0$ such that this inverse is bounded from E_0 into itself, i.e., $(\lambda I - B)^{-1} \in \mathcal{L}(E_0 \rightarrow E_0)$, then we alternatively (equivalently) view B as a densely defined, closed linear operator $B : E_0 \rightarrow E_0$ with the domain $\mathcal{D}(B) = E_1$, by the closed graph theorem, cf. Amann [6, Chapt. I, Lemma 1.1.2], p. 10. Indeed, if the graph $\mathcal{G}(B)$ of B is closed in $E_0 \times E_0$, it is closed also in $E_1 \times E_0$. In this case, the norm $\|\cdot\|_{E_1}$ on E_1 is equivalent with the *graph norm*

$$\|x\|_{\mathcal{D}(B)} := \|Bx\|_{E_0} + \|x\|_{E_0}, \quad x \in \mathcal{D}(B),$$

on $\mathcal{D}(B) = E_1$. An important class of such operators, denoted by $\text{Gen}(E) \equiv \text{Gen}(E_1 \rightarrow E_0)$, is formed by all closed linear operators $B : E_0 \rightarrow E_0$ with the domain $\mathcal{D}(B) = E_1$ that generate a strongly continuous semigroup $\{e^{tB} : t \geq 0\}$ on E_0 . We will consider only *generators* B with domain $\mathcal{D}(B) = E_1$. Finally, we denote by $\text{Hol}(E) \equiv \text{Hol}(E_1 \rightarrow E_0)$ the subset of all (infinitesimal) generators $B \in \text{Gen}(E)$ that generate a *holomorphic* (i.e., analytic) semigroup on E_0 . We refer to Amann [6, Chapt. I, §1], pp. 9–24, Pazy [72, Chapt. 1–2], pp. 1–75, or Tanabe [81, Chapt. 3, §3.1–§3.4], pp. 51–72, for details about strongly continuous (and holomorphic) semigroups.

Next, given an operator $B \in \text{Hol}(E)$, let us consider the following special (linear) case of problem (4.8), namely,

$$\begin{aligned} \frac{du}{dt} - Bu(t) &= g(t) \quad \text{for a.e. } t \in (0, T); \\ u(0) &= u_0 \in E_{1-\frac{1}{p},p}. \end{aligned} \quad (4.9)$$

Here, $u_0 \in E_{1-\frac{1}{p},p}$ is a given initial value, $g \in L^p((0, T) \rightarrow E_0)$ is a given function, $1 < p < \infty$, and $0 < T < \infty$. In analogy with our definition of the Banach spaces X_θ^p and Y_θ^p of functions $u : \mathbb{R}_+ \rightarrow E_0$ on the entire half line \mathbb{R}_+ , endowed with the norms given by eqs. (4.1) and (4.4), respectively, we introduce the corresponding Banach spaces $X_\theta^p(0, T)$ and $Y_\theta^p(0, T)$ of functions $u : [0, T] \rightarrow E_0$ on a bounded interval $[0, T]$, $0 < T < \infty$. Of course, in (4.1), the integral $\int_0^\infty \dots \frac{dt}{t}$ has to be replaced by $\int_0^T \dots \frac{dt}{t}$. It is not difficult to show that if one replaces the pair of spaces X_θ^p and Y_θ^p by $X_\theta^p(0, T)$ and $Y_\theta^p(0, T)$, respectively, in the definition of the trace space $E_{\theta,p}$ and its norm in (4.6), the *same* interpolation trace space is obtained. These facts can be inferred easily from the treatment of trace spaces in the monographs [1, 6, 65, 84] or from the original works by Lions [59, 60, 61]. In particular, we have the continuous imbedding

$$Y_{1-\frac{1}{p}}^p(0, T) = L^p((0, T) \rightarrow E_1) \cap W^{1,p}((0, T) \rightarrow E_0) \hookrightarrow C\left([0, T] \rightarrow E_{1-\frac{1}{p},p}\right), \quad (4.10)$$

see, e.g., [1, Chapt. 7], §7.67, p. 255. Thus, the (linear) trace mapping

$$\tau : u \mapsto u(0) : Y_{1-\frac{1}{p}}^p(0, T) \rightarrow E_{1-\frac{1}{p},p}$$

is continuous.

We say that a function $u : [0, T] \rightarrow E_0$ is a *strict solution* of the initial value problem (4.9) if

$$u \in Y_{1-\frac{1}{p}}^p(0, T), \quad \tau u \equiv u(0) = u_0,$$

and the differential equation in (4.9) is satisfied with all terms in $X_{1-\frac{1}{p}}^p(0, T) = L^p((0, T) \rightarrow E_0)$.

Definition 4.1. An infinitesimal generator $B \in \text{Hol}(E)$ of a holomorphic semi-group on E_0 with domain $\mathcal{D}(B) = E_1$ is said to possess the *maximal L^p -regularity property*, symbolically $B \in \text{MR}_p(E) \equiv \text{MR}_p(E_1 \rightarrow E_0)$, if for any given initial condition $u_0 \in E_{1-\frac{1}{p},p}$ and any given function $g \in L^p((0, T) \rightarrow E_0)$, problem (4.9) possesses a unique strict solution $u \in Y_{1-\frac{1}{p}}^p(0, T)$ that satisfies the following estimate:

There exists a constant $M \equiv M(p, E, B, T) > 0$, independent of u_0 and g , such that

$$\begin{aligned} & \int_0^T \left\| \frac{du}{dt} \right\| \left\| \frac{du}{dt} \right\|_{E_0}^p dt + \int_0^T \|Bu(t)\|_{E_0}^p dt \\ & \leq M \left(\|u_0\|_{E_{1-\frac{1}{p},p}}^p + \int_0^T \|g(t)\|_{E_0}^p dt \right). \end{aligned} \quad (4.11)$$

We have adopted this definition of class $\text{MR}_p(E)$ from Clément and Li [20, p. 18] and from the monograph by Ashyralyev and Sobolevskii [9, Chapt. 1], §3.5, p. 28. It may be viewed as some kind of *ellipticity hypothesis* for the linear operator $B \in \mathcal{L}(E_1 \rightarrow E_0)$ or a *stability hypothesis* for the linear parabolic problem (4.9). Equivalently, the abstract (linear evolutionary) partial differential operator

$$(\partial_t - B, \tau) : Y_{1-\frac{1}{p}}^p(0, T) \rightarrow L^p((0, T) \rightarrow E_0) \times E_{1-\frac{1}{p},p},$$

defined for every $u \in Y_{1-\frac{1}{p}}^p(0, T) = L^p((0, T) \rightarrow E_1) \cap W^{1,p}((0, T) \rightarrow E_0)$ by

$$(\partial_t - B, \tau) : u \mapsto \left(\frac{du}{dt} - Bu(t), u(0) \right), \quad (4.12)$$

possesses a bounded inverse furnished by the strict solution

$$u = (\partial_t - B, \tau)^{-1}(g, u_0) \in Y_{1-\frac{1}{p}}^p(0, T)$$

of problem (4.9); cf. Angenent [7, Lemma 2.2, p. 95] for the parallel interpolation case $p = \infty$ introduced in Da Prato and Grisvard [33].

Remark 4.2. (a) It is not difficult to show that the maximal L^p -regularity class $\text{MR}_p(E)$ is independent from a particular choice of $T \in (0, \infty)$; see Dore [24, Sect. 5, p. 310], Corollary 5.4, or Prüss [74, p. 4], remarks after Corollary 1.3. More importantly, this class is *independent* from $p \in (1, \infty)$ as well, i.e., $\text{MR}_p(E) = \text{MR}_{p_0}(E)$ holds for all $p, p_0 \in (1, \infty)$, by a classical result due to Sobolevskii [78, 79]; see, e.g., [78, §3.1, pp. 343–345]. Further details on the independence of $\text{MR}_p(E)$ from $p \in (1, \infty)$ can be found in Ashyralyev and Sobolevskii [9, Chapt. 1], §3.5, Theorem 3.6 on p. 35, Dore [24, Sect. 7, p. 313], Theorem 7.1, or Hieber [42, Corollary 4.4, p. 371], where in (4.11) one may take $M \equiv M(p) = p^2(p-1)^{-1}M(p_0) < \infty$ if the constant $M(p_0) \in (0, \infty)$ is known, by [9].

(b) We are allowed to specify the constant $M \equiv M(p, E, B, T) > 0$ in (4.11) to be the *smallest* nonnegative number $M \in \mathbb{R}_+$ for which (4.11) is valid; cf. Clément and Li [20, Proposition 2.2, p. 19]. Then, clearly, $T \mapsto M(p, E, B, T)$ is a *nondecreasing* (nonnegative) function of time $T \in (0, \infty)$. Indeed, if $T' \in (0, T)$ and $g \in L^p((0, T') \rightarrow E_0)$ is arbitrary, it suffices to apply (4.11) with the function

$$\tilde{g}(t) = \begin{cases} g(t) & \text{if } 0 \leq t \leq T'; \\ 0 & \text{if } T' < t \leq T, \end{cases}$$

in place of g in order to derive (4.11) for T' in place of T with the same constant M . Hence, $M(p, E, B, T') \leq M(p, E, B, T)$ holds for $0 < T' < T$. It is easy to see that $M(p, E, B, T) > 0$. (The case $M = 0$ would easily lead to a contradiction.) In what follows we always use this optimal value of M , i.e., $M \equiv M(p, E, B, T) > 0$.

(c) Simple perturbation theory for linear operators shows that the set $\text{Hol}(E) \equiv \text{Hol}(E_1 \rightarrow E_0)$ is open in the Banach space $\mathcal{L}(E_1 \rightarrow E_0)$. Even a more precise, relative perturbation result is valid; see Kato [52, Chapt. IX], §2.2, Theorem 2.4 on p. 499. A similar result can be derived for the class $\text{MR}_p(E) \equiv \text{MR}_p(E_1 \rightarrow E_0)$ applying the perturbation technique from either Amann [6, Chapt. III, §1.6], Proposition 1.6.3 on p. 97, or from Clément and Li [20, Proof of Theorem 2.1], pp. 19–23: The set $\text{MR}_p(E)$ is open in $\mathcal{L}(E_1 \rightarrow E_0)$; see Lemma 5.1 below. Indeed, this follows from the fact that the set of all bounded linear operators from

$$\mathcal{L}\left(Y_{1-\frac{1}{p}}^p(0, T) \rightarrow L^p((0, T) \rightarrow E_0) \times E_{1-\frac{1}{p}, p}\right)$$

that possess a bounded inverse is open in this Banach space, and the inverse $(\partial_t - B, \tau)^{-1}$ is a locally Lipschitz-continuous function of $B \in \mathcal{L}(E_1 \rightarrow E_0)$, by Lemma 5.1 below and formula (5.4) thereafter, with $B \in \text{MR}_p(E)$ being fixed and $A \in \mathcal{L}(E_1 \rightarrow E_0)$ having a sufficiently small operator norm $\|A\|_{\mathcal{L}(E_1 \rightarrow E_0)}$ depending on B .

Now we are ready to define a strict solution u to our abstract nonlinear evolutionary problem (4.8). We assume that $1 < p < \infty$, $0 < T < \infty$, $g \in L^p((0, T) \rightarrow E_0)$, $u_0 \in U$ where U is an open set in $E_{1-\frac{1}{p}, p}$, and the mappings

$$\begin{aligned} A &: (t, v) \mapsto A(t, v) : [0, T] \times U \subset [0, T] \times E_{1-\frac{1}{p}, p} \rightarrow \mathcal{L}(E_1 \rightarrow E_0), \\ f &: (t, v) \mapsto f(t, v) : [0, T] \times U \subset [0, T] \times E_{1-\frac{1}{p}, p} \rightarrow E_0 \end{aligned}$$

satisfy the following “natural” hypotheses (cf. Clément and Li [20, p. 19], (H1)–(H3)):

4.1. Hypothesis.

(H4) $A : [0, T] \times U \rightarrow \mathcal{L}(E_1 \rightarrow E_0)$ is a Lipschitz continuous mapping such that $A(t, v) \in \text{MR}_p(E)$ for all $(t, v) \in [0, T] \times U$.

(H5) $f : [0, T] \times U \rightarrow E_0$ is a Lipschitz continuous mapping.

Of course, the metric on $[0, T] \times U$ is induced by the canonical norm on $\mathbb{R} \times E_{1-\frac{1}{p}, p}$. It is a matter of a straight forward calculation to verify that both substitution mappings,

$$\begin{aligned} (v, u) &\mapsto [t \mapsto A(t, v(t))u(t)] : C([0, T] \rightarrow U) \times L^p((0, T) \rightarrow E_1) \rightarrow L^p((0, T) \rightarrow E_0), \\ v &\mapsto [t \mapsto f(t, v(t))] : C([0, T] \rightarrow U) \rightarrow L^p((0, T) \rightarrow E_0), \end{aligned}$$

are locally Lipschitz continuous with values in $L^p((0, T) \rightarrow E_0)$; see, e.g., Clément and Li [20, Proof of Theorem 2.1], pp. 19–23.

Remark 4.3. In (H4) we did not have to assume that $A(t, v) \in \text{MR}_p(E)$ holds for all $(t, v) \in [0, T] \times U$. We could assume *only* $A(0, u_0) \in \text{MR}_p(E)$; cf. results to follow below (e.g., Theorems 4.5 and 4.7 and Remark 4.6). However, the set $\text{MR}_p(E)$ being open in $\mathcal{L}(E_1 \rightarrow E_0)$, $A(0, u_0) \in \text{MR}_p(E)$ would imply that there are a number $T_0 \in (0, T]$ and an open neighborhood U_0 of u_0 in $E_{1-\frac{1}{p}, p}$, $u_0 \in U_0 \subset U$, such that $A(t, v) \in \text{MR}_p(E)$ holds for all $(t, v) \in [0, T_0] \times U_0$, by the Lipschitz continuity of A . But this statement is qualitatively the same as $A(t, v) \in \text{MR}_p(E)$ for all $(t, v) \in [0, T] \times U$ in our (H4).

Definition 4.4 (Clément and Li [20, p. 18]). Recall that U is an open set in $E_{1-\frac{1}{p}, p}$ and $u_0 \in U$. We say that a function $u : [0, T] \rightarrow E_0$ is a *strict solution* of the initial value problem (4.8) if $u \in Y_{1-\frac{1}{p}}^p(0, T)$, $u(t) \in U$ for every $t \in [0, T]$, $u(0) = u_0$, and the differential equation in (4.8) is satisfied with all terms (summands) in $L^p((0, T) \rightarrow E_0)$.

We recall that the Banach space $Y_{1-\frac{1}{p}}^p(0, T)$ has been introduced in (4.10).

The main result in [20, Theorem 2.1, p. 19] is *local in time* and reads as follows, with (H4) being somewhat weakened in the sense of our Remark 4.3 above.

Theorem 4.5. *Let $1 < p < \infty$ and $0 < T < \infty$. Let U be a nonempty open set in $E_{1-\frac{1}{p}, p}$ and $u_0 \in U$. Assume that both mappings $A : [0, T] \times U \rightarrow \mathcal{L}(E_1 \rightarrow E_0)$ and $f : [0, T] \times U \rightarrow E_0$ are Lipschitz continuous. If $A(0, u_0) \in \text{MR}_p(E)$ then there exists some time $T_1 \equiv T_1(u_0) \in (0, T]$, depending on u_0 , such that the abstract initial value problem (4.8) possesses a unique strict solution*

$$\begin{aligned} u &\in Y_{1-\frac{1}{p}}^p(0, T_1) \\ &\left(= L^p((0, T_1) \rightarrow E_1) \cap W^{1,p}((0, T_1) \rightarrow E_0) \hookrightarrow C([0, T_1] \rightarrow E_{1-\frac{1}{p}, p}) \right) \end{aligned} \quad (4.13)$$

on the time interval $[0, T_1]$. Consequently, one has $u(t) \in U$ for every $t \in [0, T_1]$.

This theorem is proved in [20], Section 2, pp. 20–23, using the Banach contraction principle in the closed ball

$$\Sigma_{\rho_1, T_1}^{(u_0)} = \{v \in Y^{T_1} : v(0) = u_0 \quad \text{and} \quad \|v - w\|_{Y^{T_1}} \leq \rho_1\}$$

of radius $\rho_1 \in (0, \infty)$ centered at the point $w \in Y^{T_1}$ in the Banach space

$$Y^{T_1} = Y_{1-\frac{1}{p}}^p(0, T_1) = L^p((0, T_1) \rightarrow E_1) \cap W^{1,p}((0, T_1) \rightarrow E_0).$$

Here, the “center” function $w \in Y^{T_1}$ is defined to be the restriction to $[0, T_1]$ of the unique strict solution $\tilde{w} \in Y^T = Y_{1-\frac{1}{p}}^p(0, T)$ to the abstract initial value problem (4.9) in the time interval $[0, T]$ with the linear operator $B = A(0, u_0) \in \text{MR}_p(E)$ and the right-hand side $g(t)$ replaced by the sum $f(t, u_0) + g(t)$,

$$\begin{aligned} \frac{d\tilde{w}}{dt} - A(0, u_0)\tilde{w}(t) &= f(t, u_0) + g(t) \quad \text{for a.e. } t \in (0, T); \\ \tilde{w}(0) &= u_0 \in E_{1-\frac{1}{p}, p}. \end{aligned} \quad (4.14)$$

Although the proof in [20] has been carried out only for $A(t, u) = A(u)$ independent from time $t \in [0, T]$, it is a matter of straight forward calculations to adapt this proof to the case of $A(t, u)$ depending on time t , cf. [20, p. 23], Remark at the end of Section 2. A detailed treatment of the latter case is presented in Prüss [74, pp. 9–13], Chapt. 3, under slightly different assumptions (see also Köhne, Prüss, and Wilke [55]).

Remark 4.6. Furthermore, one can easily conclude from the proof of Theorem 2.1 in [20, pp. 20–23] that if $\bar{B}_{R_0}(w_0)$ is any closed ball in the Banach space $E_{1-\frac{1}{p}, p}$ of radius $R_0 \in (0, \infty)$ centered at a point $w_0 \in E_{1-\frac{1}{p}, p}$, such that $\bar{B}_{R_0}(w_0) \subset U$ and $R_0 > 0$ is small enough, then the constants $\rho_1 \in (0, \infty)$ and $T_1 \in (0, T]$ can be chosen small enough to depend solely on R_0 , but not on w_0 , provided $u_0 \in \bar{B}_{R_0}(w_0) \subset U$. The estimates in [20, pp. 20–23], based on the Lipschitz constants for A and f in $[0, T] \times U$ and the estimate in (4.11), remain valid for any $u_0 \in \bar{B}_{R_0}(w_0)$. Thus, we have $T_1 \equiv T_1(R_0) \in (0, T]$ and $\rho_1 \equiv \rho_1(R_0) \in (0, \infty)$. Finally, using similar estimates, cf. [20, p. 22], (2.14)–(2.17), one can show that the (strict) solution mapping

$$u_0 \mapsto u : \bar{B}_{R_0}(w_0) \subset U \subset E_{1-\frac{1}{p}, p} \rightarrow Y^{T_1} = Y_{1-\frac{1}{p}}^p(0, T_1)$$

is Lipschitz continuous with a Lipschitz constant $L \equiv L(R_0) \in (0, \infty)$ independent from $w_0 \in E_{1-\frac{1}{p}, p}$, such that $\bar{B}_{R_0}(w_0) \subset U$ and $R_0 > 0$ is small enough. This means that if $u_1, u_2 : [0, T_1] \rightarrow E_{1-\frac{1}{p}, p}$ are two strict solutions to problem (4.8) on the time interval $[0, T_1]$, with (possibly different) initial values $u_1(0) = u_{0,1}$ and $u_2(0) = u_{0,2}$ in $\bar{B}_{R_0}(w_0) \subset U$, then one has $u_1(t), u_2(t) \in U$ for all $t \in [0, T_1]$ and

$$\|u_1 - u_2\|_{Y^{T_1}} \leq L \|u_{0,1} - u_{0,2}\|_{E_{1-\frac{1}{p}, p}}. \quad (4.15)$$

Combining this with the continuous imbedding $Y^{T_1} \hookrightarrow C([0, T_1] \rightarrow E_{1-\frac{1}{p}, p})$ in (4.10), we obtain

$$\|u_1(t) - u_2(t)\|_{E_{1-\frac{1}{p}, p}} \leq L_1 \|u_{0,1} - u_{0,2}\|_{E_{1-\frac{1}{p}, p}} \quad \text{for all } t \in [0, T_1], \quad (4.16)$$

with another Lipschitz constant $L_1 \equiv L_1(R_0) \in (0, \infty)$.

A number of sufficient conditions that guarantee the existence of a *global* weak solution $\mathbf{u} : \mathbb{R}^N \times (0, T) \rightarrow \mathbb{R}^M$ (\mathbb{C}^M) for all times $t \in (0, T)$ to the parabolic Cauchy problem (1.1) can be found in Amann [4, 5] for systems similar to ours. As we do not wish to impose those kinds of restrictive growth conditions on the reaction function \mathbf{f} on the right-hand side of (1.1), we prefer to assume the *existence* of a fixed global strict solution (cf. (4.13))

$$\begin{aligned}
 w &\in Y_{1-\frac{1}{p}}^p(0, T) \\
 &\left(= L^p((0, T) \rightarrow E_1) \cap W^{1,p}((0, T) \rightarrow E_0) \hookrightarrow C([0, T] \rightarrow E_{1-\frac{1}{p},p}) \right)
 \end{aligned}
 \tag{4.17}$$

to problem (4.8) on the whole time interval $[0, T]$, for some $T \in (0, \infty)$, with a prescribed initial value $w(0) = w_0 \in U \subset E_{1-\frac{1}{p},p}$ and such that $w(t) \in U$ and $A(t, w(t)) \in \text{MR}_p(E)$ for all $t \in [0, T]$. Then the local Theorem 4.5 and Remark 4.6 from above may be applied on any time interval $[t_0, t_0 + T_1] \subset [0, T]$ of sufficiently short length $T_1 > 0$ in order to obtain unique strict solutions u “along” the known solution w to the following abstract initial value problem:

$$\begin{aligned}
 \frac{du}{dt} - A(t, u(t))u(t) &= f(t, u(t)) + g(t) \quad \text{for a.e. } t \in (t_0, t_0 + T_1); \\
 u(t_0) &= u_0 \in E_{1-\frac{1}{p},p}.
 \end{aligned}
 \tag{4.18}$$

Here, $u_0 \in \bar{B}_{R_0}(w(t_0))$ is arbitrary, where the radius $R_0 > 0$ is small enough, as described in Remark 4.6, such that $\bar{B}_{R_0}(w(t_0)) \subset U$. By Theorem 4.5, the strict solution $u : [t_0, t_0 + T_1] \rightarrow E_{1-\frac{1}{p},p}$ satisfies $u(t) \in U$ for every $t \in [t_0, t_0 + T_1]$. The image $w([0, T]) = \{w(t) : t \in [0, T]\}$ of the solution w being compact in the open set $U \subset E_{1-\frac{1}{p},p}$, we may choose $R_0 > 0$ even smaller, such that $\bar{B}_{R_0}(w(t)) \subset U$ holds for all $t \in [0, T]$.

In addition to these claims that follow immediately from the proof of [20, Theorem 2.1, pp. 20–23], one can deduce from inequalities analogous to those in [20, p. 22], (2.14)–(2.17), cf. Remark 4.6 above, (4.16), that there exists a Lipschitz constant $L_1 \in [1, \infty)$, such that if $u_1, u_2 : [t_0, t_0 + T_1] \rightarrow E_{1-\frac{1}{p},p}$ are two strict solutions to problem (4.18) with initial values $u_1(t_0) = u_{0,1}$ and $u_2(t_0) = u_{0,2}$ in $\bar{B}_{R_0}(w(t_0)) \subset U$, then one has $u_1(t), u_2(t) \in U$ and

$$\|u_1(t) - u_2(t)\|_{E_{1-\frac{1}{p},p}} \leq L_1 \|u_1(t_0) - u_2(t_0)\|_{E_{1-\frac{1}{p},p}}
 \tag{4.19}$$

for all $t \in [t_0, t_0 + T_1]$. Consequently, fixing the smallest integer $m \in \mathbb{N}$ such that $m \geq T/T_1$ (≥ 1), we obtain, by “induction” on $k = 1, 2, 3, \dots, m$, first

$$\|u(t) - w(t)\|_{E_{1-\frac{1}{p},p}} \leq L_1^k \|u_0 - w_0\|_{E_{1-\frac{1}{p},p}} \leq L_1^k \cdot R_0/L_1^k = R_0
 \tag{4.20}$$

for all $t \in [0, \min\{kT_1, T\}]$, whenever $\|u_0 - w_0\|_{E_{1-\frac{1}{p},p}} \leq R_0/L_1^k$; also $u_1(t), u_2(t) \in \bar{B}_{R_0}(w(t)) \subset U$ and

$$\|u_1(t) - u_2(t)\|_{E_{1-\frac{1}{p},p}} \leq L_1^k \|u_{0,1} - u_{0,2}\|_{E_{1-\frac{1}{p},p}}
 \tag{4.21}$$

for all $t \in [0, \min\{kT_1, T\}]$, whenever

$$\|u_{0,j} - w_0\|_{E_{1-\frac{1}{p},p}} \leq R_k := R_0/L_1^k \quad (> 0); \quad j = 1, 2,$$

for $k = 1, 2, 3, \dots, m$.

We have thus obtained the following result, *global in time* on an arbitrary time interval (t_0, T) , $0 \leq t_0 < T$, with the constants $R_0 \in (0, \infty)$, $T_1 \equiv T_1(R_0) \in (0, T]$, and $L_1 \in [1, \infty)$ specified above in (4.19)–(4.21):

Theorem 4.7. *Let $1 < p < \infty$, $0 < T < \infty$, and $g \in L^p((0, T) \rightarrow E_0)$. Assume that U is a nonempty open subset of $E_{1-\frac{1}{p}, p}$ and A and f satisfy (H4) and (H5), respectively. Finally, assume that $w : [0, T] \rightarrow U \subset E_{1-\frac{1}{p}, p}$ is a fixed global strict solution to problem (4.8) satisfying (4.17), with a prescribed initial value $w(0) = w_0 \in U$ and such that $w(t) \in U$ and $A(t, w(t)) \in \text{MR}_p(E)$ for all $t \in [0, T]$. Then there exist some constant $R_0 \in (0, \infty)$, sufficiently small, with the following two properties, where $R_m = R_0/L_1^m \in (0, R_0]$ is the constant defined in (4.21):*

- (i) *If $t_0 \in [0, T)$ and $u_0 \in \bar{B}_{R_m}(w(t_0)) \subset U$, then the abstract initial value problem (4.8) on the time-interval (t_0, T) with $u(t_0) = u_0$ possesses a unique strict solution $u \in Y_{1-\frac{1}{p}}^p(t_0, T)$ (cf. (4.13))*

$$u \in Y_{1-\frac{1}{p}}^p(t_0, T) \left(= L^p((t_0, T) \rightarrow E_1) \cap W^{1,p}((t_0, T) \rightarrow E_0) \hookrightarrow C([t_0, T] \rightarrow E_{1-\frac{1}{p}, p}) \right)$$

such that $u(t) \in \bar{B}_{R_0}(w(t)) \subset U$ for every $t \in [t_0, T]$.

- (ii) *If $t_0 \in [0, T)$ and $u_1, u_2 : [t_0, T] \rightarrow E_{1-\frac{1}{p}, p}$ are two strict solutions to problem (4.8) on the time-interval (t_0, T) with initial values $u_1(t_0) = u_{0,1}$ and $u_2(t_0) = u_{0,2}$ in $\bar{B}_{R_m}(w(t_0)) \subset U$, then one has $u_1(t), u_2(t) \in \bar{B}_{R_0}(w(t)) \subset U$ and*

$$\|u_1(t) - u_2(t)\|_{E_{1-\frac{1}{p}, p}} \leq L_1^m \|u_{0,1} - u_{0,2}\|_{E_{1-\frac{1}{p}, p}} \quad \text{for all } t \in [t_0, T]. \tag{4.22}$$

5. ANALYTICITY IN TIME FOR THE ABSTRACT CAUCHY PROBLEM

In this section we establish a few temporal analyticity results, Theorem 5.3 being the most important among them, that will be used later (in Section 8) in order to prove Part (ii) of Theorem 3.4.

5.1. Auxiliary linear perturbation results. We begin by quoting a well-known result: If $B \in \text{Gen}(E)$ is the generator of a holomorphic semigroup on E_0 with the domain $\mathcal{D}(B) = E_1$, i.e., $B \in \text{Hol}(E)$, then so is every operator $B_\nu = (1 + i\nu)B : E_1 \subset E_0 \rightarrow E_0$, $\nu \in \mathbb{R}$, provided $|\nu|$ is small enough, $|\nu| \leq \delta_1 < 1$; see, e.g., Amann [6, Chapt. I, §1], pp. 9–24, Pazy [72, §3.2, pp. 80–81], or Tanabe [81, Chapt. 3, §3.1–§3.4], pp. 51–72. A more general perturbation theorem for generators of holomorphic semigroups is proved in Pazy [72, §3.2], Theorem 2.1 on p. 80. An analogous perturbation result for the smaller class $\text{MR}_p(E)$ ($\text{MR}_p(E) \subset \text{Hol}(E) \subset \text{Gen}(E)$) is proved in Amann [6, Chapt. III, §1.6], Proposition 1.6.3 on p. 97. Since we take advantage of the latter in an essential manner, we now give its precise formulation.

Let $1 < p < \infty$ and $0 < T < \infty$. Given any generator $B \in \text{Gen}(E)$, let us consider the bounded linear operator $\tilde{K}_B : L^1((0, T) \rightarrow E_0) \rightarrow L^\infty((0, T) \rightarrow E_0)$ defined by

$$(\tilde{K}_B g)(t) := \int_0^t e^{(t-s)B} g(s) ds \in E_0 \tag{5.1}$$

for all $t \in [0, T]$ and all $g \in L^1((0, T) \rightarrow E_0)$. It is proved in [6, Chapt. III, §1.5], Theorem 1.5.2 on p. 95, that if $B \in \text{Hol}(E)$ and B possesses the maximal

L^p -regularity property, i.e., $B \in \text{MR}_p(E) \equiv \text{MR}_p(E_1 \rightarrow E_0)$, then the restriction

$$K_B = \tilde{K}_B|_{X^T} \text{ of } \tilde{K}_B \text{ to } X^T := X_{1-\frac{1}{p}}^p(0, T) = L^p((0, T) \rightarrow E_0)$$

is a bounded linear operator from the Banach space X^T into another Banach space

$$Y^T := Y_{1-\frac{1}{p}}^p(0, T) = L^p((0, T) \rightarrow E_1) \cap W^{1,p}((0, T) \rightarrow E_0)$$

with the operator norm $\|K_B\|_{\mathcal{L}(X^T \rightarrow Y^T)} < \infty$. For the perturbed initial value problem

$$\begin{aligned} \frac{du}{dt} - (B + A)u(t) &= g(t) \quad \text{for a.e. } t \in (0, T); \\ u(0) &= u_0 \in E_{1-\frac{1}{p}, p}, \end{aligned} \quad (5.2)$$

the following result is established in [6, Chapt. III, §1.6], Proposition 1.6.3 on p. 97:

Lemma 5.1. *Assume that $B \in \text{MR}_p(E)$ and let $A \in \mathcal{L}(E_1 \rightarrow E_0)$ be arbitrary with the norm*

$$\|A\|_{\mathcal{L}(E_1 \rightarrow E_0)} \leq \gamma / \|K_B\|_{\mathcal{L}(X^T \rightarrow Y^T)} \quad \text{for some } \gamma \in (0, 1).$$

Then also the operator $B_A = B + A \in \mathcal{L}(E_1 \rightarrow E_0)$ belongs to the class $\text{MR}_p(E)$ and the operator norms of the inverses of the abstract (linear) partial differential operators

$$(\partial_t - B, \tau), (\partial_t - B - A, \tau) : Y_{1-\frac{1}{p}}^p(0, T) \rightarrow L^p((0, T) \rightarrow E_0) \times E_{1-\frac{1}{p}, p}$$

defined in (4.12) satisfy

$$\|(\partial_t - B - A, \tau)^{-1}\| \leq C \cdot (1 - \gamma)^{-1} \|(\partial_t - B, \tau)^{-1}\|, \quad (5.3)$$

where $C \equiv C(p, E, T) > 0$ is a constant independent of A, B , and γ .

More precisely, we have

$$(\partial_t - B - A, \tau)^{-1} = (I - K_B A)^{-1} (\partial_t - B, \tau)^{-1} \quad (5.4)$$

with the operator norm of the product

$$K_B A : Y^T \rightarrow Y^T = Y_{1-\frac{1}{p}}^p(0, T) = L^p((0, T) \rightarrow E_1)$$

bounded above by

$$\|K_B A\|_{\mathcal{L}(Y^T \rightarrow Y^T)} \leq \|K_B\|_{\mathcal{L}(X^T \rightarrow Y^T)} \cdot \|A\|_{\mathcal{L}(E_1 \rightarrow E_0)} \leq \gamma < 1.$$

Here, I stands for the identity mapping in $\mathcal{L}(Y^T \rightarrow Y^T)$. Hence, the Neumann series $(I - K_B A)^{-1} = \sum_{k=0}^{\infty} (K_B A)^k$ converges absolutely in $\mathcal{L}(Y^T \rightarrow Y^T)$ and $\|(I - K_B A)^{-1}\|_{\mathcal{L}(Y^T \rightarrow Y^T)} \leq (1 - \gamma)^{-1} < \infty$.

The following claims are trivial applications of this lemma: $\text{MR}_p(E) \equiv \text{MR}_p(E_1 \rightarrow E_0)$ is an open subset of the Banach space $\mathcal{L}(E_1 \rightarrow E_0)$. Furthermore, if $B \in \text{MR}_p(E)$ and $A \in \mathcal{L}(E_1 \rightarrow E_0)$ then also $B_{\nu A} = B + \nu A \in \text{MR}_p(E)$ holds for every $\nu \in \mathbb{C}$ provided $|\nu|$ is small enough,

$$|\nu| \leq \delta_1 = \gamma \|A\|_{\mathcal{L}(E_1 \rightarrow E_0)}^{-1} \|K_B\|_{\mathcal{L}(X^T \rightarrow Y^T)}^{-1} < \infty.$$

Naturally, the special case $A = iB$ is of interest.

The following perturbation lemma for problem (5.2) is related to Angenent [7, Lemma 2.5, p. 97]; see also Denk, Hieber, and Prüss [23], Proposition 4.3 on p. 44 and Theorem 4.4 on p. 45.

Lemma 5.2. *Assume that $B \in \text{MR}_p(E)$. Then there exists a number $\delta \in (0, 1)$ and a constant $C_\delta \in \mathbb{R}_+$ with the following property: If $A \in \mathcal{L}(E_1 \rightarrow E_0)$ is arbitrary with the norm*

$$\|Au\|_{E_0} \leq \delta \|Bu\|_{E_0} + C_\delta \|u\|_{E_0} \quad \text{for all } u \in E_1, \tag{5.5}$$

then also the operator $B_A = B + A \in \mathcal{L}(E_1 \rightarrow E_0)$ belongs to the class $\text{MR}_p(E)$. Furthermore, there exists a constant $\tilde{M} \equiv \tilde{M}(p, E, B, \delta, C_\delta, T) > 0$, independent of $(g, u_0) \in L^p((0, T) \rightarrow E_0) \times E_{1-\frac{1}{p}, p}$, such that the unique strict solution $v = (\partial_t - B - A, \tau)^{-1}(g, u_0)$ to the perturbed initial value problem (5.2) satisfies the inequality

$$\begin{aligned} & \int_0^T \left\| \frac{dv}{dt} \right\|_{E_0}^p dt + \int_0^T \|(B + A)v(t)\|_{E_0}^p dt \\ & \leq \tilde{M} \left(\|u_0\|_{E_{1-\frac{1}{p}, p}}^p + \int_0^T \|g(t)\|_{E_0}^p dt \right), \end{aligned} \tag{5.6}$$

whenever $u_0 \in E_{1-\frac{1}{p}, p}$ and $g \in L^p((0, T) \rightarrow E_0)$.

Proof. Step 1. First, we prove the lemma for $[0, T] \subset \mathbb{R}_+$ replaced by a sufficiently short time interval $[t_0, t_0 + T_1] \subset [0, T]$, i.e., $0 \leq t_0 < t_0 + T_1 \leq T$ with $T_1 \in (0, \infty)$ small enough. Without loss of generality, we may assume $t_0 = 0$ and $0 < T_1 \leq T$.

Let us recall our notation and the continuous imbedding (cf. (4.13))

$$\begin{aligned} Y^{T_1} &= Y_{1-\frac{1}{p}}^p(0, T_1) = L^p((0, T_1) \rightarrow E_1) \cap W^{1,p}((0, T_1) \rightarrow E_0) \\ &\hookrightarrow C\left([0, T_1] \rightarrow E_{1-\frac{1}{p}, p}\right). \end{aligned} \tag{5.7}$$

It is easy to see that a function $v \in Y^{T_1}$ is a strict solution of the perturbed initial value problem (5.2) on $(0, T_1)$ if and only if it satisfies

$$\begin{aligned} \frac{dv}{dt} - Bv(t) &= Av(t) + g(t) \quad \text{for a.e. } t \in (0, T_1); \\ v(0) &= u_0 \in E_{1-\frac{1}{p}, p}, \end{aligned} \tag{5.8}$$

in the strict sense, again. Notice that $\tilde{g} = Av + g \in L^p((0, T_1) \rightarrow E_0)$. We observe that problem (5.8) has a unique strict solution $v \in Y^{T_1}$ as soon as we have shown that the affine self mapping $F : v \mapsto \hat{v} : Y^{T_1} \rightarrow Y^{T_1}$, defined by

$$\begin{aligned} \frac{d\hat{v}}{dt} - B\hat{v}(t) &= Av(t) + g(t) \quad \text{for a.e. } t \in (0, T_1); \\ \hat{v}(0) &= u_0 \in E_{1-\frac{1}{p}, p}, \end{aligned} \tag{5.9}$$

possesses a unique fixed point $v \in Y^{T_1}$. Obviously, such a fixed point must belong to the (closed) affine subspace

$$Y_{(u_0)}^{T_1} = \{v \in Y^{T_1} : v(0) = u_0\} \quad \text{of the Banach space } Y^{T_1};$$

hence, $Y_{(u_0)}^{T_1} = u_0 + Y_{(0)}^{T_1}$. Clearly, $Y_{(0)}^{T_1} = \{v \in Y^{T_1} : v(0) = 0\}$ is a closed vector subspace of Y^{T_1} . The former one inherits the norm from the latter.

Next, we prove that $F : v \mapsto \hat{v}$ is a contraction on $Y_{(u_0)}^{T_1}$. To this end, let $v_i \in Y_{(u_0)}^{T_1}$ be arbitrary and set $\hat{v}_i = F(v_i)$; $i = 1, 2$. The differences $z = v_1 - v_2$ and

$\hat{z} = \hat{v}_1 - \hat{v}_2$ are in $Y_{(0)}^{T_1}$ and, by (5.9), they satisfy

$$\begin{aligned} \frac{d\hat{z}}{dt} - B\hat{z}(t) &= Az(t) \quad \text{for a.e. } t \in (0, T_1); \\ \hat{z}(0) &= 0 \in E_{1-\frac{1}{p}, p}. \end{aligned} \quad (5.10)$$

By Remark 4.2, Part (a), the operator $B \in \text{MR}_p(E)$ satisfies (4.11) with a constant $M(p, E, B, T_1) \leq M(p, E, B, T) \equiv M_T < \infty$. Hence, we have

$$\int_0^{T_1} \left\| \frac{d\hat{z}}{dt} \right\|_{E_0}^p dt + \int_0^{T_1} \|B\hat{z}(t)\|_{E_0}^p dt \leq M_T \int_0^{T_1} \|Az(t)\|_{E_0}^p dt.$$

Now we estimate the integrand on the right-hand side of (5.5),

$$\begin{aligned} &\int_0^{T_1} \left\| \frac{d\hat{z}}{dt} \right\|_{E_0}^p dt + \int_0^{T_1} \|B\hat{z}(t)\|_{E_0}^p dt \\ &\leq M_T \int_0^{T_1} (\delta \|Bz(t)\|_{E_0} + C_\delta \|z(t)\|_{E_0})^p dt \\ &\leq 2^{p-1} M_T \left(\delta^p \int_0^{T_1} \|Bz(t)\|_{E_0}^p dt + C_\delta^p \int_0^{T_1} \|z(t)\|_{E_0}^p dt \right). \end{aligned} \quad (5.11)$$

The integrand in the second integral on the right-hand side is estimated by Hölder's inequality:

$$\begin{aligned} \|z(t)\|_{E_0} &= \left\| \int_0^t \frac{dz}{dt}(s) ds \right\|_{E_0} \leq \int_0^t \|z'(s)\|_{E_0} ds \\ &\leq \left(\int_0^t \left\| \frac{dz}{dt}(s) \right\|_{E_0}^p ds \right)^{1/p} \left(\int_0^t ds \right)^{1/p'} \quad \text{for all } t \in [0, T_1], \end{aligned}$$

where $p' = p/(p-1) \in (1, \infty)$. Here, we have used $z(0) = 0 \in E_0$. Hence,

$$\|z(t)\|_{E_0}^p \leq t^{p/p'} \left(\int_0^t \left\| \frac{dz}{dt}(s) \right\|_{E_0}^p ds \right).$$

After integration we thus obtain, thanks to $p/p' = p-1$,

$$\int_0^{T_1} \|z(t)\|_{E_0}^p dt \leq \frac{1}{p} T_1^p \int_0^{T_1} \left\| \frac{dz}{dt}(s) \right\|_{E_0}^p ds. \quad (5.12)$$

Of course, the same inequality is valid for $\hat{z} \in Y_{(0)}^{T_1}$ in place of the function z . We apply the last inequality to the right-hand side of (5.11),

$$\begin{aligned} &\int_0^{T_1} \left\| \frac{d\hat{z}}{dt} \right\|_{E_0}^p dt + \int_0^{T_1} \|B\hat{z}(t)\|_{E_0}^p dt \\ &\leq 2^{p-1} M_T \left(\delta^p \int_0^{T_1} \|Bz(t)\|_{E_0}^p dt + \frac{1}{p} C_\delta^p T_1^p \int_0^{T_1} \left\| \frac{dz}{dt} \right\|_{E_0}^p dt \right). \end{aligned} \quad (5.13)$$

The integrals on both sides containing the generator $B \in \text{MR}_p(E)$ are estimated as follows. First, there are constants $c_1, C_1 \in (0, \infty)$ and $c_2, C_2 \in \mathbb{R}_+$ such that the inequalities

$$c_1 \|u\|_{E_1} - c_2 \|u\|_{E_0} \leq \|Bu\|_{E_0} \leq C_1 \|u\|_{E_1} + C_2 \|u\|_{E_0} \quad \text{hold for all } u \in E_1.$$

Consequently, we have

$$c_1^p \|u\|_{E_1}^p \leq 2^{p-1} (\|Bu\|_{E_0}^p + c_2^p \|u\|_{E_0}^p) \quad \text{and}$$

$$\|Bu\|_{E_0}^p \leq 2^{p-1} (C_1^p \|u\|_{E_1}^p + C_2^p \|u\|_{E_0}^p) \quad \text{for all } u \in E_1.$$

Applying these inequalities to (5.13), we arrive at

$$\begin{aligned} & \int_0^{T_1} \left\| \frac{d\hat{z}}{dt} \right\|_{E_0}^p dt + 2^{-(p-1)} c_1^p \int_0^{T_1} \|\hat{z}(t)\|_{E_1}^p dt - c_2^p \int_0^{T_1} \|\hat{z}(t)\|_{E_0}^p dt \\ & \leq 2^{p-1} M_T \cdot 2^{p-1} C_1^p \delta^p \int_0^{T_1} \|z(t)\|_{E_1}^p dt \\ & \quad + 2^{p-1} M_T \left(2^{p-1} C_2^p \delta^p \int_0^{T_1} \|z(t)\|_{E_0}^p dt + \frac{1}{p} C_\delta^p T_1^p \int_0^{T_1} \left\| \frac{dz}{dt} \right\|_{E_0}^p dt \right). \end{aligned}$$

Finally, we estimate the integrals $\int_0^{T_1} \|\hat{z}(t)\|_{E_0}^p dt$ and $\int_0^{T_1} \|z(t)\|_{E_0}^p dt$ above by (5.12), thus obtaining

$$\begin{aligned} & \left(1 - \frac{1}{p} T_1^p c_2^p \right) \int_0^{T_1} \left\| \frac{d\hat{z}}{dt} \right\|_{E_0}^p dt + 2^{-(p-1)} c_1^p \int_0^{T_1} \|\hat{z}(t)\|_{E_1}^p dt \\ & \leq 2^{2(p-1)} C_1^p \delta^p M_T \int_0^{T_1} \|z(t)\|_{E_1}^p dt \\ & \quad + \frac{2^{p-1}}{p} T_1^p M_T (2^{p-1} C_2^p \delta^p + C_\delta^p) \int_0^{T_1} \left\| \frac{dz}{dt} \right\|_{E_0}^p dt. \end{aligned} \quad (5.14)$$

We finish this step by choosing first $\delta \in (0, 1)$ then $T_1 \in (0, T]$ small enough, such that

$$\begin{aligned} 2^{2(p-1)} C_1^p \delta^p M_T & \leq \frac{1}{2} \cdot 2^{-(p-1)} c_1^p \quad \text{and} \\ \frac{2^{p-1}}{p} T_1^p M_T (2^{p-1} C_2^p \delta^p + C_\delta^p) & \leq \frac{1}{2} \left(1 - \frac{1}{p} T_1^p c_2^p \right), \end{aligned}$$

respectively, or, equivalently,

$$0 < \delta \leq 2^{-(3p-2)/p} (c_1/C_1) M_T^{-1/p} \quad \text{and} \quad (5.15)$$

$$T_1^p [2^p M_T (2^{p-1} C_2^p \delta^p + C_\delta^p) + c_2^p] \leq p. \quad (5.16)$$

With these choices of δ and T_1 , we obtain

$$\|\hat{z}\|_{Y^{T_1}}^b \leq \frac{1}{2} \|z\|_{Y^{T_1}}^b \quad (5.17)$$

in the new, equivalent norm

$$\|u\|_{Y^{T_1}}^b := \left[\left(1 - \frac{1}{p} T_1^p c_2^p \right) \int_0^{T_1} \left\| \frac{du}{dt} \right\|_{E_0}^p dt + 2^{-(p-1)} c_1^p \int_0^{T_1} \|u(t)\|_{E_1}^p dt \right]^{1/p} \quad (5.18)$$

on the abstract Sobolev space $Y^{T_1} = Y_{1-\frac{1}{p}}^p(0, T_1) = L^p((0, T_1) \rightarrow E_1) \cap W^{1,p}((0, T_1) \rightarrow E_0)$; see (4.5) and below in (5.19). Inequality (5.17) shows that $F : v \mapsto \hat{v}$ is a contraction on $Y_{(u_0)}^{T_1} (\subset Y^{T_1})$ with the Lipschitz constant $\frac{1}{2}$ with respect to the new norm $\|\cdot\|_{Y^{T_1}}^b$. Consequently, problem (5.8) has a unique strict solution $v \in Y^{T_1}$; in fact, we have $v \in Y_{(u_0)}^{T_1}$.

The following estimate for v can be proved by the same arguments as those used in our proof of contraction above: There is a constant $\Gamma \equiv \Gamma(T_1) \in (0, \infty)$,

independent from $u_0 \in E_{1-\frac{1}{p},p}$ and $g \in L^p((0, T) \rightarrow E_0)$, such that

$$\begin{aligned} (\|v\|_{Y^{T_1}}^\sharp)^p &= \int_0^{T_1} \left\| \frac{dv}{dt} \right\|_{E_0}^p dt + \int_0^{T_1} \|v(t)\|_{E_1}^p dt \\ &\leq \Gamma \left(\|u_0\|_{E_{1-\frac{1}{p},p}}^p + \int_0^{T_1} \|g(t)\|_{E_0}^p dt \right). \end{aligned} \quad (5.19)$$

Recall that $\|\cdot\|_{Y^{T_1}}^\sharp$ is an equivalent norm on the Banach space Y^{T_1} ; cf. (4.5). In analogy with Remark 4.2, Part (b), we may take the constant $\Gamma \equiv \Gamma(T_1) > 0$ in (5.19) above to be the *smallest* nonnegative number $\Gamma \in \mathbb{R}_+$ for which (5.19) is valid. It is easy to see that $\Gamma \equiv \Gamma(T_1) \in \mathbb{R}_+$ is a *nondecreasing* function of time $T_1 \in (0, T]$ and $\Gamma > 0$. The last estimate, (5.19), easily implies (5.6) with T_1 in place of T . The imbedding (5.7) being continuous, by (5.19), there is another constant $\hat{\Gamma} \equiv \hat{\Gamma}(T_1) \in [1, \infty)$, independent from $u_0 \in E_{1-\frac{1}{p},p}$ and $g \in L^p((0, T) \rightarrow E_0)$, such that

$$\|v(T_1)\|_{E_{1-\frac{1}{p},p}}^p \leq \hat{\Gamma} \left(\|u_0\|_{E_{1-\frac{1}{p},p}}^p + \int_0^{T_1} \|g(t)\|_{E_0}^p dt \right). \quad (5.20)$$

Again, similarly to $\Gamma \equiv \Gamma(T_1) > 0$ in (5.19), we may take the constant $\hat{\Gamma} \equiv \hat{\Gamma}(T_1) > 0$ in (5.20) above to be the *smallest* number $\hat{\Gamma} \in [1, \infty)$ for which (5.20) is valid. It is now easy to see that also the constant $\hat{\Gamma} \equiv \hat{\Gamma}(T_1) \geq 1$ is a *nondecreasing* function of time $T_1 \in (0, T]$.

Step 2. We may take $T_1 = T/m$ sufficiently small in Step 1 above, where $m \in \mathbb{N}$ is a sufficiently large positive integer. Next, we replace the interval $[0, T_1]$ from Step 1 by any subinterval $[t_0, t_0 + T_1] = \mathcal{J}_k = [(k-1)T_1, kT_1]$ of $[0, T]$ of length T_1 for $k = 1, 2, \dots, m$; hence, $\cup_{k=1}^m \mathcal{J}_k = [0, T]$. We make use of the existence and uniqueness of a strict solution

$$v \in Y_{1-\frac{1}{p}}^p(t_0, t_0 + T_1) = L^p((t_0, t_0 + T_1) \rightarrow E_1) \cap W^{1,p}((t_0, t_0 + T_1) \rightarrow E_0)$$

of the perturbed initial value problem (5.2) in every subinterval \mathcal{J}_k ; $k = 1, 2, \dots, m$, together with the estimates (5.19) and (5.20) on \mathcal{J}_k , by Step 1. Thus, from (5.19) and (5.20) we obtain, respectively,

$$\begin{aligned} &\int_{(k-1)T_1}^{kT_1} \left\| \frac{dv}{dt} \right\|_{E_0}^p dt + \int_{(k-1)T_1}^{kT_1} \|v(t)\|_{E_1}^p dt \\ &\leq \Gamma \left(\|v((k-1)T_1)\|_{E_{1-\frac{1}{p},p}}^p + \int_{(k-1)T_1}^{kT_1} \|g(t)\|_{E_0}^p dt \right), \end{aligned} \quad (5.21)$$

$$\|v(kT_1)\|_{E_{1-\frac{1}{p},p}}^p \leq \hat{\Gamma} \left(\|v((k-1)T_1)\|_{E_{1-\frac{1}{p},p}}^p + \int_{(k-1)T_1}^{kT_1} \|g(t)\|_{E_0}^p dt \right). \quad (5.22)$$

We recall that $\hat{\Gamma} \geq 1$. Consequently, iterating inequalities (5.22) for $k = 1, 2, \dots, \ell$, $1 \leq \ell \leq m$, we arrive at

$$\begin{aligned} \|v(\ell T_1)\|_{E_{1-\frac{1}{p},p}}^p &\leq \hat{\Gamma}^\ell \|u_0\|_{E_{1-\frac{1}{p},p}}^p + \sum_{k=1}^{\ell} \hat{\Gamma}^{\ell-k+1} \int_{(k-1)T_1}^{kT_1} \|g(t)\|_{E_0}^p dt \\ &\leq \hat{\Gamma}^\ell \left(\|u_0\|_{E_{1-\frac{1}{p},p}}^p + \int_0^{\ell T_1} \|g(t)\|_{E_0}^p dt \right). \end{aligned} \quad (5.23)$$

Next, we sum inequalities (5.21) for $k = 1, 2, \dots, m$, thus obtaining

$$\begin{aligned} \left(\|v\|_{Y_T}^\sharp\right)^p &= \int_0^T \left\|\frac{dv}{dt}\right\|_{E_0}^p dt + \int_0^T \|v(t)\|_{E_1}^p dt \\ &\leq \Gamma \left(\sum_{k=1}^m \|v((k-1)T_1)\|_{E_{1-\frac{1}{p},p}}^p\right) + \Gamma \int_0^T \|g(t)\|_{E_0}^p dt. \end{aligned} \quad (5.24)$$

To estimate the first summand on the right-hand side from above, we apply (5.23) with $\ell = k - 1$ for $k = 1, 2, \dots, m$, thus arriving at

$$\begin{aligned} &\sum_{k=1}^m \|v((k-1)T_1)\|_{E_{1-\frac{1}{p},p}}^p \\ &\leq \left(\sum_{\ell=0}^{m-1} \hat{\Gamma}^\ell\right) \|u_0\|_{E_{1-\frac{1}{p},p}}^p + \sum_{\ell=0}^{m-1} \hat{\Gamma}^\ell \int_0^{\ell T_1} \|g(t)\|_{E_0}^p dt \\ &\leq m \hat{\Gamma}^{m-1} \|u_0\|_{E_{1-\frac{1}{p},p}}^p + (m-1) \hat{\Gamma}^{m-1} \int_0^{(m-1)T_1} \|g(t)\|_{E_0}^p dt \\ &\leq \hat{M} \Gamma^{-1} \left(\|u_0\|_{E_{1-\frac{1}{p},p}}^p + \int_0^T \|g(t)\|_{E_0}^p dt\right), \end{aligned}$$

where $\hat{M} = m \hat{\Gamma}^{m-1} \Gamma \in [1, \infty)$ is a constant independent from v . We apply this estimate to the right-hand side of (5.24) to obtain

$$\left(\|v\|_{Y_T}^\sharp\right)^p \leq \hat{M} \|u_0\|_{E_{1-\frac{1}{p},p}}^p + (\hat{M} + \Gamma) \int_0^T \|g(t)\|_{E_0}^p dt. \quad (5.25)$$

We conclude the proof by applying (5.5) with $B, A \in \mathcal{L}(E_1 \rightarrow E_0)$ to the left-hand side of (5.6),

$$\begin{aligned} &\int_0^T \left\|\frac{dv}{dt}\right\|_{E_0}^p dt + \int_0^T \|(B+A)v(t)\|_{E_0}^p dt \\ &\leq \int_0^T \left\|\frac{dv}{dt}\right\|_{E_0}^p dt + (1+\delta) \int_0^T \|Bv(t)\|_{E_0}^p dt + C_\delta \int_0^T \|v(t)\|_{E_0}^p dt \\ &\leq \int_0^T \left\|\frac{dv}{dt}\right\|_{E_0}^p dt + (1+\delta) \|B\|_{\mathcal{L}(E_1 \rightarrow E_0)} \int_0^T \|v(t)\|_{E_1}^p dt + C_\delta \int_0^T \|v(t)\|_{E_0}^p dt \\ &\leq M_\delta \left(\|v\|_{Y_T}^\sharp\right)^p, \end{aligned}$$

by (5.24), with a constant $M_\delta \in (0, \infty)$ independent from v . Now we apply (5.25) to the last estimate to arrive at the desired inequality (5.6) with the constant $\tilde{M} = M_\delta(\hat{M} + \Gamma) > 0$. We have proved that the operator $B_A = B + A \in \mathcal{L}(E_1 \rightarrow E_0)$ belongs to the class $\text{MR}_p(E)$. \square

5.2. Proof of analyticity in time. Now we are ready to prove that any global strict solution $w : [0, T] \rightarrow U \subset E_{1-\frac{1}{p},p}$ to problem (4.8) that satisfies the hypotheses of Theorem 4.7 above must be analytic in time $t \in (0, T)$. Let us recall that a *strict solution* to problem (4.8) has been introduced in Definition 4.4. Indeed, below we will prove a more detailed result on a complex analytic (i.e., holomorphic) extension of $u(t)$ from the real time interval $(0, T) \subset \mathbb{R} \subset \mathbb{C}$ to the open complex domain $\Delta_{\vartheta}^{T',T}$ which is the intersection of the (open) triangle $\Delta_{\vartheta}^{(T')}$ with the (open) complex strip $\mathfrak{T}^{(r)}$ defined in (1.6), (1.7), and (1.8), respectively, where $\vartheta \in (0, \pi/2)$ is a given

angle and $0 < T' \leq T < \infty$. Here, the constants $\vartheta \in (0, \pi/2)$ and $T' \in (0, T]$ will be chosen sufficiently small, but positive; hence, we have $(0, T) \subset \Delta_{\vartheta}^{T', T}$. Finally, we denote by $\bar{\Delta}_{\vartheta}^{T', T}$ the closure of $\Delta_{\vartheta}^{T', T}$ in \mathbb{C} .

In addition to (H4) and (H5), we assume that A and f satisfy the following analyticity hypotheses (cf. Lunardi [65, Chapt. 8], §8.3.3, p. 308):

5.3. Hypothesis. Recall that both spaces, E_0 and E_1 , in the Banach couple $E = (E_0, E_1)$ are assumed to be complex Banach spaces (over the field \mathbb{C}) with $E_1 \hookrightarrow E_0$ densely and continuously. Furthermore, we assume that there are positive constants $\vartheta_0 \in (0, \pi/2)$ and $T_0 \in (0, T]$, and open sets $\mathcal{U} \subset \mathbb{C}$ and $\tilde{U} \subset E_{1-\frac{1}{p}, p}$ containing the compact set $\bar{\Delta}_{\vartheta_0}^{T_0, T}$ and the open set U , respectively, i.e., $\bar{\Delta}_{\vartheta_0}^{T_0, T} \subset \mathcal{U} \subset \mathbb{C}$ and $U \subset \tilde{U} \subset E_{1-\frac{1}{p}, p}$, such that

(H4') $A : [0, T] \times U \rightarrow \mathcal{L}(E_1 \rightarrow E_0)$ possesses a holomorphic extension $\tilde{A} : \mathcal{U} \times \tilde{U} \rightarrow \mathcal{L}(E_1 \rightarrow E_0)$ to the complex domain $\mathcal{U} \times \tilde{U}$ which satisfies $\tilde{A}(t, v) \in \text{MR}_p(E)$ for all $(t, v) \in \mathcal{U} \times \tilde{U}$.

(H5') $f : [0, T] \times U \rightarrow E_0$ possesses a holomorphic extension $\tilde{f} : \mathcal{U} \times \tilde{U} \rightarrow E_0$ to the complex domain $\mathcal{U} \times \tilde{U}$.

Again, the metric on $\mathcal{U} \times \tilde{U}$ is induced by the canonical norm on $\mathbb{C} \times E_{1-\frac{1}{p}, p}$. A precise definition of a *holomorphic* (i.e., *complex analytic*) mapping $\mathcal{F} : \mathcal{O} \subset \mathcal{X} \rightarrow \mathcal{Y}$ from an open subset \mathcal{O} of a complex Banach space \mathcal{X} into another complex Banach space \mathcal{Y} is given in Deimling [22, Definition 15.1, p. 150] (see also [22, Proposition 15.2, p. 150]).

Without assuming (H4) and (H5), we observe that (H4') and (H5') still guarantee the following claims, respectively: Given any compact set $K \subset \mathcal{U}$ and any continuous function $z : [0, T] \rightarrow K$, one can easily verify that both substitution mappings,

$$\begin{aligned} v &\mapsto [t \mapsto A(z(t), v(z(t)))] : \\ &C(K \rightarrow \tilde{U}) \rightarrow \mathcal{L}(L^p((0, T) \rightarrow E_1) \rightarrow L^p((0, T) \rightarrow E_0)) \quad \text{and} \\ v &\mapsto [t \mapsto f(z(t), v(z(t)))] : C(K \rightarrow \tilde{U}) \rightarrow L^p((0, T) \rightarrow E_0), \end{aligned}$$

the former one meaning that

$$\begin{aligned} (v, u) &\mapsto [t \mapsto A(z(t), v(z(t))) u(z(t))] : \\ &C(K \rightarrow \tilde{U}) \times L^p((0, T) \rightarrow E_1) \rightarrow L^p((0, T) \rightarrow E_0), \end{aligned}$$

are locally Lipschitz continuous, the former one with values in $\mathcal{L}(L^p((0, T) \rightarrow E_1) \rightarrow L^p((0, T) \rightarrow E_0))$ and the latter one with values in $L^p((0, T) \rightarrow E_0)$. We will take advantage of this local Lipschitz continuity in our proof of Theorem 5.3 below. We remark that the operator norm in $\mathcal{L}(L^p((0, T) \rightarrow E_1) \rightarrow L^p((0, T) \rightarrow E_0))$ of the linear substitution operator

$$u \mapsto [t \mapsto A(z(t), v(z(t))) u(t)] : L^p((0, T) \rightarrow E_1) \rightarrow L^p((0, T) \rightarrow E_0),$$

with $z \in C([0, T] \rightarrow K)$ and $v \in C(K \rightarrow \tilde{U})$ being fixed, is bounded above by the supremum norm

$$\begin{aligned} \| \| A(z(\cdot), v(z(\cdot))) \| \|_{L^\infty(0, T)} &:= \| [t \mapsto A(z(t), v(z(t)))] \|_{C([0, T] \rightarrow \mathcal{L}(E_1 \rightarrow E_0))} \\ &= \sup_{0 \leq t \leq T} \| A(z(t), v(z(t))) \|_{\mathcal{L}(E_1 \rightarrow E_0)} \quad (< \infty). \end{aligned}$$

Theorem 5.3. *Let $1 < p < \infty$, $\vartheta_0 \in (0, \pi/2)$, $0 < T_0 \leq T < \infty$, and assume that $g \in L^p((0, T) \rightarrow E_0)$ possesses a holomorphic extension $\tilde{g} : \mathcal{U} \rightarrow E_0$ to an open set $\mathcal{U} \subset \mathbb{C}$ containing $\bar{\Delta}_{\vartheta_0}^{T_0, T}$, i.e., $\bar{\Delta}_{\vartheta_0}^{T_0, T} \subset \mathcal{U} \subset \mathbb{C}$. Assume that \tilde{U} is a nonempty open subset of $E_{1-\frac{1}{p}, p}$ and \tilde{A} and \tilde{f} satisfy (H4') and (H5'), respectively, and their respective restrictions $A = \tilde{A}|_{[0, T] \times U}$ and $f = \tilde{f}|_{[0, T] \times U}$ to $[0, T] \times U \subset \mathbb{R} \times E_{1-\frac{1}{p}, p}$ satisfy (H4) and (H5) with an open set $U \subset \tilde{U} \subset E_{1-\frac{1}{p}, p}$. Finally, assume that $w : [0, T] \rightarrow U \subset E_{1-\frac{1}{p}, p}$ is a fixed global strict solution to problem (4.8) (hence, satisfying (4.17)) with a prescribed initial value $w(0) = w_0 \in U$ and such that $w(t) \in U$ and $\tilde{A}(t, w(t)) \in \text{MR}_p(E)$ for all $t \in [0, T]$.*

Then there exist constants $\vartheta' \in (0, \vartheta_0]$ and $T' \in (0, T_0]$, small enough, and a holomorphic function $\tilde{w} : \Delta_{\vartheta'}^{T', T} \rightarrow E_{1-\frac{1}{p}, p}$ with the following two properties:

- (a) $\tilde{w}(t) \in \tilde{U}$ for every $t \in \Delta_{\vartheta'}^{T', T}$ and \tilde{w} verifies the abstract nonlinear evolutionary problem (4.8) in the complex domain $\Delta_{\vartheta'}^{T', T}$, i.e.,

$$\begin{aligned} \frac{du}{dt} - \tilde{A}(t, u(t))u(t) &= \tilde{f}(t, u(t)) + \tilde{g}(t) \quad \text{for every } t \in \Delta_{\vartheta'}^{T', T}; \\ \lim_{t \rightarrow 0, t \in \Delta_{\vartheta'}^{T', T}} u(t) &= w_0 \in E_{1-\frac{1}{p}, p}. \end{aligned} \tag{5.26}$$

- (b) $\tilde{w}(t) = w(t)$ holds for a.e. $t \in (0, T)$.

Such a holomorphic extension $\tilde{w} : \Delta_{\vartheta'}^{T', T} \rightarrow \tilde{U} \subset E_{1-\frac{1}{p}, p}$ of $w : (0, T) \rightarrow U \subset E_{1-\frac{1}{p}, p}$ from $(0, T)$ to $\Delta_{\vartheta'}^{T', T}$ is unique.

Before proceeding to prove this theorem, we clarify our notation with the open sets U and \tilde{U} in $E_{1-\frac{1}{p}, p}$ as follows.

Remark 5.4. We need to take advantage of our (H4) and (H5) (with an open set $U \subset E_{1-\frac{1}{p}, p}$) and (H4') and (H5') (with another open set $\tilde{U} \subset E_{1-\frac{1}{p}, p}$) only for the values of $v = w(t) \in U \subset \tilde{U}$ ($t \in \bar{\Delta}_{\vartheta'}^{T', T}$) near the (compact) image $K = \{w(t) \in E_{1-\frac{1}{p}, p} : t \in [0, T]\}$ of the (continuous) curve $w : [0, T] \rightarrow E_{1-\frac{1}{p}, p}$. Indeed, (H4') and (H5') imply that both holomorphic extensions $\tilde{A} : \mathcal{U} \times \tilde{U} \rightarrow \mathcal{L}(E_1 \rightarrow E_0)$ and $\tilde{f} : \mathcal{U} \times \tilde{U} \rightarrow E_0$ of $A : [0, T] \times U \rightarrow \mathcal{L}(E_1 \rightarrow E_0)$ and $f : [0, T] \times U \rightarrow E_0$, respectively, are locally Lipschitz continuous. Consequently, the Cartesian product $[0, T] \times K$ being compact in the complex Banach space $\mathbb{C} \times E_{1-\frac{1}{p}, p}$, we use a finite open subcover by open balls to find two bounded open sets $\mathcal{U} \subset \mathbb{C}$ and $U = \tilde{U} \subset E_{1-\frac{1}{p}, p}$, such that both mappings \tilde{A} and \tilde{f} are Lipschitz continuous in $\mathcal{U} \times \tilde{U}$. We conclude that, in our proof of Theorem 5.3 below, we may assume that $[0, T] \subset \mathcal{U} \subset \mathbb{C}$ and $K \subset U = \tilde{U} \subset E_{1-\frac{1}{p}, p}$ with both \mathcal{U} and U being open and bounded. In particular, if the numbers $\vartheta_0 \in (0, \pi/2)$ and $T_0 \in (0, T]$ are taken sufficiently small, then we have also $\bar{\Delta}_{\vartheta_0}^{T_0, T} \subset \mathcal{U}$ together with $\tilde{w}(t) \in U$ for all $t \in \bar{\Delta}_{\vartheta'}^{T', T}$, provided $\vartheta' \in (0, \vartheta_0]$ and $T' \in (0, T_0]$ are small enough. Consequently, $\bar{\Delta}_{\vartheta'}^{T', T} \subset \bar{\Delta}_{\vartheta_0}^{T_0, T} \subset \mathcal{U}$. To simplify our notation, we work only with the holomorphic extensions $\tilde{g} : \mathcal{U} \rightarrow E_0$, $\tilde{A} : \mathcal{U} \times U \rightarrow \mathcal{L}(E_1 \rightarrow E_0)$, and $\tilde{f} : \mathcal{U} \times U \rightarrow E_0$

of the mappings g , A , and f , respectively. Hence, we may remove the “tilde” from these symbols and write simply $g = \tilde{g}$, $A = \tilde{A}$, and $f = \tilde{f}$. We also may and will assume that both mappings A and f are Lipschitz continuous in all of $\mathcal{U} \times \tilde{\mathcal{U}}$.

In our construction of the continuous extension $\tilde{w} : \bar{\Delta}_{\vartheta_0}^{T_0, T} \rightarrow \mathbb{C}$ of the strict solution $w : [0, T] \rightarrow \mathbb{C}$, holomorphic in $\Delta_{\vartheta_0}^{T_0, T}$, we take advantage of a factorization approach for the complex time variable $t = \varrho\mu$ where $\varrho \in (0, \tau_0)$ and $\mu \in \mathbb{C}$ with $|\mu - 1| < \sin \vartheta$. The numbers $\tau_0 \in (0, T)$ and $\vartheta \in (0, \vartheta_0)$ are suitable constants. Fixing such a constant μ , we obtain a *mild solution*, $\omega \equiv \omega_\mu : [0, \tau_0] \rightarrow U \subset E_{1-\frac{1}{p}, p}$, of the corresponding initial value problem with the real time variable $t = \varrho \in [0, \tau_0]$. Of course, this solution depends on the complex parameter μ from the open disc

$$D_r(1) := \{\mu \in \mathbb{C} : |\mu - 1| < r\}$$

centered at the point $1 \in \mathbb{C}$ with radius $r = \sin \vartheta$. We will complete the proof by showing that the mild solution, ω , is holomorphic with respect to μ . This factorization approach has been used earlier in Henry [40, Chapt. 3, §3.4] and Lunardi [65, Chapt. 8, §8.3.3].

Proof of Theorem 5.3. Given any two numbers $\vartheta \in (0, \pi/2)$ and $\tau_0 \in (0, \infty)$, we define a bounded open *sector* in the complex plane \mathbb{C} by

$$\mathfrak{A}_\vartheta^{(\tau_0)} := \{t = \varrho\mu \in \mathbb{C} : 0 < \varrho < \tau_0 \text{ and } \mu \in \mathbb{C} \text{ with } |\mu - 1| < \sin \vartheta\} \quad (5.27)$$

with vertex at the origin $0 \in \mathbb{C}$ and angle 2ϑ . Its closure in \mathbb{C} is denoted by $\bar{\mathfrak{A}}_\vartheta^{(\tau_0)}$. Recalling our definition of the triangle $\Delta_\vartheta^{(T)}$ by (1.6), and setting $r = \sin \vartheta$ (hence, $0 < r < 1$), we deduce that

$$\Delta_{\vartheta'}^{(T_1)} \subset \mathfrak{A}_\vartheta^{(\tau_0)} \subset \Delta_\vartheta^{(T_2)}$$

holds whenever

$$0 < T_1 \leq \tau_0, \quad 0 < \vartheta' < \arctan r, \quad (1+r)\tau_0 \leq T_2 < \infty.$$

Following this factorization of the complex time $t \in \mathbb{C}$ in $t = \varrho\mu$ with $\varrho \in (0, \tau_0)$ and $\mu \in D_r(1) = \{\mu \in \mathbb{C} : |\mu - 1| < r\}$, so that $0 < 1 - r < \Re \mu < 1 + r$ with $r = \sin \vartheta (< 1)$, we replace the complex time $t \in \Delta_\vartheta^{T', T}$ in the initial value problem (5.26) by the product $\mu t \in \mathbb{C}$ with $t \in (0, \tau_0)$ and $\mu \in D_r(1)$, where we will choose both $\vartheta \in (0, \pi/2)$ and $\tau_0 \in (0, \infty)$ sufficiently small, so that $\mathfrak{A}_\vartheta^{(\tau_0)} \subset \Delta_{\vartheta_0}^{T_0, T}$ holds, i.e., $\mu t \in \Delta_{\vartheta_0}^{T_0, T}$ for every pair $(t, \mu) \in (0, \tau_0) \times D_r(1)$. Hence, we must have

$$0 < \vartheta \leq \vartheta_0, \quad r\tau_0 \leq T_0 \cdot \tan \vartheta_0, \quad \text{and} \quad (1+r)\tau_0 \leq T.$$

Given a fixed number $\mu \in D_r(1) \subset \mathbb{C}$, we look for an unknown continuous mapping $\omega \equiv \omega_\mu : [0, \tau_0] \rightarrow U \subset E_{1-\frac{1}{p}, p}$, $\omega(t) \equiv \omega_\mu(t) = \tilde{w}(\mu t)$, that according to (5.26) must be a strict solution to the following evolutionary problem (with the tilde “ $\tilde{\sim}$ ”, marking holomorphic extensions, having been removed),

$$\begin{aligned} \frac{d\omega}{dt} - \mu A(\mu t, \omega(t))\omega(t) &= \mu [f(\mu t, \omega(t)) + g(\mu t)] \quad \text{for every } t \in (0, \tau_0); \\ \omega(0) &= w_0 \in E_{1-\frac{1}{p}, p}. \end{aligned} \quad (5.28)$$

Of course, for $\mu = 1$ we will have $\omega(t) \equiv \omega_1(t) = w(t)$ for a.e. $t \in (0, T)$, by uniqueness. We remark that, thanks to our hypothesis $w(t) \in U$ and $A(t, w(t)) \in \text{MR}_p(E)$ for all $t \in [0, T]$, we have also $\mu A(t, w(t)) \in \text{MR}_p(E)$ for all $t \in [0, T]$ and

all $\mu \in \mathbb{C}$ satisfying $|\mu - 1| < \sin \vartheta$ with $\vartheta \in (0, \pi/2)$ small enough, say, $0 < \vartheta < \pi/6$ in which case $|\mu - 1| < 1/2$. This claim follows easily from the perturbation lemma, Lemma 5.2, thanks to $\mu = 1 + \nu$ with $\nu \in \mathbb{C}$ satisfying $|\nu| < \sin \vartheta$. Even Lemma 5.1 would do if $\vartheta \in (0, \pi/2)$ were chosen sufficiently small. Clearly, it suffices to prove that, for each fixed $t \in (0, \tau_0)$, the function $\mu \mapsto \omega_\mu(t) : D_r(1) \rightarrow E_{1-\frac{1}{p},p}$ is holomorphic. This approach to the analyticity in time of solutions to semilinear parabolic problems can be found, e.g., in the monographs by Henry [40, Chapt. 3], §3.4, Theorem 3.4.4 and Corollary 3.4.6 on pp. 63–66, and by Lunardi [65, Chapt. 8], §8.3.3, p. 308.

Choosing $\vartheta \in (0, \pi/2)$ and $\tau_0 \in (0, T]$ small enough, such that $\bar{\mathfrak{A}}_\vartheta^{(\tau_0)} \subset \mathcal{U}$, and recalling $r = \sin \vartheta \in (0, 1)$, we abbreviate

$$F(t, v, \mu) := \mu[f(\mu t, v) + g(\mu t)] \quad \text{for all } (t, v, \mu) \in [0, \tau_0] \times U \times D_r(1). \quad (5.29)$$

By (H5'), the mapping $(v, \mu) \mapsto F(t, v, \mu) : U \times D_r(1) \rightarrow E_0$ is holomorphic for each $t \in [0, \tau_0]$, with all partial derivatives of F with respect to v and μ being continuous in $[0, \tau_0] \times U \times D_r(1)$. According to Amann [6, Chapt. III, §4.10], pp. 180–191, and Clément and Li [20, Sect. 2], p. 18, given a fixed parameter value $\mu \in D_r(1)$, every strict solution $\omega \equiv \omega_\mu \in Y_{1-\frac{1}{p}}^p(0, \tau_0)$ of the initial value problem (5.28) satisfies the following integral equation for the unknown function $\omega \equiv \omega_\mu \in Y_{1-\frac{1}{p}}^p(0, \tau_0)$,

$$\omega(t) = \mathcal{F}(t, \omega, \mu) \quad \text{for every } t \in [0, \tau_0], \quad (5.30)$$

with the right-hand side equal to

$$\begin{aligned} &\mathcal{F}(t, v, \mu) \\ &:= e^{\mu t A(0, w_0)} w_0 + \mu \int_0^t e^{\mu(t-s)A(0, w_0)} [A(\mu s, v(s)) - A(0, w_0)] v(s) \, ds \\ &\quad + \int_0^t e^{\mu(t-s)A(0, w_0)} F(s, v(s), \mu) \, ds \quad \text{for every } t \in [0, \tau_0] \end{aligned} \quad (5.31)$$

and for all $v \in Y_{1-\frac{1}{p}}^p(0, \tau_0)$ satisfying $v(t) \in U$ for every $t \in [0, \tau_0]$. In contrast to defining a contraction mapping using the (unique) strict solution to prove local existence in Theorem 4.5, in the case of problem (5.28) we prefer to use the (unique) mild solution defined by an integral representation (variation-of-constants formula); cf. (5.31). The equivalence between strict and mild solutions is treated in Ball [10], Henry [40, Chapt. 3], and Pazy [72, Theorem on p. 259].

Clearly, (5.30) is a fixed point equation for the unknown function $\omega \in Y_{1-\frac{1}{p}}^p(0, \tau_0)$. Here, one can choose $\tau_0 \in (0, T_1]$, where $T_1 \in (0, T]$ and $\vartheta \in (0, \pi/2)$ are sufficiently small, such that $\bar{\mathfrak{A}}_\vartheta^{(T_1)} \subset \mathcal{U}$ and the mapping

$$\Phi \equiv \Phi_\mu : v \mapsto [t \mapsto \mathcal{F}(t, v, \mu)]$$

is a contraction in a closed ball

$$\Sigma_{\rho_1, T_1}^{(w_0)} = \{v \in Y^{T_1} : v(0) = w_0 \quad \text{and} \quad \|v - w\|_{[0, T_1]} \|_{Y^{T_1}} \leq \rho_1\}$$

in the Banach space

$$Y^{T_1} = Y_{1-\frac{1}{p}}^p(0, T_1) = L^p((0, T_1) \rightarrow E_1) \cap W^{1,p}((0, T_1) \rightarrow E_0)$$

of radius $\rho_1 \in (0, \infty)$ centered at the point $w \in Y^{T_1}$. As usual, the function $w|_{[0, T_1]} \in Y^{T_1}$ denotes the restriction to $[0, T_1]$ of the strict solution $w \in Y^T =$

$Y_{1-\frac{1}{p}}^p(0, T)$ from the hypotheses of our theorem. The proof of this contraction property follows the same ideas and steps as the proof of Theorem 4.5 taken from Clément and Li [20, Theorem 2.1, p. 19]. The reader is referred to Prüss [74, pp. 9–13], Chapt. 3, for further details. Notice that the numbers $\rho_1 \in (0, \infty)$, $T_1 \in (0, T]$, and $\vartheta \in (0, \pi/2)$, if chosen small enough, such that the contraction holds with the Lipschitz constant $\frac{1}{2}$, are independent from the particular choice of the parameter $\mu \in D_r(1)$ since $\bar{\mathfrak{A}}_{\vartheta}^{(T_1)}$ is a compact subset of \mathcal{U} ; cf. our remarks before this proof (Remark 5.4) that remain valid also for the compact set $\bar{\mathfrak{A}}_{\vartheta}^{(T_1)} \times K$ in the complex Banach space $\mathbb{C} \times E_{1-\frac{1}{p}, p}$. Of course, $r = \sin \vartheta \in (0, 1)$ is sufficiently small, and both $\rho_1 \in (0, \infty)$ and $T_1 \in (0, T]$ must be also so small that $v(t) \in U$ holds for all $t \in [0, T_1]$, whenever $v \in \Sigma_{\rho_1, T_1}^{(w_0)}$. Finally, the constants ρ_1 , T_1 , and ϑ can be chosen independent from $w_0 \in K$, so that one may use them in any time interval $[t_0, t_0 + T_1] \subset [0, T]$ of sufficiently short length $T_1 > 0$; the initial condition $w(0) = w_0 \in K$ at $t = 0$ is replaced by $w(t_0) \in K$ at arbitrary time $t_0 \in [0, T - T_1]$.

Next, we analyze the holomorphy properties of the fix point mapping

$$\mathcal{F} : [0, T_1] \times \Sigma_{\rho_1, T_1}^{(w_0)} \times D_r(1) \rightarrow \Sigma_{\rho_1, T_1}^{(w_0)}$$

defined in (5.31) where we may take $\tau_0 = T_1$; more precisely, those of the mapping

$$(v, \mu) \mapsto \mathcal{F}(t, v, \mu) : \Sigma_{\rho_1, T_1}^{(w_0)} \times D_r(1) \rightarrow \Sigma_{\rho_1, T_1}^{(w_0)},$$

for each fixed $t \in [0, T_1]$. To begin with, for $0 \leq s < t \leq T_1$, $\mu \in D_r(1)$, and $v \in \Sigma_{\rho_1, T_1}^{(w_0)}$, we rewrite

$$\begin{aligned} & A(\mu s, v(s)) - A(0, w_0) \\ &= \{I - [\lambda I - A(\mu s, v(s))][\lambda I - A(0, w_0)]^{-1}\}[\lambda I - A(0, w_0)], \end{aligned}$$

where $\lambda \in (0, \infty)$ is large enough in order to guarantee that the (bounded) linear operator $\lambda I - A(0, w_0) : E_1 \rightarrow E_0$ has a bounded inverse $[\lambda I - A(0, w_0)]^{-1} : E_0 \rightarrow E_1$, and observe that the function (integrand)

$$\mu \mapsto e^{\mu(t-s)A(0, w_0)} [A(\mu s, v(s)) - A(0, w_0)] v(s) : D_r(1) \rightarrow E_0$$

is holomorphic and so is the integral

$$\mu \mapsto \int_0^t e^{\mu(t-s)A(0, w_0)} [A(\mu s, v(s)) - A(0, w_0)] v(s) ds : D_r(1) \rightarrow E_0.$$

We have used here the fact that the operator-valued function

$$\mu \mapsto e^{\mu(t-s)A(0, w_0)} : D_r(1) \rightarrow \mathcal{L}(E_0 \rightarrow E_0)$$

is holomorphic for any fixed numbers $s, t \in \mathbb{R}$ satisfying $0 \leq s < t \leq T_1$. Similarly, the function

$$\mu \mapsto e^{\mu(t-s)A(0, w_0)} F(s, v(s), \mu) : D_r(1) \rightarrow E_0$$

being holomorphic, so is the integral

$$\mu \mapsto \int_0^t e^{\mu(t-s)A(0, w_0)} F(s, v(s), \mu) ds : D_r(1) \rightarrow E_0.$$

We conclude that the sum

$$\mu \mapsto \mathcal{F}(t, v, \mu) : D_r(1) \rightarrow E_0$$

defined by (5.31) with $\tau_0 = T_1$ is holomorphic for every $t \in [0, T_1]$.

Finally, from the fixed point equation (5.30) we deduce that the function $\omega \equiv \omega_\mu : [0, \tau_0] \rightarrow U \subset E_{1-\frac{1}{p}, p}$, which is continuous thanks to $\omega \in \Sigma_{\rho_1, T_1}^{(w_0)} \subset Y^{T_1} = Y_{1-\frac{1}{p}}^P(0, T_1)$, is holomorphic in the variable $\mu \in D_r(1)$. Although this holomorphy claim follows directly from a well-known result in Deimling [22, Theorem 15.3, Chapt. 4, §15, p. 151], cf. also Krantz and Parks [58, Theorem 6.1.2, §6.1, p. 118], we sketch a constructive proof below for the sake of completeness.

Indeed, any standard proof of the Banach fixed point theorem for the (*contractive*) self mapping

$$\Phi \equiv \Phi_\mu : v \mapsto [t \mapsto \mathcal{F}(t, v, \mu)] : \Sigma_{\rho_1, T_1}^{(w_0)} \rightarrow \Sigma_{\rho_1, T_1}^{(w_0)}$$

shows that, given an arbitrary “initial” function $\varphi_0 \in \Sigma_{\rho_1, T_1}^{(w_0)}$, the iterates

$$\varphi_n = \Phi(\varphi_{n-1}) = \Phi^2(\varphi_{n-2}) = \dots = \Phi^{n-1}(\varphi_1) = \Phi^n(\varphi_0); \quad \text{for } n = 1, 2, 3, \dots,$$

form a Cauchy sequence in $\Sigma_{\rho_1, T_1}^{(w_0)}$ which converges to the unique fixed point $\omega \equiv \omega_\mu$ of Φ , namely, $\varphi_n \rightarrow \omega$ in $\Sigma_{\rho_1, T_1}^{(w_0)} \subset Y^{T_1}$ as $n \rightarrow \infty$. The convergence is uniform for $\mu \in D_r(1)$. Recalling the continuous imbedding (5.7), we have also $\varphi_n(t) \rightarrow \omega(t)$ in $E_{1-\frac{1}{p}, p}$ as $n \rightarrow \infty$, uniformly for $t \in [0, T_1]$ and $\mu \in D_r(1)$. Choosing $\varphi_0 = w|_{[0, T_1]}$, a function of time $t \in [0, T_1]$ which does not depend on the parameter $\mu \in D_r(1)$, we observe that each iterate

$$\varphi_n(t) = \mathcal{F}(t, \varphi_{n-1}, \mu); \quad n = 1, 2, 3, \dots, \quad t \in [0, T_1],$$

is a holomorphic function in the variable (parameter) $\mu \in D_r(1)$. Applying Osgood’s theorem and the Cauchy integral formula for discs to each iterate $\varphi_n(t)$ (see e.g. Krantz [57], Theorem 1.2.2 (p. 24), or John [50], Chapt. 3, Sect. 3(c), eq. (3.22c), p. 71), we conclude that also the limit function $\omega \equiv \omega_\mu$ is holomorphic in the variable $\mu \in D_r(1)$ and satisfies $\Phi(\omega) = \omega$.

We have thus verified that the strict solution $w : [0, T] \rightarrow U \subset E_{1-\frac{1}{p}, p}$ of problem (4.8) possesses a holomorphic extension to the bounded open sector $\mathfrak{A}_\vartheta^{(T_1)}$. In fact, we have proved that this claim is valid in any time shift of this sector by a number $t_0 \in [0, T - T_1]$, that is, in any sector

$$t_0 + \mathfrak{A}_\vartheta^{(T_1)} := \{t = t_0 + \varrho\mu \in \mathbb{C} : 0 < \varrho < T_1 \text{ and } |\mu - 1| < \sin \vartheta\}$$

with vertex at the point $t_0 \in \mathbb{C}$ and angle 2ϑ . We apply the last result with t_0 ranging from 0 to $T - T_1$ over the interval $[0, T - T_1]$ to conclude that the function $w : [0, T] \rightarrow U \subset E_{1-\frac{1}{p}, p}$ possesses a holomorphic extension to the bounded open set

$$\cup_{t_0 \in [0, T - T_1]} (t_0 + \mathfrak{A}_\vartheta^{(T_1)}) \subset \mathbb{C}$$

which contains the open complex domain $\Delta_{\vartheta', T}^{T', T}$ defined in (1.7), whenever $T' = T_1$ and $0 < \vartheta' \leq \arctan(\sin \vartheta)$, owing to $\Delta_{\vartheta'}^{(T_1)} \subset \mathfrak{A}_\vartheta^{(T_1)}$.

Hence, we have proved that there are constants $\vartheta' \in (0, \vartheta_0]$ and $T' \in (0, T_0]$, small enough, and a holomorphic function $\tilde{w} : \Delta_{\vartheta'}^{T', T} \rightarrow E_{1-\frac{1}{p}, p}$ with the desired properties (a) and (b) in the conclusion of our theorem. Since $(0, T) \subset \Delta_{\vartheta'}^{T', T}$, such a holomorphic function \tilde{w} must be unique. The proof is complete. \square

6. ANALYTICITY IN SPACE FOR THE CAUCHY PROBLEM IN $\mathbb{R}^N \times (0, T)$

In the previous two sections, Sections 4 and 5, we have treated the initial value problem (4.8) for a strict solution $u : (0, T) \rightarrow E_0$ with the initial condition u_0 in the real interpolation space $E_{1-\frac{1}{p},p} \equiv (E_0, E_1)_{1-\frac{1}{p},p}$ between E_1 and E_0 , $E_1 \hookrightarrow E_{1-\frac{1}{p},p} \hookrightarrow E_0$.

By Theorem 4.7, such a strict solution belongs to the abstract Sobolev space $Y^T = Y_{1-\frac{1}{p}}^p(0, T)$ introduced in (4.10). Hence, we have $u(t) \in E_{1-\frac{1}{p},p}$ for every $t \in [0, T]$.

In this section we replace the triplet of abstract (complex) Banach spaces $E_1 \hookrightarrow E_{1-\frac{1}{p},p} \hookrightarrow E_0$ by the following complex Sobolev, Besov, and Lebesgue spaces, respectively,

$$W^{2m,p}(\mathbb{R}^N) \hookrightarrow B^{s;p,p}(\mathbb{R}^N) \hookrightarrow L^p(\mathbb{R}^N), \quad s = 2m\left(1 - \frac{1}{p}\right) \in (0, 2m),$$

where

$$\begin{aligned} B^{s;p,p}(\mathbb{R}^N) &:= (L^p(\mathbb{R}^N), W^{2m,p}(\mathbb{R}^N))_{s/(2m),p} \\ &= (L^p(\mathbb{R}^N), W^{2m,p}(\mathbb{R}^N))_{1-(1/p),p} \end{aligned}$$

is the Besov space obtained by real interpolation (see, e.g., [1, 65, 84]). We recall that $2 + \frac{N}{m} < p < \infty$ which guarantees $(s-m)p > N$ and, thus, the Sobolev(-Besov) imbeddings $B^{s-m;p,p}(\mathbb{R}^N) \hookrightarrow C^0(\mathbb{R}^N) \cap L^\infty(\mathbb{R}^N)$ and $B^{s;p,p}(\mathbb{R}^N) \hookrightarrow C^m(\mathbb{R}^N) \cap W^{m,\infty}(\mathbb{R}^N)$ are continuous.

Throughout this section we restrict ourselves to the easiest case of *analytic initial conditions* that we are able to treat in our present work.

6.1. **Hypothesis.** We use the following assumptions:

(H6) The initial data $\mathbf{u}_0 : \mathbb{R}^N \rightarrow \mathbb{C}^M$, $\mathbf{u}_0 = (u_{0,1}, u_{0,2}, \dots, u_{0,M})$, can be extended to a holomorphic function $\tilde{\mathbf{u}}_0 = (\tilde{u}_{0,1}, \tilde{u}_{0,2}, \dots, \tilde{u}_{0,M}) : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M$ in a complex strip $\mathfrak{X}^{(r)} \subset \mathbb{C}^N$ defined in (2.1), for some $r \in (0, \infty)$, such that every component $\tilde{u}_{0,j} : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}$; $j = 1, 2, \dots, M$, has the following properties:

(H6.1) the function $x \mapsto \tilde{u}_{0,j}(x + iy) : \mathbb{R}^N \rightarrow \mathbb{C}$ is in the (complex) Besov space $B^{s;p,p}(\mathbb{R}^N)$,

(H6.2) the Besov norm $\|\tilde{u}_{0,j}(\cdot + iy)\|_{B^{s;p,p}(\mathbb{R}^N)}$ is uniformly bounded for all $y \in Q^{(r)}$, and

(H6.3) $y \mapsto \tilde{u}_{0,j}(\cdot + iy) : Q^{(r)} \rightarrow B^{s;p,p}(\mathbb{R}^N)$ is continuously (partially) differentiable with respect to the parameter $y = (y_1, \dots, y_N) \in Q^{(r)} = \{y \in \mathbb{R}^N : |y|_\infty < r\}$; $j = 1, 2, \dots, M$.

Equivalently, the function $x \mapsto \tilde{\mathbf{u}}_0(x + iy) : \mathbb{R}^N \rightarrow \mathbb{C}^M$ belongs to the Cartesian product $\mathbf{B}^{s;p,p}(\mathbb{R}^N) = [B^{s;p,p}(\mathbb{R}^N)]^M$, its norm $\|\tilde{\mathbf{u}}_0(\cdot + iy)\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)}$ satisfies (cf. (3.10))

$$\mathfrak{N}^{(r)}(\tilde{\mathbf{u}}_0) = \sup_{y \in Q^{(r)}} \|\tilde{\mathbf{u}}_0(\cdot + iy)\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)} < \infty,$$

and it is continuously differentiable with respect to the parameter $y \in Q^{(r)}$.

The “shift” isometry $\|\tilde{\mathbf{u}}_0(\cdot + x_0 + iy_0)\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)} = \|\tilde{\mathbf{u}}_0(\cdot + iy_0)\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)}$ is obvious for all pairs $(x_0, y_0) \in \mathbb{R}^N \times Q^{(r)}$, i.e., for all complex numbers $z_0 = x_0 + iy_0 \in \mathfrak{X}^{(r)}$.

The restriction in (H6) is motivated by the following approximation property of the Sobolev and Besov spaces, see e.g. Triebel [84, Chapt. 2].

Remark 6.1. The Fréchet space $\mathcal{S}(\mathbb{R}^N)$ of all complex-valued, rapidly decreasing infinitely differentiable functions $\varphi : \mathbb{R}^N \rightarrow \mathbb{C}$ being dense in all of the spaces $L^p(\mathbb{R}^N)$, $W^{2m,p}(\mathbb{R}^N)$, $B^{s;p,p}(\mathbb{R}^N)$, and $L^2(\mathbb{R}^N)$, by [84, Chapt. 2], §2.3, Theorem 2.3.2 on p. 172, we take $\mathbf{u}_0 : \mathbb{R}^N \rightarrow \mathbb{C}^M$ so smooth and rapidly decreasing near infinity that its holomorphic extension $\tilde{\mathbf{u}}_0 : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M$ satisfies even the following stronger regularity condition: The family of functions $x \mapsto \tilde{\mathbf{u}}_0(x + iy) : \mathbb{R}^N \rightarrow \mathbb{C}^M$, parametrized by $y \in Q^{(r)}$, belongs to a bounded subset of

$$\mathbf{L}^2(\mathbb{R}^N) \cap \mathbf{W}^{2m,p}(\mathbb{R}^N) \quad \text{for some } p \in \mathbb{R}, \quad 2 + \frac{N}{m} < p < \infty.$$

For instance, all complex linear combinations of *Hermite functions* form a dense vector subspace \mathcal{V} of the Fréchet space $\mathcal{S}(\mathbb{R}^N)$, by Reed and B. Simon [75, Chapt. V, §3], Theorem V.13 on p. 143. *Hermite functions* are entire complex functions $h : \mathbb{C}^N \rightarrow \mathbb{C}$ of the form

$$h(z) = P(z_1, z_2, \dots, z_N) \cdot \exp\left(-\frac{1}{2} \sum_{i=1}^N z_i^2\right) \quad \text{for } z = (z_i)_{i=1}^N = x + iy \in \mathbb{C}^N,$$

where $P(z)$ is a complex polynomial in N complex variables $z_i \in \mathbb{C}$; $i = 1, 2, \dots, N$, see [75, p. 142]. One may take functions from \mathcal{V} as components $\tilde{u}_{0,j}$ of $\tilde{\mathbf{u}}_0$; $j = 1, 2, \dots, M$. Indeed, notice that

$$\begin{aligned} \left| \exp\left(-\frac{1}{2} \sum_{i=1}^N z_i^2\right) \right| &= \exp\left(-\frac{1}{2} \sum_{i=1}^N x_i^2 + \frac{1}{2} \sum_{i=1}^N y_i^2\right) \\ &\leq \exp\left(\frac{1}{2} N r^2\right) \cdot \exp\left(-\frac{1}{2} |x|_2^2\right) \end{aligned}$$

holds for all $z = x + iy \in \mathfrak{X}^{(r)} \subset \mathbb{C}^N$, where

$$|x|_2 = \left(\sum_{i=1}^N |x_i|^2\right)^{1/2}, \quad |y|_\infty = \max_{i=1,2,\dots,N} |y_i| < r.$$

It is well known that all three vector spaces $\mathcal{V} \subset \mathcal{S}(\mathbb{R}^N) \subset L^2(\mathbb{R}^N)$ are invariant under the (unitary) Fourier transformation $\mathcal{F} : L^2(\mathbb{R}^N) \rightarrow L^2(\mathbb{R}^N)$. (We always consider the *unitary* Fourier transformation \mathcal{F} as described in Stein and Weiss [80, Chapt. I].) Consequently, if the Fourier transform $\mathcal{F}u_{0,j} : \mathbb{R}^N \rightarrow \mathbb{C}$ of each component of $\mathbf{u}_0 : \mathbb{R}^N \rightarrow \mathbb{C}^M$ decays at least exponentially fast at infinity, then the holomorphic extension of the function $u_{0,j} : \mathbb{R}^N \rightarrow \mathbb{C}$ to a complex strip $\mathfrak{X}^{(r)} \subset \mathbb{C}^N$, for some $r \in (0, \infty)$, is easily obtained in the form of the inverse Fourier-Laplace transform $\mathcal{F}^{-1}(\mathcal{F}u_{0,j}) : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}$ of $\mathcal{F}u_{0,j}$, by the classical Paley-Wiener-Schwartz theory, see e.g. Hörmander [44, Theorem 7.4.2, p. 192] or Stein and Weiss [80, Chapt. III], §2, pp. 91–101, and §6.12, pp. 127–128. An interested reader is referred to Takáč [82, Chapt. 5] for a brief review of the (inverse) Fourier-Laplace transform that applies to our current setting.

In regard to later applications (cf. Proposition 6.5 and Theorem 7.1), in our (H6) above we have not specified the number $r \in (0, \infty)$ corresponding to the half width of the complex strip $\mathfrak{X}^{(r)} = \mathbb{R}^N + iQ^{(r)}$, a tube in \mathbb{C}^N with the base $Q^{(r)} = (-r, r)^N$. Hypotheses (H1)–(H3) in Section 3 show that only the case

$0 < r \leq r_0$ is useful. We will comment on the choice of $r \in (0, r_0]$ in Remark 7.2 right after Theorem 7.1 below. Concerning this question of choosing (finding) a suitable half-width $r \in (0, r_0]$, we begin with the following observation.

Remark 6.2. The Hermite functions $h : \mathbb{C}^N \rightarrow \mathbb{C}$ described in Remark 6.1 are *not* the only way for approximating the initial values $\mathbf{u}(\cdot, t_0) = \mathbf{u}_0 \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ at time $t_0 = 0$ in the Besov space within the Besov norm $\|\cdot\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)}$. In our approximation procedure we need to guarantee the following “uniformity” of the half width of the complex strip $\mathfrak{X}^{(r)} = \mathbb{R}^N + iQ^{(r)}$, i.e., the same half-width $r \in (0, r_0]$ for each approximating function $\tilde{\mathbf{u}}_{0,n} : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M$; $n = 1, 2, 3, \dots$. In precise analytic terms, this means that, for any given radius $R_1 \in (0, \infty)$ of the ball $B_{R_1}(0)$ in $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$, there exists a number $r_1 \in (0, r_0]$ small enough, such that, whenever $r \in (0, r_1]$, the approximating sequence of functions $\{\tilde{\mathbf{u}}_{0,n}\}_{n=1}^\infty$ has the following properties (cf. (H6)):

- (H7.1) each function $x \mapsto \tilde{\mathbf{u}}_{0,n}(x+iy) : \mathbb{R}^N \rightarrow \mathbb{C}^M$ is in the Besov space $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$, for every $y \in Q^{(r)}$,
- (H7.2) the “proximity to \mathbf{u}_0 ” estimate $\|\tilde{\mathbf{u}}_{0,n}(\cdot + iy) - \mathbf{u}_0\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)} < R_1$ holds for all $y \in Q^{(r)}$ and $n = 1, 2, 3, \dots$,
- (H7.3) $\tilde{\mathbf{u}}_{0,n} : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M$ is holomorphic for every $n = 1, 2, 3, \dots$, and finally
- (H7.4) the restrictions $\mathbf{u}_{0,n} = \tilde{\mathbf{u}}_{0,n}|_{\mathbb{R}}$ of $\tilde{\mathbf{u}}_{0,n}$ to the real line \mathbb{R} satisfy $\|\mathbf{u}_{0,n} - \mathbf{u}_0\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)} \rightarrow 0$ as $n \rightarrow \infty$.

We keep the natural notation $\mathbf{L}^2(\mathbb{R}^N) = [L^2(\mathbb{R}^N)]^M$ etc. introduced for spaces of vector-valued functions in the Introduction (Section 1). We recall the continuous Sobolev(-Besov) imbeddings

$$\begin{aligned} \mathcal{S}(\mathbb{R}^N) &\hookrightarrow L^2(\mathbb{R}^N) \cap W^{2m,p}(\mathbb{R}^N) \hookrightarrow B^{s;p,p}(\mathbb{R}^N) \hookrightarrow C^m(\mathbb{R}^N) \cap W^{m,\infty}(\mathbb{R}^N), \\ W^{2m,p}(\mathbb{R}^N) &\hookrightarrow B^{s;p,p}(\mathbb{R}^N) \hookrightarrow L^p(\mathbb{R}^N) \cap L^\infty(\mathbb{R}^N), \quad 2 + \frac{N}{m} < p < \infty. \end{aligned}$$

We remark that $W^{2m,p}(\mathbb{R}^N) \not\subset L^2(\mathbb{R}^N)$ if $2 < p < \infty$. From now on we identify $\mathbf{u}_0 \equiv \tilde{\mathbf{u}}_0$ and drop the tilde “~” in the (unique) holomorphic extension.

By (H1)–(H3) (cf. Theorem 3.4), let us set $r = r_0 \in (0, \infty)$ above. In the Cauchy problem (1.1) we may replace the real space variable $x \in \mathbb{R}^N$ by its complex shift $z = x + x_0 + iy_0$ by a fixed complex vector $z_0 = x_0 + iy_0 \in \mathfrak{X}^{(r)} \subset \mathbb{C}^N$ with any $x_0 \in \mathbb{R}^N$ and any $y_0 \in Q^{(r)}$. In the sequel we consider $z_0 \in \mathfrak{X}^{(r)}$ to be a parameter and $x \in \mathbb{R}^N$ an independent variable in the Cauchy problem (1.1) spatially “shifted” by z_0 ,

$$\begin{aligned} \frac{\partial \mathbf{u}}{\partial t} + \mathbf{P}\left(x + z_0, t, \frac{1}{i} \frac{\partial}{\partial x}\right) \mathbf{u} &= \mathbf{f}\left(x + z_0, t; \left(\frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta}\right)_{|\beta| \leq m}\right) \\ &\text{for } (x, t) \in \mathbb{R}^N \times (0, T); \end{aligned} \tag{6.1}$$

$$\mathbf{u}(x, 0) = \mathbf{u}_0(x + z_0) \quad \text{for } x \in \mathbb{R}^N.$$

By our hypothesis on the initial data $\mathbf{u}_0 : \mathbb{R}^N \rightarrow \mathbb{C}^M$ and its holomorphic extension $\mathbf{u}_0 \equiv \tilde{\mathbf{u}}_0 : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M$ stated above, for each $z_0 \in \mathfrak{X}^{(r)}$, the “shifted” function $x \mapsto \mathbf{u}_0^{(z_0)}(x) := \mathbf{u}_0(x + z_0) : \mathbb{R}^N \rightarrow \mathbb{C}^M$ belongs to $\mathbf{L}^2(\mathbb{R}^N) \cap \mathbf{W}^{2m,p}(\mathbb{R}^N)$, where $2 + \frac{N}{m} < p < \infty$. Consequently, we have $\mathbf{u}_0^{(z_0)} \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$, and thus, we may apply the local (in time) existence and uniqueness result (Theorem 4.5) on a short time interval $[t_0, T_1] \subset [0, T]$ with the initial condition $\mathbf{u}(\cdot, t_0) = \mathbf{u}_0^{(z_0)}$ in $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$ at time $t = t_0 \in [0, T]$ to conclude that the spatially “shifted” Cauchy problem (6.1)

possesses a unique weak solution $\mathbf{u}^{(z_0)} \in C([t_0, T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$, local in time. Of course, the length of the time interval $[t_0, T_1]$ depends on the shift $z_0 \in \mathfrak{X}^{(r)}$; more precisely, on its imaginary part $y_0 = \Im z_0 \in Q^{(r)}$. However, when making use of the abstract reformulation (4.8) of the Cauchy problem (6.1), we must guarantee that the values of the (unique) strict solution $\mathbf{u}^{(z_0)} : [t_0, T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ to the Cauchy problem (6.1), the (continuous) “shifted” function $t \mapsto \mathbf{u}(\cdot + z_0, t) = \mathbf{u}^{(z_0)}(\cdot, t) : [t_0, T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$, stay in the bounded open set $U = \tilde{U} \subset E_{1-\frac{1}{p},p}$ for all times $t \in [t_0, T_1]$ (cf. Remark 5.4). To avoid this technical problem, we make the following *global existence* hypothesis, cf. Theorem 3.4:

6.2. Hypothesis.

(H8) The original Cauchy problem (1.1), i.e., problem (6.1) with $z_0 = 0$ and the initial data $\mathbf{u}_0 = \hat{\mathbf{u}}_0 \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ at $t = 0$, possesses a global weak solution $\hat{\mathbf{u}} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$.

Now define the set $U \subset E_{1-\frac{1}{p},p} = \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ that appears in (H4), (H5), (H4’), (H5’), as follows: First, let

$$U_0 = \text{conv}\left(B_{R_0}(\mathbf{0}) \cup \cup_{t \in [0, T]} B_{R_0}(\hat{\mathbf{u}}(\cdot, t))\right) \subset \mathbf{B}^{s;p,p}(\mathbb{R}^N) \quad (6.2)$$

be the convex hull of the union of open balls

$$\begin{aligned} B_{R_0}(\mathbf{0}) &:= \{\mathbf{w} \in \mathbf{B}^{s;p,p}(\mathbb{R}^N) : \|\mathbf{w}\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)} < R_0\} \subset \mathbf{B}^{s;p,p}(\mathbb{R}^N), \\ B_{R_0}(\mathbf{v}) &:= \mathbf{v} + B_{R_0}(\mathbf{0}) \text{ with } \mathbf{v} = \hat{\mathbf{u}}(\cdot, t) \text{ for } t \in [0, T], \end{aligned}$$

where their radius $R_0 \in (0, \infty)$ is an arbitrary positive number. Of course, the symbol “ $\mathbf{0}$ ” stands for the zero function in $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$. Alternatively, we may take $U_0 = B_{R_0}(\hat{\mathbf{u}}_0)$ to be any open ball in $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$ centered at $\hat{\mathbf{u}}_0 = \hat{\mathbf{u}}(\cdot, 0)$ with (sufficiently large) radius $R_0 \in (0, \infty)$, such that $\mathbf{0} \in B_{R_0}(\hat{\mathbf{u}}_0)$ and $\hat{\mathbf{u}}(\cdot, t) \in B_{R_0}(\hat{\mathbf{u}}_0)$ holds for every $t \in [0, T]$. However, this choice of R_0 would not fit in Example 9.2 in Section 9 below. Clearly, U_0 is a bounded open set in $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$. From now on we take the initial values $\mathbf{u}^{(z_0)}(\cdot, t_0) = \mathbf{u}_0^{(z_0)} = \mathbf{u}^{(z_0)}(\cdot + z_0, t_0) \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ at time $t = t_0 \in [0, T)$ in the Cauchy problem (6.1) (and similar related initial value problems) from the set U_0 only, i.e., $\mathbf{u}_0^{(z_0)} \in U_0$. This choice will guarantee that the values of the (unique) strict solution $\mathbf{u}^{(z_0)} : [t_0, T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ to the “shifted” Cauchy problem (6.1) stay for all times $t \in [t_0, T_1]$ in the bounded open set $U = \tilde{U} \subset E_{1-\frac{1}{p},p} = \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ defined next, $U \supset U_0$. We put

$$U = \cup \{B_{R_0}(\mathbf{v}) : \mathbf{v} \in U_0\} \subset \mathbf{B}^{s;p,p}(\mathbb{R}^N); \quad (6.3)$$

hence,

$$U = \{\mathbf{w} \in \mathbf{B}^{s;p,p}(\mathbb{R}^N) : \|\mathbf{w} - \mathbf{v}\|_{\mathbf{B}^{s;p,p}(\mathbb{R}^N)} < R_0 \text{ for some } \mathbf{v} \in U_0\}.$$

One may call U the *open R_0 -neighborhood* of U_0 in $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$. Also $U = \tilde{U}$ is a bounded, open, and convex set in the complex Besov space $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$ and, consequently, in $\mathbf{W}^{m,p}(\mathbb{R}^N)$ and in $\mathbf{C}_{\text{bdd}}^m(\mathbb{R}^N) = \mathbf{C}^m(\mathbb{R}^N) \cap \mathbf{W}^{m,\infty}(\mathbb{R}^N)$, as well, owing to the continuous Sobolev(-Besov) imbeddings

$$B^{s;p,p}(\mathbb{R}^N) \hookrightarrow W^{m,p}(\mathbb{R}^N) \quad \text{and} \quad B^{s;p,p}(\mathbb{R}^N) \hookrightarrow C^m(\mathbb{R}^N) \cap W^{m,\infty}(\mathbb{R}^N), \quad (6.4)$$

respectively, where

$$\left(1 + \frac{N}{2m + N}\right)m < s = 2m\left(1 - \frac{1}{p}\right) < 2m$$

thanks to the inequalities $2 + \frac{N}{m} < p < \infty$; see Adams and Fournier [1, Chapt. 7], Theorem 7.34(a,c), p. 231. This shows that, for any function $\mathbf{w} \in U$, the partial derivatives $\frac{\partial^{|\beta|}\mathbf{w}}{\partial x^\beta}$, $\beta = (\beta_1, \dots, \beta_N) \in (\mathbb{Z}_+)^N$, of order $|\beta| = \beta_1 + \dots + \beta_N \leq m$, are uniformly bounded on \mathbb{R}^N ,

$$\left|\frac{\partial^{|\beta|}\mathbf{w}(x)}{\partial x^\beta}\right| \leq C \equiv C(U) = \text{const} < \infty \quad \text{for all } x \in \mathbb{R}^N, \quad (6.5)$$

for a constant $C \in \mathbb{R}_+$ depending solely on U . These partial derivatives are arguments in the reaction function $\mathbf{f}(x, t; \left(\frac{\partial^{|\beta|}\mathbf{u}}{\partial x^\beta}\right)_{|\beta| \leq m})$ on the right-hand side of (1.1) and (6.1). In (H3) on \mathbf{f} we take $\Sigma \subset \mathbb{C}^{M\tilde{N}}$ to be the closed polydisc $\Sigma = [\bar{D}_C(0)]^{M\tilde{N}}$ where $\bar{D}_C(0) := \{z \in \mathbb{C} : |z| \leq C\}$ is a closed disc. This restriction on the values of the (unique) strict solution to the bounded open set $U = \tilde{U} \subset E_{1-\frac{1}{p}, p}$ will be used in applications to semilinear Heston-type models in ‘‘Mathematical Finance’’ treated in Section 9.

From (H3) we deduce immediately that each component $f_j : \bar{\Omega} \times [\bar{D}_C(0)]^{M\tilde{N}} \rightarrow \mathbb{C}$ of the reaction function $\mathbf{f} = (f_1, \dots, f_M)$ is continuously differentiable (i.e., of class C^1) with the time derivative $\frac{\partial}{\partial t} f_j(x, t; X)$ and all argument first-order partial derivatives

$$\frac{\partial f_j}{\partial X_{\beta, k}}(x, t; X), \quad \text{for } |\beta| \leq m \text{ and } j, k = 1, 2, \dots, M,$$

being uniformly bounded on $\Omega \times \Sigma$. Consequently, each f_j is Lipschitz continuous with respect to the variables t and $X_{\beta, k}$, uniformly on $\Omega \times \Sigma$.

Recalling the continuous Sobolev(-Besov) imbeddings (6.4), i.e.,

$$B^{s; p, p}(\mathbb{R}^N) \hookrightarrow W^{m, p}(\mathbb{R}^N) \cap C^m(\mathbb{R}^N) \cap W^{m, \infty}(\mathbb{R}^N),$$

and the L^p -integrability condition in (3.4), we have just proved the following lemma (cf. (H4), (H5), (H4’), (H5’)):

Lemma 6.3. *Assume that $\mathbf{f} : \bar{\Omega} \times \mathbb{C}^{M\tilde{N}} \rightarrow \mathbb{C}^M$ satisfies (H3), and (H8) is also satisfied. Let $U = \tilde{U} \subset E_{1-\frac{1}{p}, p} = \mathbf{B}^{s; p, p}(\mathbb{R}^N)$ be as in (6.3). Then the Nemytskii operator $\mathbf{F} : [0, T] \times U \rightarrow E_0 = \mathbf{L}^p(\mathbb{R}^N)$ defined by*

$$\mathbf{F}(t, \mathbf{v})(x) := \mathbf{f}\left(x, t; \left(\frac{\partial^{|\beta|}\mathbf{v}}{\partial x^\beta}\right)_{|\beta| \leq m}\right), \quad x \in \mathbb{R}^N, \quad (6.6)$$

for all $t \in [0, T]$ and all $\mathbf{v} \in U$, satisfies the following properties:

(a) $\mathbf{F} : [0, T] \times U \rightarrow E_0$ is a Lipschitz continuous mapping, i.e., \mathbf{F} satisfies (H5).

(b) The substitution mapping $\mathfrak{F} : C([0, T] \rightarrow U) \rightarrow L^p((0, T) \rightarrow E_0)$ defined by

$$\mathfrak{F}(\mathbf{v})(t) := \mathbf{F}(t, \mathbf{v}(t)), \quad \text{for all } t \in [0, T], \mathbf{v} \in C([0, T] \rightarrow U),$$

is Lipschitz continuous with values in $L^p((0, T) \rightarrow E_0)$.

Proof. The only claims in Parts (a) and (b) that remain to be verified are that \mathbf{F} maps $[0, T] \times U$ into E_0 and \mathfrak{F} maps $C([0, T] \rightarrow U)$ into $L^p((0, T) \rightarrow E_0)$, respectively.

(a) For each component F_j of $\mathbf{F} = (F_1, \dots, F_M)$ from (6.6) we derive

$$\begin{aligned}
 & F_j(t, \mathbf{v}_1)(x) - F_j(t, \mathbf{v}_2)(x) \\
 &= \sum_{|\beta| \leq m} \sum_{k=1}^M \left[\int_0^1 \frac{\partial f_j}{\partial X_{\beta,k}} \left(x, t; \left((1-\theta) \frac{\partial^{|\beta|} \mathbf{v}_1}{\partial x^\beta} + \theta \frac{\partial^{|\beta|} \mathbf{v}_2}{\partial x^\beta} \right)_{|\beta| \leq m} \right) d\theta \right] \\
 & \times \frac{\partial^{|\beta|}}{\partial x^\beta} (v_{1,k}(x) - v_{2,k}(x)), \quad x \in \mathbb{R}^N, j = 1, 2, \dots, M,
 \end{aligned} \tag{6.7}$$

for all $t \in [0, T]$ and for all $\mathbf{v}_1, \mathbf{v}_2 \in U$. Notice that the partial derivatives with respect to x^β emerge from the chain rule applied to the right-hand side of (6.6) using the partial derivatives $\frac{\partial}{\partial X_{\beta,k}} f_j(x, t; X)$ with respect to the argument $X_{\beta,k}$ and the convex combination $\mathbf{w} = (1-\theta)\mathbf{v}_1 + \theta\mathbf{v}_2 \in U$ for $0 \leq \theta \leq 1$, thanks to U being convex. Consequently, with this abbreviation for \mathbf{w} and our choice of the constant $C \equiv C(U)$ in (6.5), all partial derivatives $\frac{\partial^{|\beta|} \mathbf{w}}{\partial x^\beta}$, $\beta = (\beta_1, \dots, \beta_N) \in (\mathbb{Z}_+)^N$, of order $|\beta| = \beta_1 + \dots + \beta_N \leq m$, are uniformly bounded on \mathbb{R}^N , by (6.5). By (H3), cf. Remark 3.1, all partial derivatives $\frac{\partial}{\partial X_{\beta,k}} f_j(x, t; \cdot) : \Sigma \rightarrow \mathbb{C}$ are uniformly bounded,

$$\left| \frac{\partial f_j(x, t; X)}{\partial X_{\beta,k}} \right| \leq C_1 \equiv C_1(C(U)) = \text{const} < \infty \tag{6.8}$$

(cf. (3.5)) for all $(x, t) \in \Omega$ and all $X = ((X_{\beta,k})_{|\beta| \leq m})_{k=1}^M \in \Sigma$, by a constant $C_1 \in \mathbb{R}_+$ depending solely on $C(U)$. We apply these estimates to the integrands in (6.7) to conclude that there is a Lipschitz constant $L \equiv L(U) \in \mathbb{R}_+$ depending solely on U (through the constant $C_1(C(U)) \geq 0$ in (6.8) above), such that

$$|\mathbf{F}(t, \mathbf{v}_1)(x) - \mathbf{F}(t, \mathbf{v}_2)(x)| \leq L \sum_{|\beta| \leq m} \sum_{k=1}^M \left| \frac{\partial^{|\beta|}}{\partial x^\beta} (v_{1,k}(x) - v_{2,k}(x)) \right| \tag{6.9}$$

for all $x \in \mathbb{R}^N$, $t \in [0, T]$ and $\mathbf{v}_1, \mathbf{v}_2 \in U$.

We recall $U \subset B^{s,p,p}(\mathbb{R}^N)$ and the imbeddings in (6.4) to deduce from (6.9) that the mappings $\mathbf{F}(t, \cdot) : U \rightarrow E_0 = \mathbf{L}^p(\mathbb{R}^N)$ are uniformly Lipschitz continuous (with the same Lipschitz constant) for all $t \in [0, T]$. Here, we single out the special case of $\mathbf{v}_1 = \mathbf{v} \in U$ being arbitrary and $\mathbf{v}_2 = \mathbf{0} \in U$, i.e., $\mathbf{v}_2(x) \equiv \mathbf{0} \in \mathbb{C}^M$ for all $x \in \mathbb{R}^N$. Then (6.9) yields

$$|\mathbf{F}(t, \mathbf{v})(x)| \leq |\mathbf{F}(t, \mathbf{0})(x)| + L \sum_{|\beta| \leq m} \sum_{k=1}^M \left| \frac{\partial^{|\beta|}}{\partial x^\beta} v_k(x) \right|, \tag{6.10}$$

for all $x \in \mathbb{R}^N$, $t \in [0, T]$ and $\mathbf{v} \in U$, where

$$\mathbf{F}(t, \mathbf{0})(x) = \mathbf{f}(x, t; \vec{\mathbf{0}}), \quad x \in \mathbb{R}^N, \quad \vec{\mathbf{0}} = (0)_{|\beta| \leq m} \equiv (0, \dots, 0) \in \mathbb{C}^{M\tilde{N}}.$$

satisfies $\mathbf{F}(t, \mathbf{0}) \in E_0 = \mathbf{L}^p(\mathbb{R}^N)$, by the L^p -integrability condition in (3.4), i.e.,

$$\|\mathbf{F}(t, \mathbf{0})\|_{E_0} = \left(\int_{\mathbb{R}^N} |\mathbf{f}(x, t; \vec{\mathbf{0}})|^p dx \right)^{1/p} \leq K \quad \text{for all } t \in [0, T],$$

where $K \in (0, \infty)$ is a constant. Now it follows from (6.10) above that also $\mathbf{F}(t, \mathbf{v}) \in E_0$ holds for all $t \in [0, T]$ and for all $\mathbf{v} \in U$, as claimed.

(b) Analogous results for the mapping $\mathfrak{F} : C([0, T] \rightarrow U) \rightarrow L^p([0, T] \rightarrow E_0)$ follow from those we have just proved for $\mathbf{F}(t, \cdot) : U \rightarrow E_0$, $t \in [0, T]$. Namely, the

“supremum” (or “maximum”) norm on the Banach space $C([0, T] \rightarrow E_{1-\frac{1}{p}, p})$ of all continuous functions $u : [0, T] \rightarrow E_{1-\frac{1}{p}, p}$ is defined by

$$\|u\|_{L^\infty(0, T)} := \|u\|_{C([0, T] \rightarrow E_{1-\frac{1}{p}, p})} = \sup_{t \in [0, T]} \|u(t)\|_{E_{1-\frac{1}{p}, p}}. \quad \square$$

Remark 6.4. Analogous result (to Lemma 6.3) hold for the “shifted” Nemytskii operator $\mathbf{F}^{(z_0)} : [0, T] \times U \rightarrow E_0 = \mathbf{L}^p(\mathbb{R}^N)$, by a complex vector $z_0 \in \mathfrak{X}^{(r)}$, defined by

$$\mathbf{F}^{(z_0)}(t, \mathbf{v})(x) := \mathbf{f}\left(x + z_0, t; \left(\frac{\partial^{|\beta|} \mathbf{v}}{\partial x^\beta}\right)_{|\beta| \leq m}\right), \quad x \in \mathbb{R}^N, \quad (6.11)$$

for all $t \in [0, T]$ and for all $\mathbf{v} \in U$. Both constants, $C_1 \equiv C_1(C(U))$ and $L \equiv L(U)$ in inequalities (6.8) and (6.9), respectively, are independent from the shift by $z_0 \in \mathfrak{X}^{(r)}$ in case $x \in \mathbb{R}^N$ is replaced by $x + z_0$, thanks to $(x, t) \in \Omega$ where the domain $\Omega = \Gamma_T^{(T_0)}(r_0, \vartheta_0) = \mathfrak{X}^{(r_0)} \times \Delta_{\vartheta_0}^{T_0, T} \subset \mathbb{C}^N \times \mathbb{C}$ has been introduced in Section 3 before (H1)–(H3) (cf. (3.8) and Theorem 3.4).

In what follows, the initial data \mathbf{u}_0 in the Cauchy problem (6.1) at time $t_0 \in [0, T]$ have nothing to do with the initial value $\hat{\mathbf{u}}(\cdot, 0) = \hat{\mathbf{u}}_0 \in \mathbf{B}^{s; p, p}(\mathbb{R}^N)$ of the global weak solution $\hat{\mathbf{u}} \in C([0, T] \rightarrow \mathbf{B}^{s; p, p}(\mathbb{R}^N))$ in (H8), except for the restriction $\mathbf{u}_0 \in U_0$, where U_0 is determined by the values of the solution $\hat{\mathbf{u}}(\cdot, t)$ in $\mathbf{B}^{s; p, p}(\mathbb{R}^N)$ for $0 \leq t \leq T$, see (6.2).

We take advantage of Remark 6.4 to recall that, given a holomorphic function $\mathbf{u}_0 \equiv \hat{\mathbf{u}}_0 : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M$ as described before (H8), the spatially “shifted” Cauchy problem (6.1) possesses a unique weak solution $\mathbf{u}^{(z_0)} \equiv \mathbf{u}^{(z_0, t_0)} \in C([t_0, T_1] \rightarrow \mathbf{B}^{s; p, p}(\mathbb{R}^N))$, local in time, for every fixed shift $z_0 \in \mathfrak{X}^{(r)}$, locally uniformly in the complex domain, provided its real and imaginary parts, $x_0 = \Re z_0, y_0 = \Im z_0 \in \bar{Q}^{(r_1)} \subset Q^{(r)}$, are small enough, i.e., $\max\{|x_0|_\infty, |y_0|_\infty\} \leq r_1$ ($r_1 < r = r_0$). Indeed, it suffices to choose $r_1 \in (0, r)$ so small that each “shifted” function $x \mapsto \mathbf{u}_0^{(z_0)}(x) = \mathbf{u}_0(x + z_0) : \mathbb{R}^N \rightarrow \mathbb{C}^M$ (serving as the initial data at time $t = t_0 \in [0, T]$), with $z_0 = x_0 + iy_0$ satisfying $\max\{|x_0|_\infty, |y_0|_\infty\} \leq r_1$, lies in the open set U_0 specified in (6.2) after (H8), i.e., $\mathbf{u}_0^{(z_0)} = \mathbf{u}_0(\cdot + z_0) \in U_0 \subset E_{1-\frac{1}{p}, p}$. Recall that $\bar{Q}^{(r_1)}$ and $\bar{\mathfrak{X}}^{(r_1)} = \mathbb{R}^N + i\bar{Q}^{(r_1)}$ stand for the respective closures of the cube $Q^{(r_1)} \subset \mathbb{R}^N$ and the strip $\mathfrak{X}^{(r_1)} \in \mathbb{C}^N$. Indeed, our choice of $r_1 \in (0, r)$ small enough to guarantee $\mathbf{u}_0^{(z_0)} = \mathbf{u}_0(\cdot + x_0 + iy_0) \in U_0$ for every $y_0 \in Q^{(r_1)}$ (and for all $x_0 \in \mathbb{R}^N$), is possible thanks to the closed cube $\bar{Q}^{(r_1)} \subset \mathbb{R}^N$ being compact. For “small” shifts $z_0 = x_0 + iy_0 \in \mathbb{C}^N$ we introduce the complex cube

$$\begin{aligned} Q_{\mathbb{C}}^{(r_1)} &:= Q^{(r_1)} + iQ^{(r_1)} = \{z = x + iy \in \mathbb{C} : \max\{|x|_\infty, |y|_\infty\} < r_1\} \\ &\subset \mathfrak{X}^{(r_1)} = \mathbb{R}^N + iQ^{(r_1)} \end{aligned}$$

and denote by $\bar{Q}_{\mathbb{C}}^{(r_1)}$ its closure in \mathbb{C}^N ; hence,

$$\begin{aligned} \bar{Q}_{\mathbb{C}}^{(r_1)} &= \bar{Q}^{(r_1)} + i\bar{Q}^{(r_1)} \\ &= \{z = x + iy \in \mathbb{C} : \max\{|x|_\infty, |y|_\infty\} \leq r_1\} \\ &\subset \bar{\mathfrak{X}}^{(r_1)} = \mathbb{R}^N + i\bar{Q}^{(r_1)} \subset \mathfrak{X}^{(r)}, \end{aligned}$$

thanks to $0 < r_1 < r$. We will call $\bar{Q}_{\mathbb{C}}^{(r_1)}$ the small shift cube.

To obtain a *spatially holomorphic* extension $\mathbf{u} \equiv \hat{\mathbf{u}} : \mathfrak{X}^{(r_1)} \times [0, T] \rightarrow \mathbb{C}^M$ of the function $\hat{\mathbf{u}} \in C([0, T] \rightarrow \mathbf{B}^{s; p, p}(\mathbb{R}^N))$, from (H8), to the spatio-temporal strip

$\mathfrak{X}^{(r_1)} \times [0, T]$, having a (unique) continuous extension $\mathbf{u} : \bar{\mathfrak{X}}^{(r_1)} \times [0, T] \rightarrow \mathbb{C}^M$ to the closed strip $\bar{\mathfrak{X}}^{(r_1)} \times [0, T]$, such that $(y_0, t) \mapsto \mathbf{u}(\cdot + iy_0, t) : \bar{Q}^{(r_1)} \times [0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ is continuous and satisfies $\mathbf{u}(\cdot + iy_0, t) \in U$ for every $y_0 \in \bar{Q}^{(r_1)}$ and for every $t \in [0, T]$, we construct \mathbf{u} first *locally in time* on a short time interval $[t_0, t_0 + T_1] \subset [0, T]$ as follows, where $T_1 > 0$ is small enough.

Let $\mathbf{u}^{(z_0)} : t \mapsto \mathbf{u}^{(z_0)}(\cdot, t) \equiv \mathbf{u}^{(z_0, t_0)}(\cdot, t)$ from $C([t_0, t_0 + T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ be as above, such that $z_0 = iy_0$ and $\mathbf{u}^{(iy_0)}(\cdot, t) \in U$ for every pair $(y_0, t) \in \bar{Q}^{(r_1)} \times [t_0, t_0 + T_1]$. Given any $y_0 \in \bar{Q}^{(r_1)}$, let us define

$$\mathbf{u}(x + iy_0, t) := \mathbf{u}^{(iy_0)}(x, t) \quad \text{for all } (x, t) \in \mathbb{R}^N \times [t_0, t_0 + T_1]. \quad (6.12)$$

Clearly, $\mathbf{u} : \bar{\mathfrak{X}}^{(r_1)} \times [t_0, t_0 + T_1] \rightarrow \mathbb{C}^M$ is a well-defined mapping, and it has the following properties.

Proposition 6.5. *Let $M, N \geq 1$, $0 < T < \infty$, and assume that (H1)–(H3) are satisfied with constants $0 < r_0 < \infty$, $0 < T_0 \leq T$, and $0 < \vartheta_0 < \pi/2$. Furthermore, assume that $\hat{\mathbf{u}} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ is a globally defined weak solution to the original Cauchy problem (1.1), i.e., (H8) is valid. Let the sets $U_0 \subset U \subset E_{1-\frac{1}{p}, p} = \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ be specified as in (6.2) and (6.3). Given any $t_0 \in [0, T]$, let $\mathbf{u}_0 \in U_0 \subset \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ be any initial data at time $t = t_0$, such that \mathbf{u}_0 has a holomorphic extension $\tilde{\mathbf{u}}_0 : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M$ as described in (H6) and Remark 6.2 (we identify $\mathbf{u}_0 \equiv \tilde{\mathbf{u}}_0$), with $\mathbf{u}_0^{(z_0)} = \mathbf{u}_0(\cdot + z_0) \in U_0 \subset E_{1-\frac{1}{p}, p}$ whenever $z_0 = x_0 + iy_0 \in \bar{Q}_\mathbb{C}^{(r_1)} \subset \bar{\mathfrak{X}}^{(r_1)}$ for some $r_1 \in (0, r)$. (We have set $r = r_0$; r_1 may depend on \mathbf{u}_0 .)*

Then there exists a number $T_1 \in (0, T - t_0]$, depending on r_1 and U , but not on t_0 , such that the Cauchy problem (6.1) on the (local) time interval $[t_0, t_0 + T_1] \subset [0, T]$ with the initial condition $\mathbf{u}(\cdot, t_0) = \mathbf{u}_0^{(z_0)}$ possesses a unique weak solution $\mathbf{u}^{(z_0)} \equiv \mathbf{u}^{(z_0, t_0)} \in C([t_0, t_0 + T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$, such that $\mathbf{u}^{(z_0)}(\cdot, t) \in U$ holds for every $t \in [t_0, t_0 + T_1]$. The family $\mathbf{u}^{(z_0)}$, parametrized by $z_0 = x_0 + iy_0 \in \bar{Q}_\mathbb{C}^{(r_1)}$, has the following properties:

- $\mathbf{u}^{(z_0)}(x, t) = \mathbf{u}^{(iy_0)}(x + x_0, t)$ holds for all $(x, t) \in \mathbb{R}^N \times [t_0, t_0 + T_1]$ and for all $z_0 \in \bar{Q}_\mathbb{C}^{(r_1)}$; consequently, even for all $z_0 \in \bar{\mathfrak{X}}^{(r_1)}$.
- The mapping $\mathbf{u} : \bar{\mathfrak{X}}^{(r_1)} \times [t_0, t_0 + T_1] \rightarrow \mathbb{C}^M$ defined in (6.12) satisfies $\mathbf{u}(x + x_0 + iy_0, t) = \mathbf{u}^{(x_0 + iy_0)}(x, t)$ for all $(x, t) \in \mathbb{R}^N \times [t_0, t_0 + T_1]$ and for all $z_0 = x_0 + iy_0 \in \mathbb{C}^N$ with $|y_0|_\infty \leq r_1$.
- The mapping $\mathbf{u} : \bar{\mathfrak{X}}^{(r_1)} \times [t_0, t_0 + T_1] \rightarrow \mathbb{C}^M : (z, t) = (x + iy, t) \mapsto \mathbf{u}(x + iy, t)$ is continuously (partially) differentiable with respect to all the real variables x_i and y_i ($i = 1, 2, \dots, N$) in $x = (x_1, \dots, x_N)$ and $y = (y_1, \dots, y_N)$ in \mathbb{R}^N with $|y|_\infty \leq r_1$.
- For each fixed $t \in [t_0, t_0 + T_1]$, the mapping $\mathbf{u}(\cdot, t) : \bar{\mathfrak{X}}^{(r_1)} \rightarrow \mathbb{C}^M : z = x + iy \mapsto \mathbf{u}(x + iy, t)$ is holomorphic, i.e., (partially) complex differentiable with respect to all the complex variables $z_i = x_i + iy_i$ ($i = 1, 2, \dots, N$) in $z = (z_1, \dots, z_N) \in \bar{\mathfrak{X}}^{(r_1)} \subset \mathbb{C}^N$.

As for Part (d) in this proposition, there are several equivalent definitions of a holomorphic function of several complex variables used in the literature, cf. Krantz [57, Definitions I–IV, pp. 3–4]. We adopt the most widely used definition in [57], Definition II (p. 3) and Definition 1.2.1 (p. 24). From this definition it is easy to derive the existence of an absolutely convergent power series ([57, Definition III, p. 3]) and the Cauchy formula in a polydisc ([57, Definition IV, pp. 3–4]). Nevertheless,

the equivalence of [57, Definition I, p. 3] verified in part (d) in our proposition above with [57, Definitions II–IV, pp. 3–4] is a deep classical result due to Hartogs; see [57, Theorem 1.2.5, p. 25]. However, taking also part (c) into account, we observe that also [57, Definition II, p. 3] is verified in our proposition.

Proof of Proposition 6.5. Recalling our remarks on the local (in time) existence and uniqueness before this proposition, we observe that it suffices to verify only our claims in Parts (a)–(d).

(a) Clearly, given any fixed $z_0 = x_0 + iy_0 \in \bar{Q}_{\mathbb{C}}^{(r_1)}$, both functions

$$t \mapsto \mathbf{u}^{(z_0)}(\cdot, t) \quad \text{and} \quad t \mapsto \mathbf{u}^{(iy_0)}(\cdot + x_0, t) : [t_0, t_0 + T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$$

are weak solutions to our Cauchy problem (6.1) on the (sufficiently short) time interval $[t_0, t_0 + T_1] \subset [0, T]$ with the same initial data $\mathbf{u}_0^{(z_0)}(\cdot) = \mathbf{u}_0^{(iy_0)}(\cdot + x_0)$ at time $t = t_0 \in [0, T]$, for some $T_1 \in (0, T - t_0]$. The uniqueness for problem (6.1) now forces $\mathbf{u}^{(z_0)}(\cdot, t) \equiv \mathbf{u}^{(iy_0)}(\cdot + x_0, t)$ for every $t \in [t_0, t_0 + T_1]$ as claimed.

Part (b) is an immediate consequence of Part (a) applied to (6.12).

(c) At the initial time $t = t_0$, the continuous (partial) differentiability is valid by our hypotheses on the initial data $\mathbf{u}_0 : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^N$ viewed as a function

$$z_0 = x_0 + iy_0 = (x_0, y_0) \mapsto \mathbf{u}_0(\cdot + z_0) = \mathbf{u}_0^{(z_0)} : \mathfrak{X}^{(r)} = \mathbb{R}^N \times Q^{(r)} \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$$

valued in the Besov space $E_{1-\frac{1}{p},p} = \mathbf{B}^{s;p,p}(\mathbb{R}^N)$; in particular, $\mathbf{u}_0(\cdot + z_0) = \mathbf{u}_0^{(z_0)} \in U_0 \subset E_{1-\frac{1}{p},p}$ provided $z_0 = x_0 + iy_0 \in \bar{Q}_{\mathbb{C}}^{(r_1)} \subset \tilde{\mathfrak{X}}^{(r_1)} (\subset \mathfrak{X}^{(r)})$. We recall that U_0 is an open subset of $E_{1-\frac{1}{p},p}$ defined in (6.2). We may view this C^1 differentiability as (partial) differentiability with respect to the real parameters $x_{0,i}$ and $y_{0,i}$ in the complex shift $z_0 = x_0 + iy_0 \in \bar{Q}_{\mathbb{C}}^{(r_1)}$, where $x_0 = (x_{0,1}, \dots, x_{0,N})$ and $y_0 = (y_{0,1}, \dots, y_{0,N})$ are in \mathbb{R}^N with $\max\{|x_0|_{\infty}, |y_0|_{\infty}\} \leq r_1$.

We now briefly interrupt our proof of Proposition 6.5 to make the following remarks:

Remarks. The kind of theory on continuous and differentiable dependence of the solution

$$z_0 \mapsto \mathbf{u}(\cdot + z_0, t) = \mathbf{u}^{(z_0)}(\cdot, t) : \tilde{\mathfrak{X}}^{(r_1)} = \mathbb{R}^N \times \bar{Q}^{(r_1)} \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$$

for $t \in [t_0, t_0 + T_1]$, on the real parameters $x_{0,i}$ and $y_{0,i}$ in $z_0 \in \tilde{\mathfrak{X}}^{(r_1)}$, that has been developed in Henry [40, Chapt. 3], §3.4, pp. 62–70, or, alternatively, in Lunardi [65, Chapt. 8], §8.3.1, pp. 302–306, can be adapted also to our setting for the spatially “shifted” Cauchy problem (6.1), with only minor changes. We should remark that, in this approach, the following hypotheses on A and f will do; they follow from (H1)–(H3) (cf. Lemma 6.3 and its proof):

6.3. Hypothesis. In analogy to, (H4’) and (H5’) let us assume that there are positive constants $\vartheta_0 \in (0, \pi/2)$ and $T_0 \in (0, T]$, and open sets $\mathcal{U} \subset \mathbb{C}$ and $\tilde{U} \subset E_{1-\frac{1}{p},p}$ containing the compact set $\bar{\Delta}_{\vartheta_0}^{T_0, T}$ and the open set U , respectively, i.e., $\bar{\Delta}_{\vartheta_0}^{T_0, T} \subset \mathcal{U} \subset \mathbb{C}$ and $U \subset \tilde{U} \subset E_{1-\frac{1}{p},p}$, such that

(H4’’) $A : [0, T] \times \mathcal{U} \rightarrow \mathcal{L}(E_1 \rightarrow E_0)$ possesses a continuously (Fréchet-) differentiable extension (i.e., of class C^1) $\tilde{A} : \mathcal{U} \times \tilde{U} \rightarrow \mathcal{L}(E_1 \rightarrow E_0)$ to the complex domain $\mathcal{U} \times \tilde{U}$ which satisfies $\tilde{A}(t, v) \in \text{MR}_p(E)$ for all $(t, v) \in \mathcal{U} \times \tilde{U}$.

(H5'') $f : [0, T] \times U \rightarrow E_0$ possesses a continuously (Fréchet-) differentiable extension $\tilde{f} : \mathcal{U} \times \tilde{U} \rightarrow E_0$ to the complex domain $\mathcal{U} \times \tilde{U}$.

Clearly, in both these hypotheses, the mappings A and f , respectively, are extended from the domain $[0, T] \times U \subset \bar{\Delta}_{\partial_0}^{T_0, T} \times U$ to the complex domain $\mathcal{U} \times \tilde{U} \subset \mathbb{C} \times E_{1-\frac{1}{p}, p}$.

Now we continue the proof of Proposition 6.5. Proof of part (c). Recall that the metric on $\mathcal{U} \times \tilde{U}$ ($\subset \mathbb{C} \times E_{1-\frac{1}{p}, p}$) is induced by the canonical norm on $\mathbb{C} \times E_{1-\frac{1}{p}, p}$. It is evident that (H4') and (H5') imply (H4'') and (H5''), respectively.

Applying the results from [40, Chapt. 3, §3.4] or [65, Chapt. 8, §8.3.1], we now conclude that the mapping

$$z_0 \mapsto \mathbf{u}(\cdot + z_0, t) : \tilde{\mathfrak{X}}^{(r_1)} = \mathbb{R}^N \times \bar{Q}^{(r_1)} \rightarrow \mathbf{B}^{s; p, p}(\mathbb{R}^N), \quad t \in [t_0, t_0 + T_1],$$

is continuously differentiable with respect to the real parameters $x_{0,i}$ and $y_{0,i}$ in $z_0 \in \tilde{\mathfrak{X}}^{(r_1)}$. The partial derivatives,

$$\frac{\partial \mathbf{u}}{\partial x_{0,i}} \equiv \frac{\partial \mathbf{u}}{\partial x_i} = \frac{\partial \mathbf{u}}{\partial z_i}, \quad \frac{\partial \mathbf{u}}{\partial y_{0,i}} \equiv \frac{\partial \mathbf{u}}{\partial y_i} = i \cdot \frac{\partial \mathbf{u}}{\partial z_i} : \tilde{\mathfrak{X}}^{(r_1)} \times [t_0, t_0 + T_1] \rightarrow \mathbb{C}^M$$

valued in $C([t_0, t_0 + T_1] \rightarrow \mathbf{B}^{s; p, p}(\mathbb{R}^N))$, are the unique weak solutions of the following Cauchy problems derived from (6.1) by the corresponding partial differentiation, respectively:

$$\begin{aligned} & \frac{\partial}{\partial t} \left(\frac{\partial \mathbf{u}}{\partial x_i} \right) + \frac{\partial \mathbf{P}}{\partial x_i} \left(x + z_0, t, \frac{1}{i} \frac{\partial}{\partial x} \right) \mathbf{u}(x + z_0, t) \\ & + \mathbf{P} \left(x + z_0, t, \frac{1}{i} \frac{\partial}{\partial x} \right) \left(\frac{\partial \mathbf{u}}{\partial x_i} \right) (x + z_0, t) \\ & = \frac{\partial \mathbf{f}}{\partial x_i} \left(x + z_0, t; \left(\frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta} (x + z_0, t) \right)_{|\beta| \leq m} \right) \\ & + \sum_{|\beta| \leq m} \sum_{k=1}^M \frac{\partial \mathbf{f}}{\partial Z_{\beta, k}} \left(x + z_0, t; \left(\frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta} \right)_{|\beta| \leq m} \right) \frac{\partial^{|\beta|}}{\partial x^\beta} \left(\frac{\partial u_k}{\partial x_i} \right) (x + z_0, t) \end{aligned} \tag{6.13}$$

for $(x, t) \in \mathbb{R}^N \times (t_0, t_0 + T_1)$;

$$\frac{\partial \mathbf{u}}{\partial x_i} (x + z_0, 0) = \frac{\partial \mathbf{u}_0}{\partial x_i} (x + z_0) \quad \text{for } x \in \mathbb{R}^N,$$

and

$$\begin{aligned} & \frac{\partial}{\partial t} \left(\frac{\partial \mathbf{u}}{\partial y_i} \right) + \frac{\partial \mathbf{P}}{\partial y_i} \left(x + z_0, t, \frac{1}{i} \frac{\partial}{\partial x} \right) \mathbf{u}(x + z_0, t) \\ & + \mathbf{P} \left(x + z_0, t, \frac{1}{i} \frac{\partial}{\partial x} \right) \left(\frac{\partial \mathbf{u}}{\partial y_i} \right) (x + z_0, t) \\ & = \frac{\partial \mathbf{f}}{\partial y_i} \left(x + z_0, t; \left(\frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta} (x + z_0, t) \right)_{|\beta| \leq m} \right) \\ & + \sum_{|\beta| \leq m} \sum_{k=1}^M \frac{\partial \mathbf{f}}{\partial Z_{\beta, k}} \left(x + z_0, t; \left(\frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta} \right)_{|\beta| \leq m} \right) \frac{\partial^{|\beta|}}{\partial x^\beta} \left(\frac{\partial u_k}{\partial y_i} \right) (x + z_0, t) \end{aligned} \tag{6.14}$$

for $(x, t) \in \mathbb{R}^N \times (t_0, t_0 + T_1)$;

$$\frac{\partial \mathbf{u}}{\partial y_i} (x + z_0, 0) = \frac{\partial \mathbf{u}_0}{\partial y_i} (x + z_0) \quad \text{for } x \in \mathbb{R}^N,$$

where the complex variable $Z_{\beta,k}$ stands for $Z_{\beta,k} = \frac{\partial^{|\beta|}}{\partial x^\beta} u_k = i^{|\beta|} D_x^\beta u_k \in \mathbb{C}$. This proves Part (c).

Proof of part (d). We take advantage of the two equations (6.13) and (6.14), to apply the Cauchy-Riemann operator $\partial/\partial \bar{z}_i$ from (2.2) to problem (6.1) to conclude that the Cauchy-Riemann derivative

$$\bar{\partial}_{z_0,i} \mathbf{u} \equiv \frac{\partial \mathbf{u}}{\partial \bar{z}_{0,i}} \equiv \frac{\partial \mathbf{u}}{\partial \bar{z}_i} = \frac{1}{2} \left(\frac{\partial}{\partial x_i} + i \frac{\partial}{\partial y_i} \right) \mathbf{u} : \bar{\mathfrak{X}}^{(r_1)} \times [t_0, t_0 + T_1] \rightarrow \mathbb{C}^M$$

is in $C([t_0, t_0 + T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ and obeys the following homogeneous linear Cauchy problem, which is a simple linear combination $\frac{1}{2} \cdot (6.13) + \frac{i}{2} \cdot (6.14) = (6.15)$:

$$\begin{aligned} & \frac{\partial}{\partial t} (\bar{\partial}_{z_0,i} \mathbf{u}) + \mathbf{P} \left(x + z_0, t, \frac{1}{i} \frac{\partial}{\partial x} \right) (\bar{\partial}_{z_0,i} \mathbf{u})(x + z_0, t) \\ &= \sum_{|\beta| \leq m} \sum_{k=1}^M \frac{\partial \mathbf{f}}{\partial Z_{\beta,k}} \left(x + z_0, t; \left(\frac{\partial^{|\beta|} \mathbf{u}}{\partial x^\beta} \right)_{|\beta| \leq m} \right) \frac{\partial^{|\beta|}}{\partial x^\beta} (\bar{\partial}_{z_0,i} u_k)(x + z_0, t) \end{aligned} \quad (6.15)$$

for $(x, t) \in \mathbb{R}^N \times (t_0, t_0 + T_1)$;

$$(\bar{\partial}_{z_0,i} \mathbf{u})(x + z_0, 0) = \mathbf{0} \quad \text{for } x \in \mathbb{R}^N.$$

Here, we have used that both operators

$$z \mapsto \mathbf{P} \left(z, t, \frac{1}{i} \frac{\partial}{\partial x} \right) \quad \text{and} \quad z \mapsto \mathbf{f}(z, t; (Z_\beta)_{|\beta| \leq m}) : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M \quad (r = r_0)$$

are holomorphic, i.e., $\bar{\partial}_{z_i} \mathbf{P}(z, t, \frac{1}{i} \frac{\partial}{\partial x}) = \mathbf{0}$ and $\bar{\partial}_{z_i} \mathbf{f}(z, t; (Z_\beta)_{|\beta| \leq m}) = \mathbf{0}$, by (H1) and (H3), respectively. By our choice of $\mathbf{u}_0 \equiv \bar{\mathbf{u}}_0 : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M$ being holomorphic, we have also $\bar{\partial}_{z_i} \mathbf{u}_0(z) = \mathbf{0}$; $i = 1, 2, \dots, N$. Notice that (6.15) is valid only for every $z_0 \in \bar{\mathfrak{X}}^{(r_1)} (\subset \mathfrak{X}^{(r)})$.

By (H1) and (H2), the linear differential operator on the left-hand side of (6.15), i.e.,

$$\frac{\partial}{\partial t} + \mathbf{P} \left(x + z_0, t, \frac{1}{i} \frac{\partial}{\partial x} \right),$$

is uniformly parabolic of order $2m$ with smooth coefficients. It is proved in Denk, Hieber and Prüss [23, p. 67], Theorem 5.7 (cf. also [74], Theorem 2.1 (p. 8) and remarks thereafter (p. 9)) that, for every $z_0 \in \bar{\mathfrak{X}}^{(r)}$ and for every $t \in [0, T]$,

$$A^{(z_0)}(t) := -\mathbf{P} \left(x + z_0, t, \frac{1}{i} \frac{\partial}{\partial x} \right) : \mathbf{W}^{2m,p}(\mathbb{R}^N) \rightarrow \mathbf{L}^p(\mathbb{R}^N) \quad (6.16)$$

is a bounded linear operator, i.e., $A^{(z_0)}(t) \in \mathcal{L}(E_1 \rightarrow E_0)$, and it possesses the maximal L^p -regularity property, i.e., $A^{(z_0)}(t) \in \text{MR}_p(E) \equiv \text{MR}_p(E_1 \rightarrow E_0)$. Let us recall that $E_1 = \mathbf{W}^{2m,p}(\mathbb{R}^N) \hookrightarrow E_0 = \mathbf{L}^p(\mathbb{R}^N)$.

Furthermore, in view of (H3), the pointwise multiplication and differentiation operators on the right-hand side of (6.15) are of order $|\beta|$ ($|\beta| \leq m < 2m$) and all have bounded continuous coefficients, by $\mathbf{u}(\cdot, t) \in U \subset \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ for every $t \in [t_0, t_0 + T_1]$ combined with the Sobolev imbedding $B^{s;p,p}(\mathbb{R}^N) \hookrightarrow C^m(\mathbb{R}^N) \cap W^{m,\infty}(\mathbb{R}^N)$, where $2 + \frac{N}{m} < p < \infty$ and $m < s = 2m(1 - \frac{1}{p}) < 2m$. We denote their sum, which appears in (6.15), by $f^{(z_0)}(t) : \mathbf{B}^{s;p,p}(\mathbb{R}^N) \rightarrow \mathbf{L}^p(\mathbb{R}^N)$, i.e., $f^{(z_0)}(t) \in \mathcal{L}(E_{1-\frac{1}{p},p} \rightarrow E_0)$. Here, we allow any $z_0 \in \bar{\mathfrak{X}}^{(r_1)}$. Consequently, the mappings $(t, v) \mapsto A^{(z_0)}(t) : [0, T] \times U \rightarrow \mathcal{L}(E_1 \rightarrow E_0)$ and $(t, v) \mapsto f^{(z_0)}(t)v : [0, T] \times U \rightarrow E_0$

satisfy (H4) and (H5) for A and f , respectively, with $U = E_{1-\frac{1}{p},p}$, $A^{(z_0)}(t) \in \mathcal{L}(E_1 \rightarrow E_0)$ being independent from $v \in U$, and $v \mapsto f^{(z_0)}(t)v$ linear in $v \in U$.

We observe that the homogeneous linear Cauchy problem (6.15) for the Cauchy-Riemann derivative $\bar{\partial}_{z_0,i} \mathbf{u}$ of \mathbf{u} takes the following abstract linear form, whenever $z_0 \in \tilde{\mathfrak{X}}^{(r_1)}$:

$$\begin{aligned} \frac{d}{dt}(\bar{\partial}_{z_0,i} \mathbf{u}) - A^{(z_0)}(t)(\bar{\partial}_{z_0,i} \mathbf{u}) &= f^{(z_0)}(t)(\bar{\partial}_{z_0,i} \mathbf{u}) \quad \text{for a.e. } t \in (t_0, t_0 + T_1); \\ (\bar{\partial}_{z_0,i} \mathbf{u})(0) &= \mathbf{0} \in E_{1-\frac{1}{p},p}. \end{aligned} \tag{6.17}$$

This abstract linear problem corresponds to the nonlinear initial value problem (4.8) treated in Section 4. We apply the uniqueness part of Theorem 4.5 to deduce $(\bar{\partial}_{z_0,i} \mathbf{u})(x, t) \equiv \mathbf{0}$ for all $(x, t) \in \mathbb{R}^N \times [t_0, t_0 + T_1]$. This implies that the mapping $z \mapsto u_k(z, t) : \tilde{\mathfrak{X}}^{(r_1)} \rightarrow \mathbb{C}$ is holomorphic in each complex variable $z_i \in \mathbb{C}$, for every fixed time $t \in [t_0, t_0 + T_1]$; $k = 1, 2, \dots, N$. Moreover, by part (c), all complex partial derivatives $\partial_{z_i} u_k(\cdot, t)$ are continuous in $\tilde{\mathfrak{X}}^{(r_1)}$. Finally, we take advantage of the classical fact that such a function $u_k(\cdot, t) : \tilde{\mathfrak{X}}^{(r_1)} \rightarrow \mathbb{C}$ is holomorphic (Remark 2.1); see e.g. John [50, Theorem, p. 70] or Krantz [57, Definition II, p. 3]. Also Part (d) and, thus, the entire proposition is proved. \square

7. SPACE-TIME ANALYTICITY FOR THE CAUCHY PROBLEM IN $\mathbb{R}^N \times (0, T)$

We summarize the time and space analyticity results from the last two sections (Sections 5 and 6), for the mapping $\mathbf{u} : \tilde{\mathfrak{X}}^{(r_1)} \times [t_0, t_0 + T_1] \rightarrow \mathbb{C}^M$ defined in (6.12), in the following theorem.

Theorem 7.1. *Let $M, N \geq 1$, $0 < T < \infty$, and assume that (H1)–(H3) are satisfied with some constants $0 < r_0 < \infty$, $0 < T_0 \leq T$, and $0 < \vartheta_0 < \pi/2$. Furthermore, assume that $\hat{\mathbf{u}} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ is a globally defined weak solution to the original Cauchy problem (1.1), i.e., (H8) is valid. Let the sets $U_0 \subset U \subset E_{1-\frac{1}{p},p} = \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ be specified as in (6.2) and (6.3). Given any $t_0 \in [0, T)$, let $\mathbf{u}_0 \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ be any initial data at time $t = t_0$, such that $\mathbf{u}_0 \in U_0$ and \mathbf{u}_0 has a holomorphic extension $\tilde{\mathbf{u}}_0 : \tilde{\mathfrak{X}}^{(r)} \rightarrow \mathbb{C}^M$ as described before Lemma 6.3 (we identify $\mathbf{u}_0 \equiv \tilde{\mathbf{u}}_0$), with $\mathbf{u}_0^{(z_0)} = \mathbf{u}_0(\cdot + z_0) \in U_0 \subset E_{1-\frac{1}{p},p}$ whenever $z_0 \in \bar{Q}_{\mathbb{C}}^{(r_1)} \subset \tilde{\mathfrak{X}}^{(r_1)}$ for some $r_1 \in (0, r)$, cf. (H6). (We have set $r = r_0$; r_1 may depend on \mathbf{u}_0 .)*

Finally, let $\mathbf{u} : \tilde{\mathfrak{X}}^{(r_1)} \times [t_0, t_0 + T_1] \rightarrow \mathbb{C}^M$ be the continuous mapping obtained in Proposition 6.5, with $T_1 \in (0, T - t_0]$ depending on r_1 and U , but not on t_0 . Replace $T_0 \in (0, T]$ by $\min\{T_0, T_1\}$ if necessary, so that $0 < T_0 \leq T_1 \leq T$ holds. Then there exist constants $\vartheta' \in (0, \vartheta_0]$ and $T' \in (0, T_0]$, small enough, and a continuous mapping $\tilde{\mathbf{u}} : \tilde{\mathfrak{X}}^{(r_1)} \times (t_0 + \Delta_{\vartheta', T'}^{T', T_1}) \rightarrow \mathbb{C}^M$ with the following properties:

- (i) *For each $z_0 \in \tilde{\mathfrak{X}}^{(r_1)}$, the (unique) weak solution*

$$\mathbf{u}^{(z_0)} \in C([t_0, t_0 + T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$$

to the Cauchy problem (6.1) on the time interval $[t_0, t_0 + T_1] \subset [0, T]$ with the initial condition $\mathbf{u}(\cdot, t_0) = \mathbf{u}_0^{(z_0)}$ at time $t = t_0$ satisfying $z_0 \in \bar{Q}_{\mathbb{C}}^{(r_1)}$ possesses a unique holomorphic extension from $(t_0, t_0 + T_1)$ to $t_0 + \Delta_{\vartheta', T'}^{T', T_1}$, such that $\mathbf{u}^{(z_0)}(\cdot, t_0 + s) = \tilde{\mathbf{u}}(\cdot + z_0, t_0 + s) \in U$ holds for every $s \in \Delta_{\vartheta', T'}^{T', T_1}$.

- (ii) The complex function $\tilde{\mathbf{u}} : \mathfrak{X}^{(r_1)} \times (t_0 + \Delta_{\vartheta'}^{T', T_1}) \rightarrow \mathbb{C}^M$ is holomorphic (jointly) in all its variables, $z = (z_1, z_2, \dots, z_N) \in \mathfrak{X}^{(r_1)} \subset \mathbb{C}^N$ and $t \in t_0 + \Delta_{\vartheta'}^{T', T_1} \subset \mathbb{C}$.

Let us recall that, by our notation in (1.7), for $\vartheta \in (0, \pi/2)$, $0 < t_0 < T < \infty$, and $0 < T' \leq T - t_0$, we have

$$\begin{aligned} t_0 + \Delta_{\vartheta}^{T', T-t_0} &= (t_0 + \Delta_{\vartheta}^{(T-t_0)}) \cap \{t \in \mathbb{C} : |\Im t| < T' \tan \vartheta\} \\ &= \cup_{t_0 \leq \xi \leq T-T'} \{\xi + t' \in \mathbb{C} : t' \in \Delta_{\vartheta}^{(T')}\} \\ &= \cup_{t_0 \leq \xi \leq T-T'} (\xi + \Delta_{\vartheta}^{(T')}) \end{aligned} \quad (7.1)$$

with the closure $t_0 + \bar{\Delta}_{\vartheta}^{T', T-t_0}$ in \mathbb{C} .

Remark 7.2. (a) The main difference between our main result, Theorem 3.4 (Section 3), and Theorem 7.1 above is the *temporally local* character of the latter stated for the time interval $[t_0, t_0 + T_1] \subset [0, T]$ with the additional analyticity hypothesis on the initial data \mathbf{u}_0 (as in part (iii) of Theorem 3.4).

(b) Recalling our choice of the number $r \in (0, \infty)$ in (H6) (before Remark 6.1) on the complex analyticity of the initial data $\mathbf{u}_0 = \tilde{\mathbf{u}}_0 : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M$ extended to the complex strip $\mathfrak{X}^{(r)} \subset \mathbb{C}^N$, we observe that the number $r_1 \in (0, r)$, originally introduced in the spatially “shifted” Cauchy problem (6.1), is needed for sufficiently small perturbations (“shifts”) $z_0 \in \mathbb{C}^N$ of the space variable $z \in \mathfrak{X}^{(r_1)}$ in order to keep $z + z_0 \in \mathfrak{X}^{(r)}$. As we have already mentioned after Remark 6.1, (H1)–(H3) show that only the case $0 < r \leq r_0$ is useful. We now recall from Proposition 6.5 and Theorem 7.1 that, to avoid excessive notation, $r_1 \in (0, r)$ must be chosen small enough, such that $\mathbf{u}_0^{(z_0)} = \mathbf{u}_0(\cdot + z_0) \in U_0 \subset U \subset E_{1-\frac{1}{p}, p}$ for every $z_0 \in \bar{Q}_{\mathbb{C}}^{(r_1)} \subset \bar{\mathfrak{X}}^{(r_1)}$. We recall that the sets U_0 and U are defined in (6.2) and (6.3), respectively. We stress that both, U_0 and U , are open in $E_{1-\frac{1}{p}, p}$, while being determined solely by the restriction to the real line \mathbb{R} , $\mathbf{u}_0 = \tilde{\mathbf{u}}_0|_{\mathbb{R}} \in U_0 \subset E_{1-\frac{1}{p}, p}$, of the initial data $\tilde{\mathbf{u}}_0 : \mathfrak{X}^{(r)} \rightarrow \mathbb{C}^M$, i.e., by $\mathbf{u}_0 : \mathbb{R}^N \rightarrow \mathbb{C}^M$ as an element of the Besov space $E_{1-\frac{1}{p}, p} = \mathbf{B}^{s; p, p}(\mathbb{R}^N)$. Consequently, the number $r_1 \equiv r_1(\mathbf{u}_0) \in (0, \infty)$ is determined by these initial data $\mathbf{u}_0 \in U_0$; we have $0 < r_1 < r \leq r_0$ where we may choose $r = r_0$, by Remark 6.1. Such a choice of $r_1 \in (0, r)$ is possible thanks to the closed cube $\bar{Q}^{(r_1)}$ being compact in \mathbb{R}^N .

Proof of Theorem 7.1. (i) Let $z_0 = x_0 + iy_0 \in \bar{\mathfrak{X}}^{(r_1)}$ be arbitrary, but fixed, with $\max\{|x_0|_{\infty}, |y_0|_{\infty}\} \leq r_1$, i.e., $z_0 \in \bar{Q}_{\mathbb{C}}^{(r_1)}$. We apply our time analyticity result in Theorem 5.3 to the Cauchy problem (6.1) on the time interval $[t_0, t_0 + T_1] \subset [0, T]$ with the initial condition $\mathbf{u}(\cdot, t_0) = \mathbf{u}_0^{(z_0)} \in U_0$ at time $t = t_0$ to derive the conclusion of Part (i).

(ii) The second part is obtained by combining Part (i) with Proposition 6.5, particularly Part (d). Finally, the joint time and space analyticity of the complex function $\tilde{\mathbf{u}} : \mathfrak{X}^{(r_1)} \times (t_0 + \Delta_{\vartheta'}^{T', T_1}) \rightarrow \mathbb{C}^M$ is obtained by applying the classical characterization of holomorphic functions by the Cauchy-Riemann equations (Remark 2.1); see e.g. John [50, Theorem, p. 70] or Krantz [57, Definition II, p. 3]. \square

8. PROOFS OF MAIN RESULTS

Now we are ready to prove Theorem 3.4. Then Proposition 3.5 is a consequence of Theorem 3.4 and Lemma 5.2, inequality (5.6). We give its proof right after that of Theorem 3.4.

Proof of Theorem 3.4. (i) The existence and uniqueness of a weak solution $\mathbf{u} \in C([0, T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ to the Cauchy problem (1.1), local in time for $t \in [0, T_1]$ with some $T_1 \in (0, T]$, is obtained directly from the abstract result in Theorem 4.5 where $E_{1-\frac{1}{p},p} = \mathbf{B}^{s;p,p}(\mathbb{R}^N) = [B^{s;p,p}(\mathbb{R}^N)]^M$. The technical details in applying Theorem 4.5 (an abstract result) to problem (1.1) have been given in Section 6, right after problem (6.1). The linear parabolic operator on the left-hand side in (6.1) is treated by the maximal L^p -regularity described in Remark 4.2(a). The special case of p in this remark, $p_0 = 2$, is taken care of by standard parabolic regularity making use of Gårding's inequality in Corollary 3.3; see, e.g., Friedman [31, Chapt. 10]. If a weak solution $\mathbf{u} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ exists globally in time $t \in [0, T]$, then it is unique, by Theorem 4.7, Part (i).

(ii) The (unique) temporal extension of the function $\mathbf{u} : \mathbb{R}^N \times (0, T) \rightarrow \mathbb{C}^M$ to a holomorphic function $\mathbf{u}^\sharp : \Delta_{\vartheta'}^{T',T} \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ that possesses another extension to a continuous function on the closure $\bar{\Delta}_{\vartheta'}^{T',T}$, denoted again by \mathbf{u}^\sharp , is derived from Theorem 7.1, Part (i). More precisely, Part (i) of Theorem 7.1 is applied to the (global) weak solution $\mathbf{u} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ of the Cauchy problem (1.1), which is assumed to exist, in the temporal complex domain $t_0 + \Delta_{\vartheta'}^{(T')}$ with the initial value $\mathbf{u}(\cdot, t_0) \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ at every initial time $t_0 \in [0, T - T']$. Here, we have used that

$$\Delta_{\vartheta'}^{T',T} = \cup_{0 \leq t_0 \leq T - T'} (t_0 + \Delta_{\vartheta'}^{(T')}) \quad (8.1)$$

(cf. (7.1)).

(iii) We remark that (H8) is satisfied with the function

$$\hat{\mathbf{u}} = \mathbf{u} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)),$$

a globally defined weak solution to the original Cauchy problem (1.1), which exists by our hypothesis. Let us recall the definitions of the bounded, open, and convex sets U_0 and U in $\mathbf{B}^{s;p,p}(\mathbb{R}^N)$, $U_0 \subset U = \tilde{U} \subset E_{1-\frac{1}{p},p} = \mathbf{B}^{s;p,p}(\mathbb{R}^N)$, in (6.2) and (6.3), respectively, where the radius $R_0 \in (0, \infty)$ is an arbitrary positive number. We recall also our hypothesis that the initial condition $\mathbf{u}(\cdot, 0) = \mathbf{u}_0 \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ possesses a (unique) holomorphic extension $\tilde{\mathbf{u}}_0 : \mathfrak{X}^{(\kappa_0)} \rightarrow \mathbb{C}^M$ from \mathbb{R}^N to the complex domain $\mathfrak{X}^{(\kappa_0)} = \mathbb{R}^N + iQ^{(\kappa_0)} \subset \mathbb{C}^N$ (a tube), for some $\kappa_0 \in (0, r_0]$, that satisfies (3.10).

We begin with a construction of the (unique) spatial extension of the continuous function $\mathbf{u} : \mathbb{R}^N \times [0, T] \rightarrow \mathbb{C}^M$ to a continuous function $\mathbf{u}^\flat : \tilde{\mathfrak{X}}^{(\rho)} \times [0, T] \rightarrow \mathbb{C}^M$ that is holomorphic in the space variable $z = x + iy \in \tilde{\mathfrak{X}}^{(\rho)} = \mathbb{R}^N + iQ^{(\rho)}$ with some $\rho \in (0, \kappa_0]$ small enough. Let us recall our notation with the “shifted” function $x \mapsto \mathbf{u}_0^{(z_0)}(x) := \mathbf{u}_0(x + z_0) : \mathbb{R}^N \rightarrow \mathbb{C}^M$ introduced in the Cauchy problem (6.1) spatially “shifted” by $z_0 = x_0 + iy_0 \in \mathfrak{X}^{(r)} \subset \mathbb{C}^N$. The constant $r \in (0, \infty)$ has been introduced in (H6); only the case $0 < r < \kappa_0 \leq r_0 (< \infty)$ is useful. We wish to apply Proposition 6.5 with the constant $r_1 \in (0, r)$ specified there. We choose

$\rho \in (0, r_1)$ small enough, such that also

$$\|\mathbf{u}_0^{(iy)} - \mathbf{u}_0\|_{B^{s;p,p}(\mathbb{R}^N)} = \|\mathbf{u}_0(\cdot + iy) - \mathbf{u}_0\|_{B^{s;p,p}(\mathbb{R}^N)} < R_0 \quad \text{holds for all } y \in \bar{Q}^{(\rho)}.$$

Here, we have used the (Lipschitz) continuity of the mapping $y \mapsto \mathbf{u}_0^{(iy)} := \mathbf{u}_0(\cdot + iy) : \bar{Q}^{(r_1)} \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$, thanks to $0 < r_1 < \kappa_0$ supplemented by the Cauchy formula in a polydisc centered in $\tilde{\mathfrak{X}}^{(r_1)}$ (with radius $< \kappa_0 - r_1$) and contained in the complex strip $\mathfrak{X}^{(\kappa_0)} \subset \mathbb{C}^N$. From (6.2) we deduce that $\mathbf{u}_0^{(iy)} \in U_0$ for all $y \in \bar{Q}^{(\rho)}$. In analogy with our proof of Part (i) above, we apply Theorem 4.5 to conclude that the spatially shifted Cauchy problem (6.1), with the shift $z_0 = iy$ ($y \in \bar{Q}^{(\rho)}$), possesses a unique weak solution $\mathbf{u}^{(iy)} \in C([0, T_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$, local in time for $t \in [0, T_1]$ with some $T_1 \in (0, T]$, that satisfies $\mathbf{u}^{(iy)}(\cdot, t) \in U$ for every $t \in [0, T_1]$. We apply part (c) or part (d) of Proposition 6.5 with $t_0 = 0$ to conclude that there is a number $R_1 \in (0, R_0)$ small enough, such that even $\mathbf{u}^{(iy)}(\cdot, t) \in U_0 \subset U$ holds for every $t \in [0, T_1]$, provided $\rho \in (0, r_1)$ is chosen so small that also

$$\|\mathbf{u}_0^{(iy)} - \mathbf{u}_0\|_{B^{s;p,p}(\mathbb{R}^N)} < R_1 \quad (< R_0) \quad \text{holds for all } y \in \bar{Q}^{(\rho)}.$$

Here, besides the (Lipschitz) continuity of the mapping $y \mapsto \mathbf{u}_0^{(iy)} := \mathbf{u}_0(\cdot + iy) : \bar{Q}^{(r_1)} \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ mentioned above, we have used also the continuous dependence of the solution $\mathbf{u}^{(iy)}$ upon the initial data $\mathbf{u}_0^{(iy)} \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ obtained in Theorem 4.7, part (ii); see also Remark 4.6. According to (6.12), we define the function $\mathbf{u}^b : \tilde{\mathfrak{X}}^{(\rho)} \times [0, T_1] \rightarrow \mathbb{C}^M$ by the formula

$$\mathbf{u}^b(x + iy, t) := \mathbf{u}^{(iy)}(x, t) \quad \text{for all } (x, y, t) \in \mathbb{R}^N \times \bar{Q}^{(\rho)} \times [0, T_1]. \quad (8.2)$$

Clearly, by Proposition 6.5, part (c), the function $\mathbf{u}^b : \tilde{\mathfrak{X}}^{(\rho)} \times [0, T_1] \rightarrow \mathbb{C}^M$ is continuous and, by Proposition 6.5, part (d), also holomorphic with respect to the complex variable $z = x + iy \in \tilde{\mathfrak{X}}^{(\rho)} = \mathbb{R}^N + iQ^{(\rho)}$ at every time $t \in [0, T_1]$.

Next, we set $\mathbf{u}_1^{(iy)} := \mathbf{u}^{(iy)}(\cdot, T_1) \in U_0$ and repeat the procedure from above on the interval $[T_1, 2T_1]$ with the initial data $\mathbf{u}_1^{(iy)} \in U_0$ at $t_0 = T_1$ in place of $\mathbf{u}_0^{(iy)} \in U_0$ at $t_0 = 0$. We stress that the interval length $T_1 \in (0, T - t_0]$ in Proposition 6.5 is independent from the choice of the initial time $t_0 \in (0, T)$ whenever $[t_0, t_0 + T_1] \subset [0, T]$. Again, we apply part (c) or part (d) of Proposition 6.5 with $t_0 = T_1$ (in place of $t_0 = 0$) to conclude that there is a number $R_2 \in (0, R_1)$ small enough, such that even $\mathbf{u}^{(iy)}(\cdot, t) \in U_0 \subset U$ holds for every $t \in [0, 2T_1]$, provided $\rho \in (0, r_1)$ is chosen so small that also

$$\|\mathbf{u}_0^{(iy)} - \mathbf{u}_0\|_{B^{s;p,p}(\mathbb{R}^N)} < R_2 \quad (< R_1 < R_0) \quad \text{holds for all } y \in \bar{Q}^{(\rho)}.$$

The desired function \mathbf{u}^b is naturally extended from the domain $\tilde{\mathfrak{X}}^{(\rho)} \times [0, T_1]$ to $\tilde{\mathfrak{X}}^{(\rho)} \times [0, 2T_1]$ by setting (cf. (8.2))

$$\mathbf{u}^b(x + iy, t) := \mathbf{u}^{(iy)}(x, t) \quad \text{for all } (x, y, t) \in \mathbb{R}^N \times \bar{Q}^{(\rho)} \times [T_1, 2T_1]. \quad (8.3)$$

We keep repeating this procedure (by “induction” on k) with the initial data $\mathbf{u}_k^{(iy)} := \mathbf{u}^{(iy)}(\cdot, kT_1) \in U_0$ successively for every $k = 0, 1, 2, \dots, m$ until reaching the inequalities

$$(m-1)T_1 \leq T < mT_1 \quad \text{at } k = m-1.$$

In fact, setting $\mathbf{u}_{m-1}^{(iy)} := \mathbf{u}^{(iy)}(\cdot, (m-1)T_1) \in U_0$ and repeating the procedure from above on the time interval $[(m-1)T_1, mT_1]$ with the initial data $\mathbf{u}_{m-1}^{(iy)} \in U_0$ at $t_0 = (m-1)T_1$ in place of $\mathbf{u}_0^{(iy)} \in U_0$ at $t_0 = 0$, we can apply Theorem 4.5 to conclude

that the spatially shifted Cauchy problem (6.1), with the shift $z_0 = iy$ ($y \in \bar{Q}^{(\rho)}$), possesses a unique weak solution $\mathbf{u}^{(iy)} \in C([(m-1)T_1, mT_1] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$, local in time for $t \in [(m-1)T_1, mT_1]$, that satisfies $\mathbf{u}^{(iy)}(\cdot, t) \in U$ for every $t \in [(m-1)T_1, mT_1]$. Consequently, we may assume $T = mT_1$ instead of $(m-1)T_1 \leq T < mT_1$. In this way we have constructed a finite set of numbers $R_1, R_2, \dots, R_{m-1}, R_m$ such that $0 < R_m < R_{m-1} < \dots < R_1 < R_0$ and, provided $\rho \in (0, r_1)$ is chosen small enough, also

$$\|\mathbf{u}_k^{(iy)} - \mathbf{u}_0\|_{B^{s;p,p}(\mathbb{R}^N)} < R_k \quad \text{holds for all } y \in \bar{Q}^{(\rho)}, k = 0, 1, 2, \dots, m-1,$$

together with $\mathbf{u}^{(iy)}(\cdot, t) \in U_0 \subset U$ for every $t \in [0, (m-1)T_1]$ and $\mathbf{u}^{(iy)}(\cdot, t) \in U$ for every $t \in [(m-1)T_1, mT_1]$. Finally, the desired function \mathbf{u}^b is defined successively on the domains $\tilde{\mathfrak{X}}^{(\rho)} \times [(k-1)T_1, kT_1]$ for each $k = 1, 2, 3, \dots, m$ by the formula

$$\mathbf{u}^b(x + iy, t) := \mathbf{u}^{(iy)}(x, t) \quad \text{for all } (x, y, t) \in \mathbb{R}^N \times \bar{Q}^{(\rho)} \times [0, T]. \tag{8.4}$$

To summarize the result of the procedure described above, we have determined a constant $\rho \in (0, r_1)$, small enough, such that for each shift $y \in \bar{Q}^{(\rho)}$ there is a unique weak solution $\mathbf{u}^{(iy)} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ to the spatially shifted Cauchy problem (6.1) with the shift $z_0 = iy$, that satisfies $\mathbf{u}^{(iy)}(\cdot, t) \in U_0 \subset U$ for every $t \in [0, (m-1)T_1]$ and $\mathbf{u}^{(iy)}(\cdot, t) \in U$ for every $t \in [0, T]$, where $T = mT_1$. We apply Proposition 6.5, Parts (c) and (d), once again to conclude that the function $\mathbf{u}^b : \tilde{\mathfrak{X}}^{(\rho)} \times [0, T] \rightarrow \mathbb{C}^M$ constructed above in (8.4) has the desired properties: it is continuous and holomorphic in the space variable $z = x + iy \in \tilde{\mathfrak{X}}^{(\rho)} = \mathbb{R}^N + iQ^{(\rho)}$ with some $\rho \in (0, r_1)$ small enough, where $0 < r_1 < \kappa_0 \leq r_0$.

Now we are ready to finish our proof of part (iii) by further extending the (global) weak solution $\mathbf{u} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ of the Cauchy problem (1.1) from the domain $\tilde{\mathfrak{X}}^{(\rho)} \times [0, T]$ of the (unique) spatial extension $\mathbf{u}^b : \tilde{\mathfrak{X}}^{(\rho)} \times [0, T] \rightarrow \mathbb{C}^M$ to another continuous function $\tilde{\mathbf{u}} : \tilde{\mathfrak{X}}^{(\rho)} \times \bar{\Delta}_{\vartheta'}^{T',T} \rightarrow \mathbb{C}^M$ which is holomorphic in $\tilde{\mathfrak{X}}^{(\rho)} \times \Delta_{\vartheta'}^{T',T}$. We recall that the solution $\mathbf{u} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ is assumed to exist by hypothesis (in part (iii)) with the initial data $\mathbf{u}_0 \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ having a (unique) holomorphic extension $\tilde{\mathbf{u}}_0 : \tilde{\mathfrak{X}}^{(\kappa_0)} \rightarrow \mathbb{C}^M$ from \mathbb{R}^N to the complex domain $\tilde{\mathfrak{X}}^{(\kappa_0)} \subset \mathbb{C}^N$, for some $\kappa_0 \in (0, r_0]$.

We apply Theorem 7.1 to the function $\mathbf{u}^b : \tilde{\mathfrak{X}}^{(\rho)} \times [0, T] \rightarrow \mathbb{C}^M$ on every time interval $[t_0, t_0 + T_1] \subset [0, T]$. We remark that the number $T_1 \in (0, T - t_0]$ depends on r_1 and U , but not on $t_0 \in [0, T)$, provided $[t_0, t_0 + T_1] \subset [0, T]$. In fact, making use of the same argument as above, where we have extended the function \mathbf{u}^b from the domain $\tilde{\mathfrak{X}}^{(\rho)} \times [0, T]$ to $\tilde{\mathfrak{X}}^{(\rho)} \times [0, mT_1]$ in case $(m-1)T_1 \leq T < mT_1$, we can extend \mathbf{u}^b from the domain $\tilde{\mathfrak{X}}^{(\rho)} \times [t_0, T]$ to $\tilde{\mathfrak{X}}^{(\rho)} \times [t_0, t_0 + T_1]$ in case $0 \leq t_0 < T < t_0 + T_1$. Thus, if $T < t_0 + T_1$ then we may replace T by $T = t_0 + T_1$ and, hence, assume that $[t_0, t_0 + T_1] \subset [0, T]$. Consequently, by Theorem 7.1, the (unique) weak solution $\mathbf{u} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ to the Cauchy problem (1.1) possesses a unique holomorphic extension from the time interval $(t_0, t_0 + T_1)$ to the complex temporal domain $t_0 + \Delta_{\vartheta'}^{T',T_1}$, such that the continuous mapping $\tilde{\mathbf{u}} : \tilde{\mathfrak{X}}^{(r_1)} \times (t_0 + \Delta_{\vartheta'}^{T',T_1}) \rightarrow \mathbb{C}^M$ constructed in Theorem 7.1 is holomorphic in the space-time domain $\tilde{\mathfrak{X}}^{(r_1)} \times (t_0 + \Delta_{\vartheta'}^{T',T_1})$. Since the interval $[t_0, t_0 + T_1] \subset [0, T]$ is arbitrary, both statements (iii₁) and (iii₂) and the complex analyticity statement (iii₃) in part (iii) of Theorem 3.4 follow from (8.1). More precisely, the desired holomorphic extension $\tilde{\mathbf{u}} : \tilde{\mathfrak{X}}^{(r_1)} \times \Delta_{\vartheta'}^{T',T} \rightarrow \mathbb{C}^M$ is obtained by shifting the temporal

domain $t_0 + \Delta_{\vartheta'}^{T', T_1}$ with the vertex t_0 ranging from the left to the right over the time interval $[0, T - T_1]$. In this process, the uniqueness result in Theorem 7.1, Part (i), guarantees that the function $\tilde{\mathbf{u}}(\cdot + iy, t) = \mathbf{u}^{(iy)}(t) = \tilde{\mathbf{u}}(\cdot + iy, t_0 + s) = \mathbf{u}^{(iy)}(t) \in \mathbf{B}^{s;p,p}(\mathbb{R}^N)$, with $0 \leq s = t - t_0 \leq T_1$, is well defined for all $(y, t) \in Q^{(r_1)} \times \Delta_{\vartheta'}^{T', T}$ independently from the particular choice of the vertex $t_0 \in [0, T - T_1]$ of the the complex temporal domain $t_0 + \Delta_{\vartheta'}^{T', T_1}$. Furthermore, in part (iii), (iii₁), (3.11) holds with ρ in place of κ_0 , whereas $r' \in (0, \kappa_0]$ has to be replaced by ρ , as well.

This concludes our proof of Theorem 3.4. \square

We conclude this section with the proof of the estimate in (3.14).

Proof of Proposition 3.5. We recall from Section 6 that the shifted continuous function $\mathbf{u}^{(iy)} : t \mapsto \mathbf{u}^{(iy)}(t) = \tilde{\mathbf{u}}(\cdot + iy, t) : [0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ is a unique weak solution of the spatially shifted Cauchy problem (6.1) with the shift $z_0 = iy$ ($y \in Q^{(r_1)}$). Consequently, $\mathbf{u}^{(iy)} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}^N))$ is also a strict solution (cf. Definition 4.4) to the following abstract initial value problem, for every $y \in Q^{(r_1)}$:

$$\begin{aligned} \frac{d}{dt} \mathbf{u}^{(iy)} - A^{(iy)}(t) \mathbf{u}^{(iy)} &= \mathbf{F}^{(iy)}(t, \mathbf{u}^{(iy)}(t)) \quad \text{for a.e. } t \in (0, T); \\ \mathbf{u}^{(iy)}(0) &= \tilde{\mathbf{u}}_0(\cdot + iy) \in E_{1-\frac{1}{p}, p} = \mathbf{B}^{s;p,p}(\mathbb{R}^N), \end{aligned} \quad (8.5)$$

cf. (6.1) and (6.17). Here, $A^{(z_0)}(t) \in \mathcal{L}(E_1 \rightarrow E_0)$ is the bounded linear (partial differential) operator introduced in (6.16), satisfying $A^{(z_0)}(t) \in \text{MR}_p(E_1 \rightarrow E_0)$ for every $t \in [0, T]$, and $\mathbf{F}^{(z_0)} : [0, T] \times U \rightarrow E_0 = \mathbf{L}^p(\mathbb{R}^N)$ stands for the ‘‘shifted’’ Nemytskii operator defined in Remark 6.4, (6.11). We recall that both constants, $C_1 \equiv C_1(C(U))$ and $L \equiv L(U)$ in inequalities (6.8) and (6.9), respectively, are independent from the shift by $z_0 \in \mathfrak{X}^{(r_1)}$ in case $x \in \mathbb{R}^N$ is replaced by $x + z_0$; with $z_0 = iy$ ($y \in Q^{(r_1)}$) in our case.

We now derive an estimate analogous to (4.11) for our shifted Cauchy problem (8.5) in place of the (original) abstract problem (4.9). Inspecting the proof of Theorem 2.1 in Clément and Li [20, pp. 20–23] and combining it with our linear perturbation result in Lemma 5.2 and the estimate in (5.6), we arrive at the following estimate for our shifted Cauchy problem (8.5) in place of the abstract problem (4.9),

$$\begin{aligned} &\int_0^T \left\| \frac{d\mathbf{u}^{(iy)}}{dt} \right\|_{E_0}^p dt + \int_0^T \left\| A^{(iy)}(0) \mathbf{u}^{(iy)}(t) \right\|_{E_0}^p dt \\ &\leq M_{p,T} \left(\left\| \mathbf{u}^{(iy)}(0) \right\|_{E_{1-\frac{1}{p}, p}}^p + \int_0^T \left\| \mathbf{F}^{(iy)}(t, \mathbf{u}^{(iy)}(t)) \right\|_{E_0}^p dt \right), \end{aligned} \quad (8.6)$$

where $M_{p,T} \in (0, \infty)$ is a constant independent from the initial data $\mathbf{u}^{(iy)}(0) = \tilde{\mathbf{u}}_0(\cdot + iy) \in E_{1-\frac{1}{p}, p} = \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ and the right-hand side $\mathbf{F}^{(iy)}(t, \mathbf{u}^{(iy)}(t))$ of (8.5), as well. Since $A^{(iy)}(0) \in \text{MR}_p(E_1 \rightarrow E_0) \subset \text{Hol}(E_1 \rightarrow E_0)$ holds by the proof of Proposition 6.5, part (d), there is a number $\lambda_0 \in \mathbb{R}_+ = [0, \infty)$, sufficiently large, such that the bounded linear operator $\lambda_0 I - A^{(iy)}(0) : E_1 \rightarrow E_0$ is an isomorphism of E_1 onto E_0 . Hence, its inverse satisfies $(\lambda_0 I - A^{(iy)}(0))^{-1} \in \mathcal{L}(E_0 \rightarrow E_1)$. We conclude that there are constants $c_1, C_1 \in (0, \infty)$ and $c_2, C_2 \in \mathbb{R}_+$ (both sufficiently large, depending on $\lambda_0 \geq 0$) such that both inequalities

$$c_1 \|u\|_{E_1} - c_2 \|u\|_{E_0} \leq \|A^{(iy)}(0)u\|_{E_0} \leq C_1 \|u\|_{E_1} + C_2 \|u\|_{E_0}$$

hold for all $u \in E_1$. Consequently, we have (respectively)

$$2^{-(p-1)}c_1\|u\|_{E_1}^p \leq \|A^{(iy)}(0)u\|_{E_0}^p + c_2^p\|u\|_{E_0}^p \quad \text{and} \quad (8.7)$$

$$\|A^{(iy)}(0)u\|_{E_0}^p \leq 2^{p-1} (C_1^p\|u\|_{E_1}^p + C_2^p\|u\|_{E_0}^p) \quad \text{for all } u \in E_1. \quad (8.8)$$

Furthermore, for every $t \in [0, T]$, we split the expression

$$\mathbf{F}^{(iy)}\left(t, \mathbf{u}^{(iy)}(t)\right) = \mathbf{F}^{(iy)}(t, \mathbf{0}) + [\mathbf{F}^{(iy)}(t, \mathbf{u}^{(iy)}(0)) - \mathbf{F}^{(iy)}(t, \mathbf{0})]$$

and apply Lemma 6.3 and Remark 6.4 to derive the following analogue of (6.10) (where we insert $\mathbf{v} = \mathbf{u}^{(iy)}(t)$):

$$|\mathbf{F}^{(iy)}\left(t, \mathbf{u}^{(iy)}(t)\right)(x)| \leq |\mathbf{F}^{(iy)}(t, \mathbf{0})(x)| + L \sum_{|\beta| \leq m} \sum_{k=1}^M \left| \frac{\partial^{|\beta|}}{\partial x^\beta} u_k^{(iy)}(x, t) \right|, \quad (8.9)$$

for all $x \in \mathbb{R}^N, t \in [0, T]$. Here,

$$\mathbf{F}^{(iy)}(t, \mathbf{0})(x) = \mathbf{f}(x + iy, t; \vec{\mathbf{0}}), \quad x \in \mathbb{R}^N, \quad \vec{\mathbf{0}} = (0)_{|\beta| \leq m} \equiv (0, \dots, 0) \in \mathbb{C}^{M\tilde{N}},$$

satisfies $\mathbf{F}^{(iy)}(t, \mathbf{0}) \in E_0 = \mathbf{L}^p(\mathbb{R}^N)$, by the L^p -integrability condition in (3.4), i.e.,

$$\|\mathbf{F}^{(iy)}(t, \mathbf{0})\|_{E_0} = \left(\int_{\mathbb{R}^N} |\mathbf{f}(x + iy, t; \vec{\mathbf{0}})|^p dx \right)^{1/p} \leq K \quad \text{for all } t \in [0, T],$$

where $K \in (0, \infty)$ is a constant. Each term $\frac{\partial^{|\beta|}}{\partial x^\beta} u_k^{(iy)}(\cdot, t) \in L^p(\mathbb{R}^N)$ on the right-hand side of (8.9) above belongs to the Besov space $B^{s-|\beta|; p, p}(\mathbb{R}^N)$ which, thanks to $|\beta| \leq m < s = 2m \left(1 - \frac{1}{p}\right) < 2m$, is continuously imbedded into another Besov space, $B^{s-|\beta|; p, p}(\mathbb{R}^N) \hookrightarrow B^{s-m; p, p}(\mathbb{R}^N) \hookrightarrow L^p(\mathbb{R}^N)$. Applying these estimates to the right-hand side of (8.9), we thus obtain

$$\|\mathbf{F}^{(iy)}(t, \mathbf{u}^{(iy)}(t))\|_{E_0}^p \leq 2^{p-1} (K^p + \gamma_{N, m, M} L^p \cdot \|\mathbf{u}^{(iy)}(t)\|_{B^{s; p, p}(\mathbb{R}^N)}^p) \quad (8.10)$$

for all $t \in [0, T]$ and every $y \in Q^{(r_1)}$, where $\gamma_{N, m, M} \in (0, \infty)$ is a numerical constant depending only on N, m , and M . Recalling $\mathbf{u}^{(iy)}(t) \in U$ for all $(y, t) \in Q^{(r_1)} \times [0, T]$ and the definition of the set $U \subset E_{1-\frac{1}{p}, p} = \mathbf{B}^{s; p, p}(\mathbb{R}^N)$ in (6.3), we conclude that the right-hand side of (8.10) can be estimated from above by a constant $C \equiv C(K, U) \in (0, \infty)$ independent from $(y, t) \in Q^{(r_1)} \times [0, T]$:

$$\|\mathbf{F}^{(iy)}(t, \mathbf{u}^{(iy)}(t))\|_{E_0}^p \leq C(K, U) \quad \text{for all } (y, t) \in Q^{(r_1)} \times [0, T]. \quad (8.11)$$

Finally, we combine this estimate with Theorem 7.1 to arrive at

$$\|\mathbf{F}^{(iy)}(t, \tilde{\mathbf{u}}(\cdot + iy, t))\|_{E_0}^p \leq \tilde{C}(K, U) \quad \text{for all } (y, t) \in Q^{(r_1)} \times \Delta_{y'}^{T', T}, \quad (8.12)$$

where $\tilde{C} \equiv \tilde{C}(K, U) \in (0, \infty)$ is a constant independent from $(y, t) \in Q^{(r_1)} \times \Delta_{y'}^{T', T}$. Recall from our proof of Theorem 3.4, part (iii), that the function $\tilde{\mathbf{u}} : \tilde{\mathfrak{X}}^{(r_1)} \times \Delta_{y'}^{T', T} \rightarrow \mathbb{C}^M$ stands for the unique holomorphic temporal extension of the function $\mathbf{u}^b : \tilde{\mathfrak{X}}^{(r_1)} \times (0, T) \rightarrow \mathbb{C}^M$ defined in formula (8.4). This extension, $\tilde{\mathbf{u}}$, satisfies $\tilde{\mathbf{u}}(x, y, t) = \mathbf{u}^{(iy)}(x, t)$ for all $(x, y, t) \in \mathbb{R}^N \times \bar{Q}^{(r_1)} \times [0, T]$ and $\tilde{\mathbf{u}}(\cdot + iy, t) \in U$ for all

$(y, t) \in \bar{Q}^{(r_1)} \times \Delta_{\vartheta'}^{T', T}$. We employ the norm defined in (3.10) and (8.12) to estimate the right-hand side of (8.6):

$$\begin{aligned} & \int_0^T \left\| \frac{d\mathbf{u}^{(iy)}}{dt} \right\|_{E_0}^p dt + \int_0^T \|A^{(iy)}(0)\mathbf{u}^{(iy)}(t)\|_{E_0}^p dt \\ & \leq M_{p,T} \left(\sup_{y \in Q^{(r_1)}} \|\tilde{\mathbf{u}}_0(\cdot + iy)\|_{B^{s;p,p}(\mathbb{R}^N)} + \tilde{C}(K, U)T \right) \end{aligned} \quad (8.13)$$

for all $(y, t) \in Q^{(r_1)} \times [0, T]$. However, in the norm

$$\mathfrak{N}^{(r_1)}(\tilde{\mathbf{u}}_0) := \sup_{y \in Q^{(r_1)}} \|\tilde{\mathbf{u}}_0(\cdot + iy)\|_{B^{s;p,p}(\mathbb{R}^N)} < \infty$$

we have $\mathbf{u}^{(iy)}(0) = \tilde{\mathbf{u}}_0(\cdot + iy) \in U \subset E_{1-\frac{1}{p}, p} = \mathbf{B}^{s;p,p}(\mathbb{R}^N)$ for every $y \in Q^{(r_1)}$ owing to our choice of the number $r_1 \in (0, r)$ being sufficiently small in (and before) Proposition 6.5. Consequently, we can estimate the right-hand side of (8.13) above by another constant $\tilde{C}' \equiv \tilde{C}'(p, T, K, U) \in (0, \infty)$ independent from $(y, t) \in Q^{(r_1)} \times [0, T]$:

$$\int_0^T \left\| \frac{d\mathbf{u}^{(iy)}}{dt} \right\|_{E_0}^p dt + \int_0^T \|A^{(iy)}(0)\mathbf{u}^{(iy)}(t)\|_{E_0}^p dt \leq \tilde{C}'(p, T, K, U).$$

We estimate the left-hand side of this inequality by (8.7) combined with $\mathbf{u}^{(iy)}(t) \in U$ for all $(y, t) \in Q^{(r_1)} \times [0, T]$, thus arriving at

$$\begin{aligned} & \int_0^T \left\| \frac{d\mathbf{u}^{(iy)}}{dt} \right\|_{E_0}^p dt + 2^{-(p-1)}c_1 \int_0^T \|\mathbf{u}^{(iy)}(t)\|_{E_1}^p dt \\ & \leq \tilde{C}'(p, T, K, U) + c_2^p \int_0^T \|\mathbf{u}^{(iy)}(t)\|_{E_0}^p dt \leq \hat{C}(p, T, K, U), \end{aligned} \quad (8.14)$$

where $\hat{C} \equiv \hat{C}(p, T, K, U) \in (0, \infty)$ is a constant independent from $(y, t) \in Q^{(r_1)} \times [0, T]$.

Within the restriction to the real time $t \in [0, T]$, the desired estimate in (3.13) is derived directly from (8.14) above for all pairs $(y, t) \in Q^{(r_1)} \times [0, T] \subset Q^{(r_1)} \times \bar{\Delta}_{\vartheta'}^{T', T}$.

To extend (3.13) to the complex time $t \in \bar{\Delta}_{\vartheta'}^{T', T}$ with $t = \sigma + i\tau$ ($\sigma, \tau \in \mathbb{R}$), we take advantage of Theorem 7.1 once again. We will consider the function $\tilde{\mathbf{u}} : \bar{\mathfrak{X}}^{(r_1)} \times \Delta_{\vartheta'}^{T', T} \rightarrow \mathbb{C}^M$ constructed in our proof of Theorem 3.4, part (iii), along the complex temporal path $\tilde{\theta} : [0, T] \rightarrow \Delta_{\vartheta'}^{T', T} : s \mapsto \tilde{\theta}(s) := s + i\varsigma_1(s/T')\tau$ which consists of two straight line segments,

$$\tilde{\theta}_1 : [0, T'] \rightarrow \Delta_{\vartheta'}^{T', T} : s \mapsto \tilde{\theta}_1(s) := \left(1 + i\frac{\tau}{T'}\right)s \quad \text{with } 0 \leq s \leq T',$$

$$\tilde{\theta}_2 : [T', T] \rightarrow \Delta_{\vartheta'}^{T', T} : s \mapsto \tilde{\theta}_2(s) := s + i\tau \quad \text{with } T' \leq s \leq T.$$

Notice that $\tilde{\theta}_1(0) = 0$, $\tilde{\theta}_1(T') = \tilde{\theta}_2(T') = T' + i\tau$, and $\tilde{\theta}_2(T) = T + i\tau$. We replace the (complex) time variable t in the original Cauchy problem (1.1) by the new (real) time variable $s \in [0, T]$, thus obtaining two new abstract differential equations with the time derivatives

$$\frac{\partial \mathbf{u}}{\partial s} = \left(1 + i\frac{\tau}{T'}\right) \frac{\partial \mathbf{u}}{\partial t}(x, t) \Big|_{t=\tilde{\theta}_1(s)} \quad \text{for } 0 \leq s \leq T', \quad (8.15)$$

$$\frac{\partial \mathbf{u}}{\partial s} = \frac{\partial \mathbf{u}}{\partial t}(x, t) \Big|_{t=\tilde{\theta}_2(s)} \quad \text{for } T' \leq s \leq T, \quad (8.16)$$

respectively.

Finally, we apply Lemma 5.2 and the estimate in (5.6) to the new problem for $0 \leq s \leq T'$, thus arriving at the desired estimate in (3.13) for $0 \leq \sigma \leq T'$. For $T' \leq s \leq T$ we can use the definition of a strict solution (Definition 4.4) directly and combine it with the estimate in (4.11) to obtain the estimate in (3.13) for $T' \leq \sigma \leq T$. We conclude that (3.13) is valid also for all pairs $(y, t) \in Q^{(r')} \times \bar{\Delta}_{\vartheta'}^{T', T}$ with $t = \sigma + i\tau$ ($\sigma, \tau \in \mathbb{R}$).

The desired estimate in (3.14) now follows from (3.13) by applying (4.5) and (4.6). Proposition 3.5 is proved. \square

9. AN APPLICATION TO A RISK MODEL IN MATHEMATICAL FINANCE

Standard models in derivative pricing, including the *Black-Scholes model* (see Black and Scholes [13] and Merton [70]) and the *Heston model* (see Heston [41]) take advantage of risk neutral valuation methods for the arbitrage-free (“fair”) price of the derivative. The methods are economically justified by riskless hedging arguments introduced in [13, 70]; see also Fouque, Papanicolaou, and Sircar [29] and Hull [45] for detailed explanations of these arguments. An important assumption of these models, which is used in most of the hedging arguments, is the possibility to borrow and lend any amount of money at a risk-free interest rate. This crucial conjecture has been questioned as a consequence of the financial crisis that started in 2007 and resulted in the bankruptcy of major financial entities like Lehman Brothers. Enron’s bankruptcy in 2001 is briefly described in [45, p. 537], Business Snapshot 23.1. Namely, traders have to take into consideration the increased chance of a *default*. For this reason many trades contain a *collateral* against default and also the pricing of non-collateralized derivatives has to be adjusted. A standard book on risk management has been written by Hull [46]. Piterbarg [73] discusses the differences or convexity adjustments between the price processes of collateralized and non-collateralized contracts which could result in *funding value adjustments* of the price processes. It is only natural that traders have different *funding costs* for transactions and try to include them in the price of the contract. Hull and White reasoned in [47] that there exists no theoretical basis for such a *funding value adjustment* (FVA). Also Burgard and Kjaer [16, 18] came to a similar conclusion, by using different arguments. However, since these theoretical arguments are not convincing from a practitioner’s point of view, but traders make the adjustments anyway, Hull and White studied the consequences of funding value adjustment in a more practice-oriented way in [48, 49]. Further common adjustments of the no-default value of a derivative are *credit value adjustments* (CVA) and the related *debit value adjustments* (DVA); see, e.g., [16, 17, 18]. In a particular trade both parties have to take the possibility of default of the counterparty into account which is the *bilateral counterparty risk*. Price-reducing credit value adjustments are made by the trader to have a collateral against default of the counterparty (e.g., a bank), whereas debit value adjustments are the corresponding adjustments made by the counterparty. The *sum of all* adjustments to the value of the derivative evaluated in the absence of default is often referred to as XVA with

$$\text{XVA} = \text{FVA} - \text{CVA} + \text{DVA}.$$

A general partial differential equation for the *adjusted value* under the bilateral counterparty risk and funding value adjustments has been derived in Burgard and

Kjaer [17, Section 3] using hedging arguments. This partial differential equation is *nonlinear* if the mark-to-market value at default is considered to be the total value of the derivative including all value adjustments (see [17, Section 4]) and *linear* if the mark-to-market value is given by the no-default value of the derivative (see [17, Section 5]). In both cases, the partial differential equations are well suited for numerical calculations of the adjusted value of the derivative; see, e.g., Arregui, Salvador, and Vázquez [8] for recent results. Following notation introduced in [17], let us consider a derivative contract with payoff H between a *trader*, B , and a *counterparty*, C , on an asset S , e.g., a stock, that is not affected in case of a default by one of the two counterparties and follows the stochastic dynamics,

$$dS_t = \mu(t)S_t dt + \sigma(t)S_t dW_t, \quad (9.1)$$

where the drift $\mu : [0, \infty) \rightarrow \mathbb{R}$ and the volatility $\sigma : [0, \infty) \rightarrow \mathbb{R}$ are positive deterministic (Borel measurable) functions and $(W_t)_{t \geq 0}$ is a one-dimensional Brownian motion. Let V denote the *fair price* (the “risk-less” value) of the derivative in the setting without default and let \hat{V} denote the *adjusted price* (the “risky” value) including *funding value adjustments* (= FVA) and *bilateral counterparty risk* (= $-CVA + DVA$),

$$\hat{V} = V + XVA = V + FVA - CVA + DVA.$$

By Itô’s formula, the generator \mathcal{A}_t of the Markov process (9.1) is the partial differential operator

$$\mathcal{A}_t = \frac{1}{2}\sigma^2 S^2 \frac{\partial^2}{\partial S^2} + (q_S - \gamma_S)S \frac{\partial}{\partial S}, \quad (9.2)$$

where γ_S is the *dividend income rate* of S and q_S represents the *financing costs* that depend on the *risk-free rate* r and repo-rate of the asset (e.g., under the Fed Repurchase Agreement (Repo)). The decisive variable in the bilateral counterparty risk models studied in [8, 16, 17, 18] is the *mark-to-market value* (cf. “close-out”), M , introduced in [17, Sect. 3, Eq. (24)]. Only two different values of M seem to be of significant interest, namely, $M = \hat{V}$ and $M = V$ as described below:

If we set the mark-to-market value at default $M = \hat{V}$, then the total value \hat{V} satisfies the *nonlinear* partial differential equation

$$\frac{\partial}{\partial t} \hat{V} + \mathcal{A}_t \hat{V} - r\hat{V} = -(1 - R_B)\lambda_B \hat{V}^- + (1 - R_C)\lambda_C \hat{V}^+ + s_F \hat{V}^+, \quad (9.3)$$

with the final value $\hat{V}(S, T) = H(S)$ at maturity time $t = T$, by [17, Sect. 4, Eq. (26)]. Here, we have abbreviated $x^+ := \max\{x, 0\}$ and $x^- := \max\{-x, 0\}$ for $x \in \mathbb{R}$; hence, $x = x^+ - x^-$. We remark that the definition of the negative part x^- of $x \in \mathbb{R}$ often differs in the literature ([8, 16, 17]); it may be used with the negative sign, i.e., $x^- = \min\{x, 0\} (\leq 0)$, whence $x = x^+ + x^-$. We will respect this convention below only when approximating the function $x \mapsto x^-$ within the sum $x = x^+ + x^-$. Otherwise we use $x^- = \max\{-x, 0\} (\geq 0)$. The parameters λ_B and λ_C are given by $\lambda_B = r_B - r$ and $\lambda_C = r_C - r$, where r_B and r_C are the *yields on recovery-less bonds* for B and C , respectively. R_B and R_C , respectively, are *recovery rates* on the derivatives’ mark-to-market value at default and $s_F = r_F - r$ is the *funding spread* between the *sellers funding rate* r_F for borrowed cash and the risk-free rate r . We refer the interested reader to Burgard and Kjaer [17, Sect. 2, pp. 2–4] for further details concerning *recovery-less bonds*.

In contrast, if we assume the mark-to-market value $M = V$, then the resulting partial differential equation for \hat{V} is *linear*, albeit inhomogeneous with source terms on the right-hand side,

$$\frac{\partial}{\partial t} \hat{V} + \mathcal{A}_t \hat{V} - (r + \lambda_B + \lambda_C) \hat{V} = (R_B \lambda_B + \lambda_C) V^- - (\lambda_B + R_C \lambda_C) V^+ + s_F V^+, \quad (9.4)$$

with the final value $\hat{V}(S, T) = H(S)$, by [17, Sect. 5, Eq. (46)]. Of course, it is assumed that the fair price of the derivative, V , is known. It is claimed in [17, Sect. 3] that the vast majority of papers on valuation of counterparty risk uses this choice ($M = V$) for contracts that follow the well known “2002 ISDA Master Agreement” initiated by the International Swaps and Derivatives Association (ISDA). From the mathematical point of view, also any *convex combination* $M = (1 - \theta) \cdot V + \theta \cdot \hat{V} = V + \theta \cdot (XVA)$ of V and \hat{V} , with a constant $\theta \in [0, 1]$, might be of economic interest, as well.

We would like to investigate the question of *market completeness* for the nonlinear model (9.3) raised for related financial market models in Davis and Obłój [21]. There, the authors have shown that the problem of *market completeness* in Mathematical Finance is closely connected to (in fact, follows from) the analyticity of the derivative price. We refer to Takáč [82, Section 8, pp. 74–83] for a survey of results regarding the correlation between market completeness and the analyticity of the solution and an application of analyticity results to the stochastic volatility model in Fouque, Papanicolaou, and Sircar [29, p. 47]. The *Heston stochastic volatility model* (Heston [41], which is more popular) is treated in Alziary and Takáč [3]. Market completeness for other stochastic volatility models is discussed in [82, Remark 8.7, pp. 82–83]. In our present work, we are primarily interested in analyticity of the solution for the nonlinear partial differential equation (9.3) since the linear case (9.4) can be studied by applying the results from [82].

The nonlinearities in (9.3) are uniformly Lipschitz continuous which enables us to apply standard existence and uniqueness results for regular, strongly parabolic semilinear Cauchy problems from, e.g., Eidel’man [25], Friedman [30, 31, 32], Pazy [72, Chapt. 6, §6.1, pp. 183–191], or Tanabe [81]. Due to the fact that the nonlinearities $\hat{V} \mapsto \hat{V}^\pm : \mathbb{R} \rightarrow \mathbb{R}$ are not real-analytic, we cannot expect any analyticity of the solution $(S, t) \mapsto \hat{V}(S, t) : (0, \infty) \times (0, \infty) \rightarrow \mathbb{R}$, neither in space nor in time. In our approach we therefore modify the functions $\hat{V} \mapsto \hat{V}^\pm : \mathbb{R} \rightarrow \mathbb{R}$ as follows: We approximate them by real-analytic functions with complex-analytic extensions to a domain $(\supset \mathbb{R})$ in the complex plane \mathbb{C} . We attempt to justify this rather “nonrigorous” step by arguing that we deal with a model in Social Sciences (Economics) where a precise nonlinear response (i.e., the reaction function of type $\hat{V} \mapsto \hat{V}^\pm : \mathbb{R} \rightarrow \mathbb{R}$) is hard to determine, while facing the dominant influence of stochastic (and possibly also random) phenomena. In our example with a single equation in one space dimension ($M = N = 1$), see (1.4), we thus replace the nonlinearity $f(u) = f^+(u) + f^-(u)$ by a suitable linear combination of real-analytic approximations of the functions $u \mapsto u^+ : \mathbb{R} \rightarrow \mathbb{R}$ and $u \mapsto -u^-$ having the same asymptotic behavior at $\pm\infty$ (as $u \rightarrow \pm\infty$) and denoted by $f^{(+)}(u)$ and $f^{(-)}(u)$, respectively, with a complex variable $u \in \mathbb{C}$. We postpone explaining the details of this modification until Example 9.2 below.

It is a common approach to replace \hat{V} as a function of (S, t) by the function $\hat{v}(X, \tau) = \hat{V}(S, t)$ that depends on the logarithm of the *asset price* $X = \ln S$ and

the *time to maturity* $\tau = T - t$. Hence, by (9.3), this function satisfies the initial value problem

$$\frac{\partial}{\partial \tau} \hat{v} - \tilde{\mathcal{A}}_t \hat{v} + r \hat{v} = -(1 - R_B) \lambda_B f^{(-)}(\hat{v}) - (1 - R_C) \lambda_C f^{(+)}(\hat{v}) - s_F f^{(+)}(\hat{v}) \quad (9.5)$$

with the initial value $\hat{v}(X, 0) = \hat{H}(X) := H(\exp(X))$ and the partial differential operator

$$\tilde{\mathcal{A}}_t = \frac{1}{2} \sigma^2 \frac{\partial^2}{\partial X^2} + \left(q_S - \gamma_S + \frac{1}{2} \sigma^2 \right) \frac{\partial}{\partial X}. \quad (9.6)$$

As an alternative to the previous variable substitution it is also possible to directly alter the stochastic processes by

$$\tilde{S}_t = e^{r(T-t)} S_t \quad \text{and} \quad \tilde{X}_t = \ln \tilde{S}_t = X_t + r(T - t),$$

which yields the same partial differential equation (9.5) and allows for a financial interpretation of the new variables.

Since the coefficients of the operator $\tilde{\mathcal{A}}_t$ defined in (9.6) are independent of the variables X and τ , the analyticity of the solution can be studied by means of the Green function; see Takáč et al. [83]. But if we replace the stochastic process (9.1) that drives the value process of the asset by a stochastic volatility process, e.g. the mean-reverting process from the classical paper of Heston [41],

$$\begin{aligned} dS_t &= \mu S_t dt + \sqrt{V_t} S_t dW_t^S, \\ dV_t &= \kappa(\theta - V_t) dt + \sigma^V \sqrt{V_t} dW_t^V, \\ \rho dt &= dW_t^S dW_t^V, \end{aligned} \quad (9.7)$$

the coefficients of the generator depend on the variables and we can no longer calculate the Green function. In the volatility process (9.7)₂ above, the parameters κ , θ , and the volatility of volatility σ^V are positive constants and $(W_t^S)_{t \geq 0}$ and $(W_t^V)_{t \geq 0}$ are one-dimensional Brownian motions correlated by a correlation factor $\rho \in [-1, 1]$ through eq. (9.7)₃. This model has been treated recently in Salvador and Oosterlee [76, 77].

Remark 9.1. Hypotheses (H1) and (H2) on the partial differential operator are consistent with the hypotheses (H1) and (H2) in Takáč [82, p. 56]. As mentioned above, the operator connected to the stochastic volatility model of Fouque, Papanicolaou, and Sircar [29, p. 47], which is parallel to (but not more general than) the Heston model (9.7), has been studied in detail in [82, Sect. 8, pp. 74–83]. Various other stochastic volatility models have been discussed in [82, Remark 8.7, pp. 82–83], as well. Hence, (H1) and (H2) are satisfied for these models; we refer the reader to [82] for further details.

According to this remark, hypotheses (H1) and (H2) are fulfilled for (9.5) even if we consider a stochastic volatility model, e.g., like (9.7), instead of (9.1). We would like to give an example for suitable nonlinearities $f^{(+)}$ and $f^{(-)}$ that satisfy the remaining hypothesis (H3) and approximate the functions $u \mapsto u^+ : \mathbb{R} \rightarrow \mathbb{R}$ and $u \mapsto -u^-$, respectively.

Our example is motivated by Takáč [82, Example 8.2, pp. 79–80]. We define the complex planar domains

$$\nabla_{\vartheta}^{(r)} := \{ \zeta = \xi e^{i\theta} + i\eta \in \mathbb{C} : \xi \in \mathbb{R}, \eta \in (-r, r), \text{ and } |\theta| < \vartheta \}, \quad (9.8)$$

$$\begin{aligned} \nabla_0^{(r)} &:= \{\zeta = \xi + i\eta \in \mathbb{C} : \xi \in \mathbb{R}, \eta \in (-r, r)\} \\ &= \cap_{0 < \vartheta < \pi/2} \nabla_\vartheta^{(r)} = \mathbb{R} + i(-r, r) \end{aligned} \tag{9.9}$$

for $r \in (0, \infty)$ and $0 < \vartheta < \pi/2$ with their respective closures in \mathbb{C} denoted by $\bar{\nabla}_\vartheta^{(r)}$ and $\bar{\nabla}_0^{(r)}$; both contain the origin $0 \in \mathbb{C}$. For any given numbers $r \in (0, \infty)$ and $0 < \vartheta < \pi/2$, the domain $\nabla_\vartheta^{(r)}$ is of the form

$$\nabla_\vartheta^{(r)} = \nabla_\vartheta^{(0)} + i(-r, r) = \cup_{\eta \in (-r, r)} (i\eta + \nabla_\vartheta^{(0)}) \subset \mathbb{C},$$

where

$$\nabla_\vartheta^{(0)} := \{\zeta = \xi e^{i\vartheta} \in \mathbb{C} : \xi \in \mathbb{R} \text{ and } |\theta| < \vartheta\} = (\Delta_\vartheta) \cup (-\Delta_\vartheta) \cup \{0\}$$

is a symmetric sector in \mathbb{C} and the open sector Δ_ϑ is defined as in (1.5). We notice that $\nabla_0^{(r)}$ is a strip in \mathbb{C} and $\mathfrak{X}^{(r)} = (\nabla_0^{(r)})^N \subset \mathbb{C}^N$ for every $r \in (0, \infty)$. At last, we define the domain

$$\mathcal{O}_1 := \mathbb{C} \setminus \{iy : y \in (-\infty, -1] \cup [1, \infty)\}$$

that contains the closure $\bar{\nabla}_\vartheta^{(r)}$ whenever $0 < r < 1$ and $0 < \vartheta < \pi/2$. The definitions of these domains follow [82, pp. 78–79].

We now give an example for functions $f^{(+)}$ and $f^{(-)}$ (approximating $v \mapsto v^+$ and $v \mapsto -v^-$, respectively) that are analytic in \mathcal{O}_1 and whose first derivatives are bounded in $\nabla_{\vartheta_0}^{(r_0)}$, whenever $r_0 \in (0, \infty)$ and $0 < \vartheta_0 < \pi/2$.

Example 9.2. We consider the functions

$$f^{(+)}(v) = v \left[\frac{1}{2} + \frac{1}{\pi} \arctan(v) \right], \quad f^{(-)}(v) = v \left[\frac{1}{2} - \frac{1}{\pi} \arctan(v) \right] \tag{9.10}$$

defined for every $v \in \mathbb{R}$. We have chosen $f^{(-)}$ such that $f^{(-)}(v) = -f^{(+)}(-v)$ and $f^{(+)}(v) + f^{(-)}(v) = v$ hold for all $v \in \mathbb{R}$ since the nonlinearities v^+ and $-v^-$ in the original equation (9.3) satisfy the same relations, $v^- = (-v)^+$ and $v^+ - v^- = v$, respectively. In addition, by (9.10), we have

$$f^{(+)}(v) = \frac{1}{\pi} v \int_{-\infty}^v \frac{dt}{1+t^2} \quad \text{and} \quad f^{(-)}(v) = \frac{1}{\pi} v \int_v^{+\infty} \frac{dt}{1+t^2} \tag{9.11}$$

defined for every $v \in \mathbb{R}$, which yields

$$f^{(+)}(v) > -\frac{1}{\pi} \quad \text{and} \quad f^{(-)}(v) < \frac{1}{\pi} \quad \text{with the limits} \quad \lim_{v \rightarrow \mp\infty} f^{(\pm)}(v) = \mp \frac{1}{\pi},$$

respectively. We could immediately extend these two (real analytic) functions to holomorphic (i.e., complex analytic) functions $f^{(+)}, f^{(-)} : \mathbb{C} \setminus \{-i, i\}$ by replacing the Lebesgue integrals $\int_{-\infty}^v \dots dt$ and $\int_v^{+\infty} \dots dt$ in (9.11) over the real domains $(-\infty, v]$ and $[v, +\infty)$ in (9.11), respectively, by the complex path integrals along some suitable (continuously differentiable) paths

$$\gamma_+ : (-\infty, 0] \rightarrow \mathbb{C} \setminus \{-i, i\} \quad \text{and} \quad \gamma_- : [0, +\infty) \rightarrow \mathbb{C} \setminus \{-i, i\}$$

connecting the points $-\infty$ with v and v with $+\infty$, respectively, whenever $v \in \mathbb{C} \setminus \{-i, i\}$, where the paths γ_+ and γ_- do neither pass through nor wind around the points $\pm i \in \mathbb{C}$, i.e., they have the winding numbers $\text{Ind}_{\gamma_+}(\pm i) = \text{Ind}_{\gamma_-}(\pm i) = 0$. As a consequence, this extension procedure could produce *multi-valued* analytic functions which is not desirable. Therefore, we prefer to perform this holomorphic

extension of the functions $f^{(+)}$ and $f^{(-)}$ below, by formulas in (9.12), as they meet our current goals better.

We calculate the derivatives

$$\begin{aligned}(f^{(+)}(v))' &= \frac{1}{2} + \frac{1}{\pi} \arctan(v) + \frac{1}{\pi} \frac{v}{1+v^2} > 0, \\ (f^{(-)}(v))' &= \frac{1}{2} - \frac{1}{\pi} \arctan(v) - \frac{1}{\pi} \frac{v}{1+v^2} < 0,\end{aligned}$$

with

$$\begin{aligned}(f^{(+)}(v))' + (f^{(-)}(v))' &= 1, \quad \lim_{v \rightarrow \infty} (f^{(+)}(v))' = \lim_{v \rightarrow -\infty} (f^{(-)}(v))' = 1, \\ \lim_{v \rightarrow -\infty} (f^{(+)}(v))' &= \lim_{v \rightarrow \infty} (f^{(-)}(v))' = 0.\end{aligned}$$

By

$$(f^{(+)}(v))'' = \frac{1}{\pi} \frac{2}{(1+v^2)^2} > 0 \quad \text{and} \quad (f^{(-)}(v))'' = -\frac{1}{\pi} \frac{2}{(1+v^2)^2} < 0,$$

respectively, $f^{(+)}$ is a strictly monotone increasing and strictly convex function, whereas $f^{(-)}$ is strictly monotone decreasing and strictly concave. We can use Takáč [82, Example 8.2, pp. 79–80] and extend $f^{(+)}$ and $f^{(-)}$ to holomorphic functions $\tilde{f}^{(+)}$ and $\tilde{f}^{(-)}$ on the domain \mathcal{O}_1 via the formulas

$$\begin{aligned}\tilde{f}^{(+)}(z) &= z \left[\frac{1}{2} + \frac{i}{2\pi} \log \left(\frac{1-iz}{1+iz} \right) \right] = z \left[\frac{1}{2} + \frac{1}{\pi} \arctan(z) \right], \\ \tilde{f}^{(-)}(z) &= z \left[\frac{1}{2} - \frac{i}{2\pi} \log \left(\frac{1-iz}{1+iz} \right) \right] = z \left[\frac{1}{2} - \frac{1}{\pi} \arctan(z) \right]\end{aligned}\tag{9.12}$$

for every $z \in \mathcal{O}_1 = \mathbb{C} \setminus \{\pm iy : y \in [1, \infty)\}$, thanks to the argument and logarithm formulas

$$\arg(1+iy) = \arctan(y) \quad \text{for } y \in \mathbb{R} \quad \text{and} \quad \log \left(\frac{1-iz}{1+iz} \right) = -2i \cdot \arctan(z)$$

for $z \in \mathcal{O}_1$. The extensions $\tilde{f}^{(+)}$ and $\tilde{f}^{(-)}$ have the restrictions $\tilde{f}^{(\pm)}|_{\mathbb{R}} = f^{(\pm)}$ to the real axis \mathbb{R} , respectively, and they are holomorphic on the domain \mathcal{O}_1 since the argument restriction $\arg \left(\frac{1-iz}{1+iz} \right) \in (-\pi, \pi)$ holds for $z \in \mathcal{O}_1$. We refer to [82, Eqs. (76)–(78), p. 79] for further discussion of the behavior of $\log \left(\frac{1-iz}{1+iz} \right)$. As a consequence of the arguments in [82, Example 8.2, pp. 79–80], we obtain $(\tilde{f}^{(+)}(z))' + (\tilde{f}^{(-)}(z))' = 1$ for $z \in \mathcal{O}_1$ together with limits

$$\begin{aligned}(\tilde{f}^{(+)}(z))' &\rightarrow 1 \quad \text{as } |z| \rightarrow \infty \text{ with } z \in \mathbb{C}, \Re z > 0, \\ (\tilde{f}^{(-)}(z))' &\rightarrow 0 \quad \text{as } |z| \rightarrow \infty \text{ with } z \in \mathbb{C}, \Re z > 0, \\ (\tilde{f}^{(+)}(z))' &\rightarrow 0 \quad \text{as } |z| \rightarrow \infty \text{ with } z \in \mathbb{C}, \Re z < 0, \\ (\tilde{f}^{(-)}(z))' &\rightarrow 1 \quad \text{as } |z| \rightarrow \infty \text{ with } z \in \mathbb{C}, \Re z < 0.\end{aligned}\tag{9.13}$$

The domain $\mathcal{O}_1 = \mathbb{C} \setminus \pm i[1, \infty)$ contains the strip $\mathbb{R} \times i(-r_0, r_0)$ for every $0 < r_0 < 1$ and the imaginary parts of $(\tilde{f}^{(+)}(z))'$ and $(\tilde{f}^{(-)}(z))'$ are uniformly bounded for $|\Im z| < r_0$. Consequently, it suffices to verify (9.13) only for $z = \Re z = x > 0$ and $x < 0$, respectively, by Cauchy's integral theorem applied to the integral formula for the function $\arctan(z)$ in \mathcal{O}_1 .

Finally, in order to approximate the functions $u \mapsto u^+ : \mathbb{R} \rightarrow \mathbb{R}$ and $u \mapsto -u^-$, we take $\varepsilon \in (0, 1)$ small enough and use the functions

$$\begin{aligned} \tilde{f}_\varepsilon^{(+)}(z) &:= \varepsilon \tilde{f}^{(+)}\left(\frac{z}{\varepsilon}\right) = z \left[\frac{1}{2} + \frac{i}{2\pi} \log \left(\frac{1 - i(z/\varepsilon)}{1 + i(z/\varepsilon)} \right) \right], \\ \tilde{f}_\varepsilon^{(-)}(z) &:= \varepsilon \tilde{f}^{(-)}\left(\frac{z}{\varepsilon}\right) = z \left[\frac{1}{2} - \frac{i}{2\pi} \log \left(\frac{1 - i(z/\varepsilon)}{1 + i(z/\varepsilon)} \right) \right], \end{aligned} \tag{9.14}$$

for every $z \in \mathcal{O}_\varepsilon := \mathbb{C} \setminus \{\pm iy : y \in [\varepsilon, \infty)\} = \varepsilon \mathcal{O}_1$, respectively. Notice that $\tilde{f}_\varepsilon^{(+)}(z) + \tilde{f}_\varepsilon^{(-)}(z) = z$. In particular, for every $u \in \mathbb{R}$ we obtain the approximation

$$\tilde{f}_\varepsilon^{(+)}(u) \rightarrow u^+ \quad \text{and} \quad \tilde{f}_\varepsilon^{(-)}(u) \rightarrow -u^- \quad \text{as } \varepsilon \rightarrow 0+.$$

This convergence is uniform on any compact interval $[-R, R] \subset \mathbb{R}$ with $0 < R < +\infty$.

Example 9.3. Another example for real analytic functions $f^{(+)}$ and $f^{(-)}$ could be obtained by means of Takáč [82, Example 8.3, p. 80]. For this purpose, we could consider the real analytic functions

$$f^{(\pm)}(v) = \frac{1}{2} v \pm \frac{1}{2} \log(\cosh(v)) \quad \text{for every } v \in \mathbb{R},$$

which have similar properties as the functions defined in (9.10).

We wish to apply our main result, Theorem 3.4, to the semilinear initial value problem in (9.5) where we choose $f^{(+)}$ and $f^{(-)}$ as in (9.10), $\tilde{f}^{(+)}$ and $\tilde{f}^{(-)}$ as in (9.12), and $\tilde{f}_\varepsilon^{(+)}$ and $\tilde{f}_\varepsilon^{(-)}$ as in (9.14), respectively. On the right-hand side of (9.5) we replace the (nonlinear) functions $\hat{v} \mapsto \hat{v}^+ : \mathbb{R} \rightarrow \mathbb{R}$ and $\hat{v} \mapsto -\hat{v}^-$, respectively, by the pair of functions $\hat{v} \mapsto \tilde{f}_\varepsilon^{(+)}(\hat{v}) : \mathbb{R} \rightarrow \mathbb{R}$ and $\hat{v} \mapsto \tilde{f}_\varepsilon^{(-)}(\hat{v})$ defined in (9.14). Here, we take $\varepsilon \in (0, 1)$ small enough, but fixed. The initial data in (9.5) are given by a payoff function $\hat{H} \in B^{s;p,p}(\mathbb{R}^N)$ with $p > 2 + \frac{N}{m} = 3$ and $s = 2m(1 - \frac{1}{p}) = 2(1 - \frac{1}{p})$ ($M = N = m = 1$). We assume that these initial data possess a holomorphic extension $\tilde{H} : \mathfrak{X}^{(\kappa_0)} \rightarrow \mathbb{C}^1$ from \mathbb{R}^1 to the complex domain $\mathfrak{X}^{(\kappa_0)} \subset \mathbb{C}^1$, for some $\kappa_0 \in (0, r_0]$, such that the function

$$\tilde{H}(\cdot + iy) : x \mapsto \tilde{H}(x + iy) : \mathbb{R}^1 \rightarrow \mathbb{C}^1$$

belongs to $\mathbf{B}^{s;p,p}(\mathbb{R}^1)$ for each $y \in Q^{(\kappa_0)}$ and has finite norm $\mathfrak{N}^{(\kappa_0)}(\tilde{H}) < \infty$ which has been defined in (3.10). In the case of a simple European call or put option, i.e., $\hat{H}(x) = (e^x - K)^+$ or $\hat{H}(x) = (K - e^x)^+$, $x \in \mathbb{R}^1$, respectively, one may use the functions $\tilde{f}_\varepsilon^{(+)}$ and $\tilde{f}_\varepsilon^{(-)}$ in order to find the desired holomorphic extension \tilde{H} of \hat{H} that satisfies the hypotheses required in Theorem 3.4, part (iii).

According Example 9.2, the nonlinearity

$$f(\hat{v}) = -(1 - R_B)\lambda_B f^{(-)}(\hat{v}) - (1 - R_C)\lambda_C f^{(+)}(\hat{v}) - s_F f^{(+)}(\hat{v})$$

on the right-hand side of (9.5) possesses a holomorphic extension

$$\tilde{f}_\varepsilon(\hat{v}) = -(1 - R_B)\lambda_B \tilde{f}_\varepsilon^{(-)}(\hat{v}) - (1 - R_C)\lambda_C \tilde{f}_\varepsilon^{(+)}(\hat{v}) - s_F \tilde{f}_\varepsilon^{(+)}(\hat{v}) \tag{9.15}$$

for all $\hat{v} \in \mathcal{O}_\varepsilon$, where $\mathcal{O}_\varepsilon = \varepsilon \cdot \mathcal{O}_1 = \mathbb{C} \setminus \{iy : y \in (-\infty, -\varepsilon] \cup [\varepsilon, \infty)\}$.

Now let us recall the definition of the complex planar domain $\nabla_\vartheta^{(r)} \subset \mathbb{C}$ in (9.8) for $r \in (0, \infty)$ and $0 < \vartheta < \pi/2$ with the closure $\bar{\nabla}_\vartheta^{(r)}$ in \mathbb{C} . We have $\bar{\nabla}_\vartheta^{(r)} \subset \mathcal{O}_\varepsilon$ whenever $0 < r < \varepsilon$ and $0 < \vartheta < \pi/2$. Moreover, both \tilde{f}_ε and its complex

derivative \tilde{f}'_ε are uniformly bounded in $\bar{\nabla}_\vartheta^{(r)}$. Thus, fixing any number $\varepsilon \in (0, 1)$ and taking $r \in (0, \varepsilon)$, we observe that both \tilde{f}_ε and \tilde{f}'_ε are uniformly bounded in $\bar{\nabla}_\vartheta^{(r)}$ and, consequently, also in the complex strip $\mathfrak{X}^{(r)}$, $\mathfrak{X}^{(r)} \subset \bar{\nabla}_\vartheta^{(r)} \subset \mathcal{O}_\varepsilon$. We stress that the number $\varepsilon \in (0, 1)$ may be chosen arbitrarily small in order to achieve a sufficiently precise approximation of the reaction function $f = f(\hat{v})$ by the holomorphic function $\tilde{f}_\varepsilon = \tilde{f}_\varepsilon(\hat{v})$ as described in Example 9.2 above. Naturally, the choice of a smaller number $\varepsilon \in (0, 1)$ diminishes the width of the strip $\mathfrak{X}^{(r)}$ according to $0 < r < \varepsilon$.

We have $\tilde{f}_\varepsilon^{(\pm)} \in A(\mathcal{O}_\varepsilon)$ and f' is bounded in $\nabla_{\vartheta_0}^{(r_0)}$ for every $r_0 \in (0, \varepsilon/2)$ and $0 < \vartheta_0 < \pi/2$. The technical estimate (3.4) is trivially satisfied, owing to $\tilde{f}_\varepsilon^{(\pm)}(0) = 0$.

Following the discussion in Section 6, we recall that the (unique) strict solution is restricted to the bounded open set $U \subset B^{s;p,p}(\mathbb{R}^N)$ defined in (6.3), which is, due to the continuous Sobolev imbedding $B^{s;p,p}(\mathbb{R}^N) \hookrightarrow L^\infty(\mathbb{R}^N)$, bounded in $L^\infty(\mathbb{R}^N)$, as well. Hence, it is convenient to loosen Hypothesis (H3) in the sense that we replace the complex plane \mathbb{C} in the assumptions by smaller domains. In particular, the function \tilde{f}_ε in (9.15) fulfills Hypothesis (H3) with such weakened assumptions. Since all requirements are satisfied, we can apply Theorem 3.4 to the initial value problem (9.5) and obtain the real analyticity of the solution.

Indeed, we apply our main result, Theorem 3.4, to the initial value problem (9.5), where we choose $f^{(+)}$ and $f^{(-)}$ as follows: We replace the functions $f^{(+)}$ and $f^{(-)}$, respectively, by their complexifications $\tilde{f}_\varepsilon^{(+)}(z)$ and $\tilde{f}_\varepsilon^{(-)}(z)$, respectively, defined in formulas (9.14) for $z \in \mathcal{O}_\varepsilon = \varepsilon \cdot \mathcal{O}_1$. Here, $\varepsilon > 0$ is as small as needed. The initial data is given by the payoff function $\tilde{H} \in B^{s;p,p}(\mathbb{R}^N)$ for $p > 3$ and $s = 2(1 - 1/p)$. The partial differential operator $\tilde{\mathcal{A}}_t$ defined in (9.6) satisfies (H1) and (H2) (with $N = 1$). If we replace $\tilde{\mathcal{A}}_t$ by the generator of a stochastic volatility process like (9.7), then (H1) and (H2) (with $N = 2$) are still fulfilled, according to Remark 9.1. For the nonlinearity $f(\hat{v})$ in (9.5), extended in (9.15) as $\tilde{f}_\varepsilon(\hat{v})$ for $\hat{v} \in \mathcal{O}_\varepsilon$, we have $\tilde{f}_\varepsilon \in A(\mathcal{O}_\varepsilon)$ and \tilde{f}'_ε is bounded in $\nabla_{\vartheta_0}^{(r_0)}$ for every $r_0 \in (0, \varepsilon/2)$ and $0 < \vartheta_0 < \pi/2$. The technical estimate (3.4) is trivially satisfied, owing to $\tilde{f}_\varepsilon^{(\pm)}(0) = 0$.

10. HISTORICAL REMARKS AND COMMENTS

The questions we studied in this paper are clearly related to the classical Cauchy-Kowalewski theorem (John [50], Chapt. 3, Sect. 3(d), pp. 73–77). It has been known since the work by Holmgren [43] that even the heat equation (in one space dimension(!)) has solutions that are *not* real analytic in the time variable (cf. Bilodeau [12, pp. 124–125]). This phenomenon is due to a possibly very rapid growth of the solutions as the spatial variable $x \in \mathbb{R}$ escapes to $\pm\infty$; to eliminate it one needs to restrict the function space, where the solutions are considered at each time moment $t \in \mathbb{R}_+$, in order to prevent a too rapid growth of the solutions as $x \rightarrow \pm\infty$. This is precisely what has been done also in our present article.

Here, the emphasis is on the analytic dependence in time t and the Cauchy problem (1.1) is viewed as an evolutionary equation in some suitable function space, e.g., $L^2(\mathbb{R})$ or $L^2(\mathbb{R}^N)$. Consequently, the solution is viewed as a vector-valued function $u : (0, T) \rightarrow L^2(\mathbb{R}^N)$ and, thus, regularity results (including analyticity results) have been obtained in this setting. The interested reader is referred to Takáč [82,

Sect. 9, pp. 83–85] for a number of pertinent references and their description; for example, Kato and Tanabe [53], Komatsu [56], Massey III [66], Masuda [68], and in particular Tanabe [81] and the references therein.

Investigation of the *smoothing* (or *regularizing*) effect in evolutionary equations of parabolic type has a long history; see e.g. Eidel'man [25], Friedman [30, 31, 32], Pazy [72], and Tanabe [81] and numerous references therein. *Analytic* smoothing (or regularizing) effects, similar to those treated in our present article, in the space (x) and/or time (t) variable(s), have been obtained somewhat later, beginning with the theory of analytic semigroups (in an abstract Banach space), see e.g. the monographs by Kato [52], Lions [62], Pazy [72], and Tanabe [81], and applying (extending) it to nonautonomous analytic evolutionary equations, see e.g. Kato and Tanabe [53], Komatsu [56], Masuda [68], and Tanabe [81]. Evolutionary equations exhibiting analytic smoothing effects may be split into the following two classes: *dissipative* and *dispersive*. Again, we refer to [82, Sect. 9, pp. 83–85] for greater details about these two classes. The results for dissipative evolutionary equations establish only analyticity with respect to the time variable $t \in (0, T) \subset \mathbb{R}$. Hayashi and Kato [39] establish an analogous time-analyticity result for the nonlinear Schrödinger equation (NLS). The early (general) treatments on the analytic smoothing effect with respect to the space variable $x \in \mathbb{R}^N$ are given in Kahane [51] and Foias and Temam [27, 28].

Finally, we mention the analyticity results by Komatsu [54] obtained for solutions to elliptic and parabolic problems in a *bounded* spatial domain $\Omega \subset \mathbb{R}^N$ (with analytic boundary $\partial\Omega$). Analyticity in the space variable x and 2-nd Gevrey class regularity (weaker than analyticity) in the time variable t are established in Cavallucci [19, Teorema 6.1, p. 166] for linear parabolic equations. Some results about the analyticity of solutions of nonlinear parabolic systems, which are related to ours, are stated in Friedman [30, Theorems 3 and 4] without proofs, and for linear elliptic systems in Morrey, Jr., and Nirenberg [71]. For the Navier-Stokes equations, such analyticity results have been established in Masuda [69] and, with respect to the space variable $x \in \mathbb{R}^N$ only, earlier in Kahane [51] and Masuda [67]. These results state local analyticity of infinitely differentiable solutions without any description of their domain of holomorphy (i.e., domain of complex analyticity). Our present article provides such description in Theorem 3.4 and so do Refs. [14, 15]. More results of global nature on the space analyticity can be found in Bardos and Benachour [11] and Grujić and Kukavica [34].

11. DISCUSSION AND POSSIBLE GENERALIZATIONS

In contrast to the analytic smoothing results established in Takáč [82, Theorem 3.3, p. 59] (for a linear parabolic problem) and Takáč et al. [83, Theorem 2.1, p. 429] (for a semilinear parabolic problem), in the present work we have focused on *preserving* the spatial analyticity of the initial data, \mathbf{u}_0 , for all times $t \in [0, T]$ as long as a (global) weak solution $\mathbf{u} \in C([0, T] \rightarrow \mathbf{B}^{s;p,p}(\mathbb{R}))$ to the Cauchy problem (1.1) exists, that is, loosely written, $\mathbf{u}(\cdot, 0) = \mathbf{u}_0$ is spatially analytic (at $t = 0$) implies $\mathbf{u}(\cdot, t)$ is spatially analytic at all times $t \in (0, T]$, even for all $t \in (0, T + T_1]$ with some $T_1 > 0$ small enough.

However, also a spatial analytic smoothing result analogous to those in [82, Theorem 3.3, p. 59] and [83, Theorem 2.1, p. 429] should hold in our present setting in the Besov space $\mathbf{B}^{s;p,p}(\mathbb{R})$, by arguments similar to those used in [83,

pp. 434–435], proof of Lemma 3.4. The Banach contraction principle can then be used in analogy with [83, pp. 437–438], Step 4 in the proof of Theorem 3.1. This approach requires separation of the linear part of the Cauchy problem (1.1) (cf. [82]) followed by an application of the Banach contraction principle to the full semilinear parabolic problem in (1.1) (cf. [83, Theorem 2.1, p. 429]).

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