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Global strong solutions for a class of heterogeneous catalysis models



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ABSTRACT

We consider a mathematical model for heterogeneous catalysis in a finite three-dimensional pore of cylinder-like geometry, with the lateral walls acting as a catalytic surface. The system under consideration consists of a diffusion–advection system inside the bulk phase and a reaction–diffusion–sorption system modeling the processes on the catalytic wall and the exchange between bulk and surface. We assume Fickian diffusion with constant coefficients, sorption kinetics with linear growth bound and a network of chemical reactions which possesses a certain triangular structure. Our main result gives sufficient conditions for the existence of a unique global strong L^2 -solution to this model, thereby extending by now classical results on reaction–diffusion systems to the more complicated case of heterogeneous catalysis.

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1. Introduction

Catalysis is a key technology in Chemical Engineering, employed not only to increase the speed of chemical reactions by up to several orders of magnitude, but also to change the selectivity in favor of a desired product against other possible output components of a chemical reaction network. In heterogeneous catalysis, the catalytic substance forms a separate phase which is advantageous concerning the separation of the products from the catalytic material. A prototypical setting, which also underlies the mathematical model below, consists of a solid phase catalyst brought into contact with a gas or liquid which carries the educts as well as the product species inside the chemical reactor. In this case, the overall chemical conversion consists of the following steps:

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- (1) the educt species are transported to the surface of the catalytic substance;
- (2) molecules of at least one educt species adsorb at the catalyst surface;
- (3) adsorbed molecules react, either with other adsorbed molecules or with molecules in the bulk phase directly adjacent to the surface;
- (4) the product molecules are desorbed.

Of course, further processes will usually appear as well. For instance, adsorbed educt molecules may desorb back into the bulk before a chemical reaction occurs, or they can be transported along the surface by means of surface diffusion processes. For a recent view on the complexity of heterogeneous catalysis modeling see [12].

In the present paper, we only consider the case of pure surface chemistry, i.e. chemical reactions are only allowed between adsorbed species. This is actually no restriction, since one may otherwise introduce an artificial adsorbed form of the reaction partner which is in the bulk adjacent to the surface and assign to it an infinite adsorption rate such that all arriving bulk molecules immediately adsorb and, hence, are available for surface reaction.

In order for a heterogeneous catalytic process to be efficient, a large surface area is required. Therefore, in classical heterogeneous catalysis with solid phase catalyst, the latter is often provided as a porous structure, e.g. in so-called packed-bed reactors. In this case, the smallest unit is a single pore, into which the educts have to be transported in order to reach the pore wall, i.e. the catalytic surface. More information on this classical reactor concept can be found, e.g., in [4,16] or [30]. In recent years, with the advent of microreactor engineering technology, new reactor designs became feasible. Due to the large area-to-volume ratio at the micro scale, multichannel microreactors with catalytic wall coatings can replace classical porous structures and still provide fast and intense diffusive transport to the channel walls in order to facilitate the reaction speed or selectivity enhancement; see, e.g., [29,26]. Since the given and precise structure of microreactors together with modern control and measurement techniques allows for defined and reproducible operating conditions, this approach is much better accessible for detailed quantitative modeling and simulation; cf. [7,8]. Structured catalytic microreactors are also employed for efficient screening of potential catalysts for new reaction pathways; see, e.g., [17]. To tap the full potential of such microsystems approaches and to intensify also more classical heterogeneous catalysis processes, realistic and sound mathematical models are required as the basis for any numerical simulation. The most fundamental question then is whether a given model is well-posed, a necessary requirement to enable any reasonable numerical treatment.

In what follows, we consider a single pore as a prototypical element, where we allow for convection through the pore with a solenoidal velocity field which is assumed to be known and to satisfy the no-slip boundary condition at fixed walls. We focus on pores having smoothly bounded cross shapes. Let therefore $\Omega := A \times (-h, h) \subset \mathbb{R}^3$ denote a finite three-dimensional cylinder of height 2h > 0 with cross section $A \subset \mathbb{R}^2$ being a bounded simply connected C^2 -domain, such that ∂A is a closed regular C^2 -curve. The boundary of Ω decomposes into bottom $\Gamma_{\rm in}$, top $\Gamma_{\rm out}$ and lateral surface Σ , standing for inflow area, outflow area and active surface. The mathematical model consists of the partial mass balances for all involved chemical components, both within the bulk phase Ω (representing the interior of the pore) and on the active surface Σ (representing the catalytic surface). Inside the bulk phase, the species mass fluxes are due to advection and diffusion, where we assume the latter to be governed by Fick's law. On the active surface, we only consider diffusive fluxes along the surface, again assuming Fick's law to be a reasonable constitutive relation. We allow for different diffusivities but the model ignores cross-diffusion effects. Let us note in passing that for high surface coverage, cross-effects between the transport of different constituents will appear which are not accounted for by our model. The mass exchange between bulk and active surface is due to ad- and desorption phenomena, which are usually modeled via kinetic relations in analogy to chemical reaction kinetics. Examples will be discussed below.

Insertion of the flux relations into the partial mass balances for continua yields the following mathematical model for the unknown concentrations (c_i, c_i^{Σ}) with i = 1, ..., N:

$$\begin{cases} \partial_{t}c_{i} + (u \cdot \nabla)c_{i} - d_{i}\Delta c_{i} = 0 & \text{in } (0,T) \times \Omega, \\ \partial_{t}c_{i}^{\Sigma} - d_{i}^{\Sigma}\Delta_{\Sigma}c_{i}^{\Sigma} = r_{i}^{\text{sorp}}(c_{i}, c_{i}^{\Sigma}) + r_{i}^{\text{ch}}(c^{\Sigma}) \text{ on } (0,T) \times \Sigma, \\ (u \cdot \nu)c_{i} - d_{i}\partial_{\nu}c_{i} = g_{i}^{\text{in}} & \text{on } (0,T) \times \Gamma_{\text{in}}, \\ -d_{i}\partial_{\nu}c_{i} = r_{i}^{\text{sorp}}(c_{i}, c_{i}^{\Sigma}) & \text{on } (0,T) \times \Sigma, \\ -d_{i}\partial_{\nu}c_{i} = 0 & \text{on } (0,T) \times \Gamma_{\text{out}}, \\ -d_{i}^{\Sigma}\partial_{\nu_{\Sigma}}c_{i}^{\Sigma} = 0 & \text{on } (0,T) \times \partial \Sigma, \\ c_{i}|_{t=0} = c_{0,i} & \text{in } \Omega, \\ c_{i}^{\Sigma}|_{t=0} = c_{0,i}^{\Sigma} & \text{on } \Sigma. \end{cases}$$

$$(1)$$

In (1), the $d_i, d_i^{\Sigma} > 0$ are given constant diffusivities, u = u(t, x) denotes the velocity, r_i^{sorp} are the sorption rate functions, r_i^{ch} the rates of molar mass production due to surface chemistry, g_i^{in} the inflow rates of molar mass, $c_{0,i}$ and $c_{0,i}^{\Sigma}$ the initial concentrations in the bulk and on the active surface and ν is the outer normal to Ω .

 (A^{vel}) Throughout this paper we assume that the velocity field satisfies

$$u \in \mathbb{U}_p^{\Omega}(T) := W^{1,p}((0,T), L^p(\Omega, \mathbb{R}^3)) \cap L^p((0,T), W^{2,p}(\Omega, \mathbb{R}^3)), \tag{2}$$

for given time T > 0. Moreover, we assume

$$u \cdot \nu \leq 0$$
 on $\Gamma_{\rm in}$, $u \cdot \nu = 0$ on Σ , $u \cdot \nu \geq 0$ on $\Gamma_{\rm out}$,

and $\operatorname{div} u = 0$ in the distributional sense.

1.1. Examples for sorption and reaction rates

We give a few examples for sorption and reaction rate functions.

(S1) Let $k_i^{\text{ad}}, k_i^{\text{de}} > 0$ denote adsorption and desorption rate constants. The simplest sorption rate is given by the linear *Henry law*, i.e.

$$r_{H,i}^{\text{sorp}}(c_i, c_i^{\Sigma}) = k_i^{\text{ad}} c_i - k_i^{\text{de}} c_i^{\Sigma}.$$

This law only applies for dilute systems.

(S2) For moderate concentrations, Langmuir's law given by

$$r_{L,i}^{\mathrm{sorp}}(c_i, c_i^{\Sigma}) = k_i^{\mathrm{ad}} c_i \left(1 - \frac{c_i^{\Sigma}}{c_{\infty,i}^{\Sigma}} \right) - k_i^{\mathrm{de}} c_i^{\Sigma}$$

may be employed. Here $c_{\infty,i}^{\Sigma} > 0$ denotes the maximum capacity for adsorption of species *i*. In an application of our main results, we actually consider a modified version; see Remark 6.4, which satisfies all of our assumptions on the sorption rate stated in Section 5.

(R1) A standard example considers a reversible chemical reaction of type $A + B \rightleftharpoons P$ with N = 3 components. If mass action kinetics is employed, the mass productions are governed by the rate function

$$r^{\mathrm{ch}}(c^{\Sigma}) = \begin{pmatrix} -k^{\mathrm{re}}(c_1^{\Sigma}c_2^{\Sigma} - \kappa c_3^{\Sigma}) \\ -k^{\mathrm{re}}(c_1^{\Sigma}c_2^{\Sigma} - \kappa c_3^{\Sigma}) \\ +k^{\mathrm{re}}(c_1^{\Sigma}c_2^{\Sigma} - \kappa c_3^{\Sigma}) \end{pmatrix}.$$

Here $k^{\rm re} > 0$ denotes the rate constant of the forward reaction, while κ is the equilibrium constant for this reaction, determined as the ratio between forward and backward reaction rates.

Due to the nonlinear coupling between bulk and surface in (1), the extension of local and global existence results from classical bulk reaction—diffusion systems (see [24] for a recent survey) to the considered advection—diffusion—sorption—reaction system is not straightforward and there are only few papers dealing with related models. In [13], a similar system but without chemical reactions has been studied. The authors have shown that an L^1 -contraction principle holds for the evolution operator and thereby weak solutions are unique. In [18], the case of fast sorption is considered in which local equilibrium between the adsorbed concentration and the adjacent bulk concentration yields an algebraic relation between these quantities. Surface chemical reactions are not included. As mentioned above, heterogeneous catalysis processes are often performed in porous media in which case homogenization is a useful technique to obtain scale-reduced models. The mathematical details of such a homogenization for periodic porous media have been worked out in [11]. Another scale-reduced model of heterogeneous catalysis itself has been analyzed in [5] concerning the existence of time-periodic solutions.

2. Main results

The main results of this paper are the local-in-time existence of a unique nonnegative strong L^p -solution and the global-in-time existence of a unique nonnegative strong L^2 -solution. In all sections, (0,T) denotes a finite time interval. For a bounded simply connected C^2 -domain $A \subset \mathbb{R}^2$, such that ∂A is a closed regular C^2 -curve, let $\Omega := A \times (-h,h)$ be a finite cylinder in \mathbb{R}^3 . The boundary of Ω decomposes into three parts, which are the lateral surface $\Sigma = \partial A \times (-h,h)$, an inflow area $\Gamma_{\rm in} = A \times \{-h\}$ and an outflow area $\Gamma_{\rm out} = A \times \{h\}$, where $\Gamma_{\rm in}, \Gamma_{\rm out}$ are the bottom and the top of the cylinder Ω , respectively. In particular, we have $\partial \Omega = \Gamma_{\rm in} \cup \overline{\Sigma} \cup \Gamma_{\rm out}$. The local existence result reads as follows. For the assumptions imposed on the sorption and reaction rates see the beginning of Section 5.

Theorem 2.1. (Local existence) Let $J=(0,T')\subset\mathbb{R}$ and $5/3< p<\infty$ with $p\neq 3$. Suppose u satisfies (A^{vel}) , r^{sorp} satisfies (A_F^{sorp}) , (A_M^{sorp}) , (A_B^{sorp}) and r^{ch} fulfills (A_F^{ch}) , (A_N^{ch}) , (A_P^{ch}) . Then for every set of data

$$\begin{split} g_i^{in} &\in W_p^{1/2-1/2p}(J, L^p(\Gamma_{in})) \cap L^p(J, W_p^{1-1/p}(\Gamma_{in})), \\ c_{0,i} &\in W_p^{2-2/p}(\Omega), \\ c_{0,i}^\Sigma &\in W_p^{2-2/p}(\Sigma), \end{split}$$

which, if p > 3, satisfies the compatibility conditions

$$\begin{split} c_{0,i}u(0) \cdot \nu - d_i\partial_{\nu}c_{0,i} &= g_i^{in}(0) & on \ \Gamma_{in}, \\ - d_i\partial_{\nu}c_{0,i} &= r_i^{sorp}(c_{0,i}, c_{0,i}^{\Sigma}) & on \ \Sigma, \\ - d_i\partial_{\nu}c_{0,i} &= 0 & on \ \Gamma_{out}, \\ - d_i^{\Sigma}\partial_{\nu_{\Sigma}}c_{0,i}^{\Sigma} &= 0 & on \ \partial \Sigma, \end{split}$$

there exists a $T^* \in (0, T')$ and a unique strong solution (c_i, c_i^{Σ}) of (1) satisfying

$$c_i \in W^{1,p}((0,T), L^p(\Omega)) \cap L^p((0,T), W^{2,p}(\Omega)),$$

 $c_i^{\Sigma} \in W^{1,p}((0,T), L^p(\Sigma)) \cap L^p((0,T), W^{2,p}(\Sigma)),$

for all $T \in (0, T^*)$. If, in addition, $g_i^{in} \leq 0$ on Γ_{in} , $c_{0,i} \geq 0$ in Ω and $c_{0,i}^{\Sigma} \geq 0$ on Σ , then c_i and c_i^{Σ} are nonnegative.

Remark 2.2. If c_i, c_i^{Σ} are not continuous functions, then $c_i \geq 0$ has to be understood in the a.e. sense with respect to Lebesgue measure on Ω , and $c_i^{\Sigma} \geq 0$ with respect to the surface measure on Σ .

The global existence is given by

Theorem 2.3. (Global existence) Let the assumptions of Theorem 2.1 be satisfied for p=2 and $T'=\infty$. Additionally, assume that r^{ch} fulfills (A_S^{ch}) (see (28)) and that $g_i^{in} \in BC((0,T) \times \Gamma_{in})$, $c_{0,i} \in BC(\Omega)$, $c_{0,i}^{\Sigma} \in BC(\Sigma)$, and that $-g_i^{in}$, $c_{0,i}$, and $c_{0,i}^{\Sigma}$ are nonnegative. Then the local solution (c_i, c_i^{Σ}) extends to a global solution of (1), i.e., for p=2 the assertions of Theorem 2.1 hold for every finite T>0.

3. Notation

Let X be a Banach space and let $A: D(A) \subset X \to X$ denote a linear and densely defined operator. Let Y be another Banach space. We denote the space of bounded linear operators mapping from X to Y by $\mathcal{L}(X,Y)$.

All appearing constants, e.g. C, M > 0 denote generic constants which may vary from line to line, as long as it is not explicitly stated otherwise.

When working in time-space sets we make use of the notation $\Omega_T := (0,T) \times \Omega$ and $\Sigma_T := (0,T) \times \Sigma$ for a finite T > 0. When dealing with half infinite cylinders we write, $\Omega_{(-\infty,h)} := A \times (-\infty,h) \subset \mathbb{R}^3$ or $\Sigma_{(-\infty,h)} := \partial A \times (-\infty,h)$ for its lateral surface.

For a Banach space X, a domain $G \subset \mathbb{R}^n$, $m \in \mathbb{N}_0 := \mathbb{N} \cup \{0\}$, and $s \in (0, \infty) \setminus \mathbb{N}$, $W^{m,p}(G, X)$ denotes the usual X-valued Sobolev space and $W_p^s(G, X)$ denotes the X-valued Sobolev–Slobodeckij space. The norm of X will be denoted by $\|\cdot\|_X$. We also set $H^k := W^{k,2}$. In the same manner we employ C^m for m-times continuously differentiable functions, and BC^m for those with bounded derivatives up to order $m \in \mathbb{N}_0$.

We denote by $f^+ := \max\{0, f\}$, $f^- := -\min\{0, f\}$ the positive and negative part of a function f. Moreover, we use the superscripts \pm to denote sets of functions, whose elements are nonnegative or non-positive; e.g., we write $L^{\infty}(\Omega)^+$ for functions which admit a bounded essential supremum on Ω and which are nonnegative a.e. in Ω . With corresponding meaning we employ, e.g., $L^{\infty}(\Omega)^-$.

Throughout this work, let $\nabla_{\Sigma} u := (\nabla u)|_{\partial\Omega} - \nu(\nu \cdot (\nabla u)|_{\partial\Omega})$ denote the surface gradient and let $\Delta_{\Sigma} u = \nabla_{\Sigma} \cdot \nabla_{\Sigma} u$ denote the Laplace–Beltrami operator on Σ .

3.1. Maximal regularity spaces

For $1 , we employ the following maximal regularity spaces. The solution spaces for the unknown functions <math>c_i, c_i^{\Sigma}$ are given by

$$\mathbb{E}_p^{\Omega}(T) := W^{1,p}((0,T), L^p(\Omega)) \cap L^p((0,T), W^{2,p}(\Omega)),$$

$$\mathbb{E}_p^{\Sigma}(T) := W^{1,p}((0,T), L^p(\Sigma)) \cap L^p((0,T), W^{2,p}(\Sigma)).$$

For the data spaces we first establish appropriate regularity classes. Then we give necessary compatibility conditions in order to guarantee well-posedness. We set

$$\begin{split} \mathbb{F}_p^{\Omega}(T) &:= L^p((0,T) \times \Omega), \\ \mathbb{F}_p^{\Sigma}(T) &:= L^p((0,T) \times \Sigma), \\ \mathbb{G}_p^{\text{in}}(T) &:= W_p^{1/2-1/2p}((0,T), L^p(\Gamma_{\text{in}})) \cap L^p((0,T), W_p^{1-1/p}(\Gamma_{\text{in}})), \\ \mathbb{G}_p^{\Sigma}(T) &:= W_p^{1/2-1/2p}((0,T), L^p(\Sigma)) \cap L^p((0,T), W_p^{1-1/p}(\Sigma)), \\ \mathbb{G}_p^{\text{out}}(T) &:= W_p^{1/2-1/2p}((0,T), L^p(\Gamma_{\text{out}})) \cap L^p((0,T), W_p^{1-1/p}(\Gamma_{\text{out}})), \end{split}$$

$$\mathbb{I}_p^{\Omega} := W_p^{2-2/p}(\Omega),$$

$$\mathbb{I}_p^{\Sigma} := W_p^{2-2/p}(\Sigma).$$

We define the tuple data space for the heterogeneous catalysis equations without initial data through

$$\mathbb{F}_p^{\Omega,\Sigma}(T) := \mathbb{F}_p^{\Omega}(T) \times \mathbb{F}_p^{\Sigma}(T) \times \mathbb{G}_p^{\mathrm{in}}(T) \times \mathbb{G}_p^{\Sigma}(T) \times \mathbb{G}_p^{\mathrm{out}}(T) \times \{0\}$$

and the tuple data space with initial spaces through

$$\mathbb{F}_{p,I}^{\Omega,\Sigma}(T) := \mathbb{F}_p^{\Omega,\Sigma}(T) \times \mathbb{I}_p^{\Omega} \times \mathbb{I}_p^{\Sigma}.$$

In some statements we also employ the Dirichlet trace space on Σ , which is given by

$$\mathbb{H}^{\Sigma}_{p}(T):=W^{1-1/2p}_{p}((0,T),L^{p}(\Sigma))\cap L^{p}((0,T),W^{2-1/p}_{p}(\Sigma)).$$

Moreover, we need subspaces of functions having zero time trace. For instance, the X-valued Sobolev–Slobodeckij space with zero time trace is defined as

$$_{0}W_{p}^{s}((0,T),X) := \{c_{i} \in W_{p}^{s}((0,T),X) : c_{i}|_{t=0} = 0\}$$
 (3)

for $s \in (0, \infty) \setminus \mathbb{N}$ and sp > 1. We use this notation for all appearing spaces whenever zero time trace makes sense and write, e.g., ${}_{0}\mathbb{E}_{p}^{\Omega}(T)$, ${}_{0}\mathbb{E}_{p}^{\Sigma}(T)$ etc.

4. Linear equations

We discuss a suitable linearization of the heterogeneous catalysis equations as it is given below in (4) and show maximal regularity by means of cylindrical L^p -theory, the surjectivity of the Neumann trace operator, and a perturbation argument. For given data

$$(f_i, f_i^{\Sigma}, g_i^{\text{in}}, g_i^{\Sigma}, g_i^{\text{out}}, 0, c_{0,i}, c_{0,i}^{\Sigma}) \in \mathbb{F}_{n,I}^{\Omega,\Sigma}(T),$$

we consider the linear system:

$$\begin{cases}
\partial_{t}c_{i} + u \cdot \nabla c_{i} - d_{i}\Delta c_{i} = f_{i} & \text{in } (0, T) \times \Omega, \\
\partial_{t}c_{i}^{\Sigma} - d_{i}^{\Sigma}\Delta_{\Sigma}c_{i}^{\Sigma} = f_{i}^{\Sigma} & \text{on } (0, T) \times \Sigma, \\
c_{i}u \cdot \nu - d_{i}\partial_{\nu}c_{i} = g_{i}^{\Sigma} & \text{on } (0, T) \times \Gamma_{\text{in}}, \\
-d_{i}\partial_{\nu}c_{i} = g_{i}^{\Sigma} & \text{on } (0, T) \times \Sigma, \\
-d_{i}\partial_{\nu}c_{i} = g_{i}^{\text{out}} & \text{on } (0, T) \times \Gamma_{\text{out}}, \\
-d_{i}^{\Sigma}\partial_{\nu_{\Sigma}}c_{i}^{\Sigma} = 0 & \text{on } (0, T) \times \partial \Sigma, \\
c_{i}(0) = c_{0,i} & \text{in } \Omega, \\
c_{i}^{\Sigma}(0) = c_{0,i}^{\Sigma} & \text{in } \Sigma.
\end{cases} \tag{4}$$

Our purpose is to solve (4) for the unknown concentrations c_i, c_i^{Σ} . The main result of this section is given by

Proposition 4.1. Let $5/3 with <math>p \neq 3$ and let T > 0 be finite. Suppose the velocity field u satisfies assumption (A^{vel}) . Then (4) admits a unique solution

$$(c_i, c_i^{\Sigma}) \in \mathbb{E}_p^{\Omega}(T) \times \mathbb{E}_p^{\Sigma}(T)$$

if and only if the data satisfy the regularity condition

$$(f_i, f_i^{\Sigma}, g_i^{in}, g_i^{\Sigma}, g_i^{out}, 0, c_{0,i}, c_{0,i}^{\Sigma}) \in \mathbb{F}_{p,I}^{\Omega,\Sigma}(T),$$

and in case of p > 3 the compatibility conditions

$$\begin{split} (u|_{t=0} \cdot \nu) c_{0,i} - d_i \partial_\nu c_{0,i} &= g_i^{in}|_{t=0} \quad on \, \Gamma_{in}, \\ - d_i \partial_\nu c_{0,i} &= g_i^{\Sigma}|_{t=0} \quad on \, \Sigma, \\ - d_i \partial_\nu c_{0,i} &= g_i^{out}|_{t=0} \quad on \, \Gamma_{out}, \\ - d_i^{\Sigma} \partial_{\nu_{\Sigma}} c_{0,i}^{\Sigma} &= 0 \qquad on \, \partial \Sigma. \end{split}$$

Additionally, the corresponding solution operator ${}_{0}\mathcal{S}_{T}$ with respect to zero time trace spaces satisfies

$$\|_{0}\mathcal{S}_{\tau}\|_{\mathscr{L}\left(_{0}\mathbb{F}_{p}^{\Omega,\Sigma}(\tau)^{N},_{0}\mathbb{E}_{p}^{\Omega}(\tau)^{N}\times_{0}\mathbb{E}_{p}^{\Sigma}(\tau)^{N}\right)} \leq M$$

with a constant M > 0 independent of $\tau < T$.

Proposition 4.1 shows that the map $\mathcal{L}: \mathbb{E} \to \mathbb{F}$, where $\mathbb{E} = \mathbb{E}_p^{\Omega}(T)^N \times \mathbb{E}_p^{\Sigma}(T)^N$ and \mathbb{F} is the subspace of $\mathbb{F}_{p,I}^{\Omega,\Sigma}(T)^N$ containing all elements which satisfy the required compatibility conditions, is an isomorphism between the Banach spaces \mathbb{E} and \mathbb{F} . In this context, one speaks of maximal regularity of problem (4), or of the operator $\mathcal{L}: \mathbb{E} \to \mathbb{F}$, acting on (c_i, c_i^{Σ}) .

Plan of the proof: System (4) decomposes into two systems: One for the bulk concentrations c_i in Ω and one for the surface concentrations c_i^{Σ} on Σ . In the first step we neglect the velocity terms $(u \cdot \nabla)c_i$ and $(u \cdot \nu)c_i$ – playing the role of perturbation terms – and consider only homogeneous boundary data, i.e. we start with

$$\begin{cases}
\partial_t c_i - d_i \Delta c_i = f_i & \text{in } (0, T) \times \Omega, \\
-d_i \partial_\nu c_i = 0 & \text{on } (0, T) \times \partial \Omega, & (i = 1, ..., N)
\end{cases}$$

$$c_i(0) = c_{0,i} & \text{in } \Omega,$$

$$(5)$$

and

$$\begin{cases}
\partial_t c_i^{\Sigma} - d_i^{\Sigma} \Delta_{\Sigma} c_i^{\Sigma} = f_i^{\Sigma} & \text{on } (0, T) \times \Sigma, \\
-d_i^{\Sigma} \partial_{\nu_{\Sigma}} c_i^{\Sigma} = 0 & \text{on } (0, T) \times \partial \Sigma, & (i = 1, ..., N)
\end{cases}$$

$$c_i^{\Sigma}(0) = c_{0,i}^{\Sigma} & \text{on } \Sigma.$$
(6)

We proceed as follows: We solve (5) and (6) separately via cylindrical L^p -theory. For more information on this topic, see [22,23,20] and [21]. Then a symmetric extension in axial direction of Ω yields the surjectivity of the Neumann trace operator and, consequently, the solvability of the inhomogeneous initial boundary value problem. By perturbation arguments the obtained result carries over to (4).

4.1. Maximal regularity of the Laplacian

We define

$$A_i := -d_i \Delta, \quad A_i : D(A_i) \subset L^p(\Omega) \to L^p(\Omega),$$

$$D(A_i) := \{c_i \in W^{2,p}(\Omega) : -d_i \partial_{\nu} c_i = 0 \text{ on } \partial \Omega\}, \quad (i = 1, ..., N)$$

and, analogously,

$$\begin{split} A_i^{\Sigma} &:= -d_i^{\Sigma} \Delta_{\Sigma}, \quad A_i^{\Sigma} : D(A_i^{\Sigma}) \subset L^p(\Sigma) \to L^p(\Sigma), \\ D(A_i^{\Sigma}) &:= \{c_i^{\Sigma} \in W^{2,p}(\Sigma) : \ -d_i^{\Sigma} \partial_{\nu_{\Sigma}} c_i^{\Sigma} = 0 \text{ on } \partial \Sigma\} \quad (i = 1, ..., N). \end{split}$$

In this section we will show that A_i admits a bounded \mathcal{H}^{∞} -calculus and that A_i^{Σ} is \mathcal{R} -sectorial with angles strictly less than $\pi/2$ which implies the desired maximal regularity, cf. [9,14,28]. For A_i we directly apply [21, Theorem 4.1]. To this end, we employ the following cylindrical decomposition. We set $V_1 := A \subset \mathbb{R}^2$ and $V_2 := (-h, h) \subset \mathbb{R}$ for the cylinder Ω with height 2h > 0.

• We first consider the problem which results on the sections for fixed x_3 , yielding the following problems on $V_1 = A \subset \mathbb{R}^2$:

$$A_{i,1}: D(A_{i,1}) \subset L^p(V_1) \to L^p(V_1), \quad A_{i,1}c_i := -d_i(\partial_{x_1}^2 + \partial_{x_2}^2)c_i,$$
$$D(A_{i,1}) := \{c_i \in W^{2,p}(V_1) : B_{i,1}c_i := -d_i\partial_{\nu}c_i = 0 \text{ on } \partial V_1\}.$$

• The direction along the cylinder leads to an operator on an interval V_2 :

$$A_{i,2}: D(A_{i,2}) \subset L^p(V_2) \to L^p(V_2), \quad A_{i,2}c_i := -d_i\partial_{x_3}^2 c_i,$$
$$D(A_{i,2}) := \{c_i \in W^{2,p}(V_2): B_{i,2}c_i := -d_i\partial_{\nu}c_i = 0 \text{ on } \partial V_2\}.$$

It is a well-known fact that both operators, $A_{i,1}$ and $A_{i,2}$ admit a bounded \mathcal{H}^{∞} -calculus with zero \mathcal{H}^{∞} -angle. For these results and an introduction of the \mathcal{H}^{∞} -calculus we refer to e.g. [9] and [14]. Thus, we are exactly in the setting of [21] in the case of the strong Neumann–Laplacian given on both intersections V_1, V_2 . Therefore [21, Theorem 4.1 a)] yields that $A_i + \delta$ for some $\delta > 0$ admits a bounded \mathcal{H}^{∞} -calculus on $L^p(\Omega)$ with \mathcal{H}^{∞} -angle $\phi_{A_i+\delta}^{\infty} < \frac{\pi}{2}$. This implies maximal regularity for (5) on finite intervals (0,T), see e.g. [9].

We continue with the discussion of A_i^{Σ} . Here we first employ a parametrization of the lateral surface and afterwards apply a result from [20]. By the assumption on the cross section of Ω , ∂A constitutes a closed regular C^2 -curve. Hence we may choose a parametrization $\psi^{\Sigma}: (0, 2\pi) \to \mathbb{R}^2$ of ∂A . Let

$$\Psi^{\Sigma}: L^p(\partial V_1) \to L^p((0, 2\pi)), \quad \Psi^{\Sigma} c_i^{\Sigma} := c_i^{\Sigma} \circ \psi^{\Sigma}$$

denote the induced pull-back. Then the Ψ^{Σ} -transformed Laplace–Beltrami operator $\Delta_{\Sigma,1}$ on ∂V_1 is a second-order elliptic differential operator on $L^p((0,2\pi))$ subject to periodic boundary conditions:

$$\begin{split} A_{i,1}^{\Sigma} : D(A_{i,1}^{\Sigma}) \subset L^{p}((0,2\pi)) \to L^{p}((0,2\pi)), \quad A_{i,1}^{\Sigma} c_{i}^{\Sigma} := \Psi^{\Sigma}(-d_{i}^{\Sigma} \Delta_{\Sigma,1} c_{i}^{\Sigma}), \\ D(A_{i,1}^{\Sigma}) = \{c_{i}^{\Sigma} \in W^{2,p}((0,2\pi)) : \ \partial_{\varphi}^{j} c_{i}^{\Sigma}|_{\varphi=0} = \partial_{\varphi}^{j} c_{i}^{\Sigma}|_{\varphi=2\pi} \ \ (j=0,1)\}. \end{split}$$

We have

Lemma 4.2. Let the pull-back Ψ^{Σ} be given as above, then

- (1) Ψ^{Σ} is an isomorphism.
- (2) $\Psi^{\Sigma}: D(\Delta_{\Sigma,1}) \to D(A_{i,1}^{\Sigma})$ is an isomorphism.

Analogously to A_i , we resolve the transformation of A_i^{Σ} into two cylindrical parts: $A_{i,1}^{\Sigma}$ defined as above and $A_{i,2}^{\Sigma} := A_{i,2}$.

Applying this time [20, Theorem 8.10] we infer that there is a shift $\delta > 0$ such that $\delta + A_i^{\Sigma}$ is \mathcal{R} -sectorial with $\phi_{A_i^{\Sigma} + \delta}^{\mathcal{R}} < \frac{\pi}{2}$. This implies maximal regularity of (6) on finite intervals (0, T) by [28, Theorem 4.2].

We summarize the results of Subsection 4.1 in

Lemma 4.3. Let $1 with <math>p \neq 3$ and T > 0 be given.

- a) System (5) admits a unique solution $c_i \in \mathbb{E}_p^{\Omega}(T)$ if and only if the data satisfies the regularity condition $f_i \in \mathbb{F}_p^{\Omega}(T)$, $c_{0,i} \in \mathbb{F}_p^{\Omega}$ and, in case of p > 3, the compatibility condition $-d_i \partial_{\nu} c_{0,i} = 0$ on $\partial \Omega$.
- b) System (6) admits a unique solution $c_i^{\Sigma} \in \mathbb{E}_p^{\Sigma}(T)$ if and only if the data satisfies the regularity condition $f_i^{\Sigma} \in \mathbb{F}_p^{\Sigma}(T)$, $c_{0,i}^{\Sigma} \in \mathbb{F}_p^{\Sigma}$ and, in case of p > 3, the compatibility condition $-d_i^{\Sigma} \partial_{\nu_{\Sigma}} c_{0,i} = 0$ on $\partial \Sigma$.

4.2. Inhomogeneous Neumann boundary conditions

We turn to the discussion of inhomogeneous boundary values. We show surjectivity of the trace operator which leads to the solvability of the corresponding inhomogeneous boundary value problem.

Lemma 4.4. Let $1 with <math>p \neq 3$ and let T > 0. Then the Neumann trace operator

$$\gamma_1: \mathbb{E}_p^{\Omega}(T) \to \mathbb{G}_p^{in}(T) \times \mathbb{G}_p^{\Sigma}(T) \times \mathbb{G}_p^{out}(T)$$
$$c \mapsto (\partial_{\nu} c|_{\Gamma_{out}}, \partial_{\nu} c|_{\Sigma}, \partial_{\nu} c|_{\Gamma_{in}})$$

is a retraction.

Proof. Let $g^{\text{in}} \in \mathbb{G}_p^{\text{in}}(T), \ g^{\Sigma} \in \mathbb{G}_p^{\Sigma}(T)$ and $g^{\text{out}} \in \mathbb{G}_p^{\text{out}}(T)$. We show that there exists a $c \in \mathbb{E}_p^{\Omega}(T)$ such that

$$\partial_{\nu}c = g^{\rm in} \quad \text{on } \Gamma_{\rm in},$$
 (7)

$$\partial_{\nu}c = g^{\Sigma} \quad \text{on } \Sigma,$$
 (8)

$$\partial_{\nu}c = g^{\text{out}} \quad \text{on } \Gamma_{\text{out}}.$$
 (9)

Note that we skip the index i for better readability. As before we work with $\Omega = A \times (-h, h)$, such that $\Gamma_{\text{in}} = A \times \{-h\}$ and $\Gamma_{\text{out}} = A \times \{h\}$ in the present section. Let us define the halfspaces

$$H_{-h} := \mathbb{R}^2 \times (-h, \infty), \quad H_h := \mathbb{R}^2 \times (-\infty, h).$$

We proceed in two steps.

Step 1. Due to [1, Chapter 5] there is an extension of g^{in} to

$$\tilde{g}^{\mathrm{in}} \in \mathbb{G}_p^{\partial H_{-h}}(T)$$

and of g^{out} to

$$\tilde{g}^{\mathrm{out}} \in \mathbb{G}_p^{\partial H_h}(T).$$

We have

$$\tilde{g}^{\mathrm{in}}(0) \in W^{1-3/p}_p(\partial H_{-h}), \quad \tilde{g}^{\mathrm{out}}(0) \in W^{1-3/p}_p(\partial H_h).$$

In case that p>3 we may choose a $\tilde{c}_0^{\text{in}}\in W_p^{2-2/p}(H_{-h})$ and $\tilde{c}_0^{\text{out}}\in W_p^{2-2/p}(H_h)$ such that

$$\partial_{\nu} \tilde{c}_{0}^{\text{in}} = \tilde{g}^{\text{in}}(0) \quad \text{on } \partial H_{-h},$$

$$\partial_{\nu} \tilde{c}_{0}^{\text{out}} = \tilde{g}^{\text{out}}(0) \quad \text{on } \partial H_{h},$$

since $\partial_{\nu}: W_p^{2-2/p}(H_{-h}) \to W_p^{1-3/p}(\partial H_{-h})$ is a retraction due to [19, Theorem 2]. In case p < 3 we let $\tilde{c}_0^{\rm in} \in W_p^{2-2/p}(H_{-h})$ and $\tilde{c}_0^{\rm out} \in W_p^{2-2/p}(H_h)$ be arbitrary. Due to [10, Proposition 4.6] we can solve the parabolic problem

$$\begin{cases} \partial_t v^{\text{in}} - \Delta v^{\text{in}} = 0 & \text{in } (0, T) \times H_{-h} \\ \partial_{\nu} v^{\text{in}} = \tilde{g}^{\text{in}} & \text{on } (0, T) \times \partial H_{-h} \\ v^{\text{in}}(0) = \tilde{c}_0^{\text{in}} & \text{in } H_{-h} \end{cases}$$

with a unique

$$v^{\mathrm{in}} \in \mathbb{E}_p^{H_{-h}}(T)$$

and analogously we solve

$$\begin{cases} \partial_t v^{\mathrm{out}} - \Delta v^{\mathrm{out}} = 0 & \text{in } (0, T) \times H_h \\ \partial_\nu v^{\mathrm{out}} = \tilde{g}^{\mathrm{out}} & \text{on } (0, T) \times \partial H_h \\ v^{\mathrm{out}}(0) = \tilde{c}_0^{\mathrm{out}} & \text{in } H_h \end{cases}$$

with a unique

$$v^{\mathrm{out}} \in \mathbb{E}_p^{H_h}(T).$$

Let the cut-off function $\zeta \in C^{\infty}(\mathbb{R}, [0, 1])$ satisfy

$$\zeta|_{(-\infty,-h/3)} = 1, \quad \zeta|_{(h/3,\infty)} = 0.$$
 (10)

Then the convex combination

$$v := \zeta v^{\mathrm{in}}|_{\Omega} + (1 - \zeta)v^{\mathrm{out}}|_{\Omega} \in \mathbb{E}_{p}^{\Omega}(T)$$

fulfills the boundary conditions on top and bottom of Ω by construction.

Step 2. It remains to show that there exists a $w \in \mathbb{E}_p^{\Omega}(T)$ such that

$$\begin{split} \partial_{\nu} w &= 0 & \text{on } \Gamma_{\text{in}}, \\ \partial_{\nu} w &= g^{\Sigma} - \partial_{\nu} v & \text{on } \Sigma, \\ \partial_{\nu} w &= 0 & \text{on } \Gamma_{\text{out}}. \end{split}$$

This problem can be reduced to an equation on a bounded C^2 -domain, which works as follows. To this end, first define Ω_{-h} as the domain resulting from extending Ω in some way boundedly and smoothly (at least in the C^2 -sense) on the top. We also set $\Sigma_{-h} := \partial \Omega_{-h} \setminus \overline{\Gamma}_{\text{in}}$. In a similar manner we define Ω_{+h} and Σ_{+h} by extending Ω suitably at the bottom. Then, let G_{\pm} denote the domains resulting from reflecting $\Omega_{\pm h}$ at h_{\pm} and set $\Gamma_{\pm} := \partial G_{\pm}$. For instance, if the cross-section A of Ω is a circle, we connect half of a ball to Ω at Γ_{out} resp. Γ_{in} . Then G_{\pm} has the form of a 'pill'. It is clear that this way we always can find a suitable extension such that G_{\pm} is of class C^2 .

Let ζ be the cut-off function with the properties given in (10). We extend $\zeta(g^{\Sigma} - \partial_{\nu}v)$ to

$$\begin{split} \tilde{g}_{-}^{\Sigma} &\in \mathbb{G}_{p}^{\Sigma_{-h}}(T) \\ &= W_{p}^{1/2 - 1/2p}((0, T), L^{p}(\Sigma_{-h})) \cap L^{p}((0, T), W_{p}^{1 - 1/p}(\Sigma_{-h})) \end{split}$$

by 0 and \tilde{g}_{-}^{Σ} to

$$\hat{g}_{-}^{\Sigma} \in \mathbb{G}_{p}^{\Gamma_{-}}(T)$$

by even reflection at -h. Note that the extension by even reflection conserves the regularity $W_p^{1-1/p}$ in axial direction. Next, note that $\hat{g}_-^{\Sigma}(0) \in W_p^{1-1/3p}(\Gamma_-)$, if p > 3. Hence we can use [19, Theorem 2] to obtain $\hat{c}_{\text{in}}^{\Sigma} \in W_p^{2-2/p}(G_-)$ with

$$\partial_{\nu}\hat{c}_{\rm in}^{\Sigma} = \hat{g}_{-}^{\Sigma}(0) \quad \text{on } \Gamma_{-},$$

if p > 3. Here we can assume w.l.o.g. that $\hat{c}_{\rm in}^{\Sigma}$ is even with respect to $\Gamma_{\rm in}$. We solve the problem

$$\begin{cases}
\partial_t w^{\text{in}} - \Delta w^{\text{in}} = 0 & \text{in } (0, T) \times G_- \\
\partial_\nu w^{\text{in}} = \hat{g}_-^{\Sigma} & \text{on } (0, T) \times \Gamma_- \\
w^{\text{in}}(0) = \hat{c}_{\text{in}}^{\Sigma} & \text{in } G_-
\end{cases}$$
(11)

by [10, Theorem 2.1] to obtain

$$w^{\mathrm{in}} \in \mathbb{E}_p^{G_-}(T).$$

Defining

$$\tilde{w}^{\mathrm{in}} := w^{\mathrm{in}}|_{\Omega_{-h}} \in \mathbb{E}_p^{\Omega_{-h}}(T),$$

we have $\partial_{\nu}\tilde{w}^{\text{in}} = \tilde{g}_{-}^{\Sigma}$ on Σ_{-h} . Since \hat{g}_{-}^{Σ} and $\hat{c}_{\text{in}}^{\Sigma}$ are even in axial direction we have $\partial_{\nu}\tilde{w}^{\text{in}} = 0$ on Γ_{in} . Analogously we proceed with Γ_{out} . Here we extend $(1 - \zeta)(g^{\Sigma} - \partial_{\nu}v)$ to obtain a

$$\tilde{w}^{\mathrm{out}} \in \mathbb{E}_p^{\Omega_{+h}}(T)$$

with $\partial_{\nu}\tilde{w}^{\text{out}} = (1 - \zeta)(g^{\Sigma} - \partial_{\nu}v)$ on Σ and $\partial_{\nu}\tilde{w}^{\text{out}} = 0$ on Γ_{out} . Let the cut-off functions $\tilde{\zeta}_{1}, \tilde{\zeta}_{2} \in C^{\infty}(\mathbb{R}, [0, 1])$ satisfy

$$\tilde{\zeta}_1 = \begin{cases} 1 : (-\infty, h/2) \\ 0 : (2h/3, \infty) \end{cases}, \qquad \tilde{\zeta}_2 = \begin{cases} 0 : (-\infty, -2h/3) \\ 1 : (-h/2, \infty) \end{cases}.$$

Then the combination

$$w := \tilde{\zeta}_1 \tilde{w}^{\mathrm{in}}|_{\Omega} + \tilde{\zeta}_2 \tilde{w}^{\mathrm{out}}|_{\Omega} \in \mathbb{E}_p^{\Omega}(T)$$

satisfies by construction $\partial_{\nu}w = 0$ on $(0, T) \times \Gamma_{\text{in}}$ and $\partial_{\nu}w = 0$ on $(0, T) \times \Gamma_{\text{out}}$. The remaining inhomogeneous boundary condition on Σ is satisfied as well, since

$$\begin{split} \partial_{\nu}w &= \tilde{\zeta}_{1}\partial_{\nu}\tilde{w}^{\mathrm{in}} + \tilde{\zeta}_{2}\partial_{\nu}\tilde{w}^{\mathrm{out}} \\ &= \tilde{\zeta}_{1}\zeta(g^{\Sigma} - \partial_{\nu}v) + \tilde{\zeta}_{2}(1 - \zeta)(g^{\Sigma} - \partial_{\nu}v) \\ &= \zeta(g^{\Sigma} - \partial_{\nu}v) + (1 - \zeta)(g^{\Sigma} - \partial_{\nu}v) = g^{\Sigma} - \partial_{\nu}v \quad \text{on } (0, T) \times \Sigma. \end{split}$$

Putting together step (i) and (ii) we define $c := v + w \in \mathbb{E}_p^{\Omega}(T)$ and obtain that c satisfies (7)–(9). Thus, we have proved that there exists a bounded linear right-inverse to γ_1 which yields that the trace operator $\gamma_1 : \mathbb{E}_p^{\Omega}(T) \to \mathbb{G}_p^{\text{in}}(T) \times \mathbb{G}_p^{\Sigma}(T) \times \mathbb{G}_p^{\text{out}}(T)$ in fact is a retraction. \square

We turn to the fully inhomogeneous Neumann system

$$\begin{cases}
\partial_{t}c_{i} - d_{i}\Delta c_{i} = f_{i} & \text{in } (0, T) \times \Omega, \\
\partial_{t}c_{i}^{\Sigma} - d_{i}^{\Sigma}\Delta_{\Sigma}c_{i}^{\Sigma} = f_{i}^{\Sigma} & \text{on } (0, T) \times \Sigma, \\
-d_{i}\partial_{\nu}c_{i} = g_{i}^{\text{in}} & \text{on } (0, T) \times \Gamma_{\text{in}}, \\
-d_{i}\partial_{\nu}c_{i} = g_{i}^{\Sigma} & \text{on } (0, T) \times \Sigma, \\
-d_{i}\partial_{\nu}c_{i} = g_{i}^{\text{out}} & \text{on } (0, T) \times \Gamma_{\text{out}}, \\
-d_{i}^{\Sigma}\partial_{\nu\Sigma}c_{i}^{\Sigma} = 0 & \text{on } (0, T) \times \partial\Sigma, \\
c_{i}(0) = c_{0,i} & \text{in } \Omega, \\
c_{i}^{\Sigma}(0) = c_{0,i}^{\Sigma} & \text{in } \Sigma.
\end{cases} (i = 1, ..., N) \tag{12}$$

Lemma 4.5. Let $1 with <math>p \neq 3$ and let T > 0 be given. Then (12) admits a unique solution

$$(c_i, c_i^{\Sigma}) \in \mathbb{E}_p^{\Omega}(T) \times \mathbb{E}_p^{\Sigma}(T)$$

if and only if the data satisfies the regularity conditions

$$(f_i, f_i^{\Sigma}, g_i^{in}, g_i^{\Sigma}, g_i^{out}, 0, c_{0,i}, c_{0,i}^{\Sigma}) \in \mathbb{F}_{p,I}^{\Omega,\Sigma}(T)$$

and, in case of p > 3, the compatibility conditions

$$\begin{split} -d_i\partial_\nu c_{0,i} &= g_i^{in}|_{t=0} \quad on \ \Gamma_{in}, \\ -d_i\partial_\nu c_{0,i} &= g_i^\Sigma|_{t=0} \quad on \ \Sigma, \\ -d_i\partial_\nu c_{0,i} &= g_i^{out}|_{t=0} \quad on \ \Gamma_{out}, \\ -d_i^\Sigma\partial_{\nu_\Sigma} c_{0,i}^\Sigma &= 0 \qquad on \ \partial \Sigma. \end{split}$$

Additionally, the corresponding solution operator $_0S_T$ with respect to homogeneous initial values satisfies

$$\|_{0}\mathcal{S}_{\tau}\|_{\mathscr{L}(_{0}\mathbb{F}_{p}^{\Omega,\Sigma}(\tau)^{N},_{0}\mathbb{E}_{p}^{\Omega}(\tau)^{N}\times_{0}\mathbb{E}_{p}^{\Sigma}(\tau)^{N})} \leq M \qquad (0 < \tau < T)$$

$$\tag{13}$$

for a constant M > 0 independent of τ .

Proof. By Lemma 4.4 for given $g_i^{\text{in}} \in \mathbb{G}_p^{\text{in}}(T), g_i^{\Sigma} \in \mathbb{G}_p^{\Sigma}(T), g_i^{\text{out}} \in \mathbb{G}_p^{\text{out}}(T)$ there exists a $c_i^1 \in \mathbb{E}_p^{\Omega}(T)$ with

$$\begin{split} -d_i\partial_\nu c_i^1 &= g_i^{\rm in} \quad \text{on } (0,T) \times \Gamma_{\rm in}, \\ -d_i\partial_\nu c_i^1 &= g_i^\Sigma \quad \text{on } (0,T) \times \Sigma, \\ -d_i\partial_\nu c_i^1 &= g_i^{\rm out} \quad \text{on } (0,T) \times \Gamma_{\rm out}. \end{split}$$

Secondly, due to Lemma 4.3 for $f_i \in \mathbb{F}_p^{\Omega}(T)$, $c_{0,i} \in \mathbb{I}_p^{\Omega}(T)$ we find a unique $c_i^2 \in \mathbb{E}_p^{\Omega}(T)$ such that

$$\begin{cases} \partial_t c_i^2 - d_i \Delta c_i^2 = f_i - (\partial_t - d_i \Delta) c_i^1 & \text{in } (0, T) \times \Omega, \\ -d_i \partial_\nu c_i^2 = 0 & \text{on } (0, T) \times \Gamma_{\text{in}}, \\ -d_i \partial_\nu c_i^2 = 0 & \text{on } (0, T) \times \Sigma, \\ -d_i \partial_\nu c_i^2 = 0 & \text{on } (0, T) \times \Gamma_{\text{out}}, \\ c_i^2|_{t=0} = c_{0,i} - c_i^1|_{t=0} & \text{in } \Omega. \end{cases}$$

By construction $c_i := c_i^1 + c_i^2 \in \mathbb{E}_p^{\Omega}(T)$ satisfies (12). By employing the extension operator in zero time trace spaces from [25, Proposition 6.1] and that its norm is independent of $\tau < T$ the estimate for the solution operator readily follows. \square

4.3. Advection terms

Let $5/3 with <math>p \neq 3$ and T > 0 be given. Assume that u satisfies (A^{vel}) . We prove Proposition 4.1 by a perturbation argument. To this end, we have to show that the results obtained in Lemma 4.5 carry over when adding the two perturbation terms $(u \cdot \nabla)c_i$ and $(u \cdot \nu)c_i$. Let $\mathbb{U}_p^{\Omega}(T)$ be given as in (2) and set

$$\mathbb{U}_p^{\text{in}}(T) := W_p^{1-1/2p}((0,T), L^p(\Gamma_{\text{in}}, \mathbb{R}^3)) \cap L^p((0,T), W_p^{2-1/p}(\Gamma_{\text{in}}, \mathbb{R}^3))$$

for the Dirichlet trace space.

Proof of Proposition 4.1. The proof is carried out in three steps.

Step 1. We estimate both perturbation terms occurring in (4). Let $\varepsilon \in (0,1)$ be sufficiently small. Then the following algebra properties hold for all 5/3 :

$$\begin{split} & \mathbb{U}_p^{\Omega}(T) \cdot W_p^{1/2 - \varepsilon/2}((0,T), L^p(\Omega,\mathbb{R}^3)) \cap L^p((0,T), W_p^{1-\varepsilon}(\Omega,\mathbb{R}^3)) \hookrightarrow \mathbb{F}_p^{\Omega}(T), \\ & \mathbb{U}_p^{\mathrm{in}}(T) \cdot W_p^{1 - 1/2p - \varepsilon/2}((0,T), L^p(\Gamma_{\mathrm{in}})) \cap L^p((0,T), W_p^{2 - 1/p - \varepsilon}(\Gamma_{\mathrm{in}})) \hookrightarrow \mathbb{G}_p^{\mathrm{in}}(T). \end{split}$$

The first embedding follows by a direct calculation and the second by taking trace of the embedding

$$\mathbb{U}_p^{\Omega}(T) \cdot W_p^{1-\varepsilon/2}((0,T), L^p(\Omega, \mathbb{R}^3)) \cap L^p((0,T), W_p^{2-\varepsilon}(\Omega, \mathbb{R}^3))$$
$$\hookrightarrow W_p^{1/2}((0,T), L^p(\Omega)) \cap L^p((0,T), W^{1,p}(\Omega)),$$

which follows by a straight forward calculation, too. For $0 < \tau < T$ we infer the following estimates:

$$\begin{aligned} \|(u \cdot \nabla)c_{i}\|_{0\mathbb{F}_{p}^{\Omega}(\tau)} &\leq C\|u\|_{\mathbb{U}_{p}^{\Omega}(T)}\|\nabla c_{i}\|_{0W_{p}^{1/2-\varepsilon/2}((0,\tau),L^{p}(\Omega,\mathbb{R}^{3}))\cap L^{p}((0,\tau),W_{p}^{1-\varepsilon}(\Omega,\mathbb{R}^{3}))} \\ &\leq C\|u\|_{\mathbb{U}_{p}^{p}(T)}\|c_{i}\|_{0W_{p}^{1-\varepsilon/2}((0,\tau),L^{p}(\Omega))\cap L^{p}((0,\tau),W_{p}^{2-\varepsilon}(\Omega))} \\ &\leq C\tau^{\eta}\|c_{i}\|_{0\mathbb{E}_{n}^{\Omega}(\tau)} \qquad (c_{i} \in {}_{0}\mathbb{E}_{p}^{\Omega}(\tau)) \end{aligned}$$

with a constant C>0 and an exponent $\eta>0$ both being independent of τ . Analogously

$$\begin{aligned} \|(u \cdot \nu)c_{i}\|_{0\mathbb{G}_{p}^{\text{in}}(\tau)} &\leq C \|u\|_{\mathbb{U}_{p}^{\text{in}}(T)} \|c_{i}\|_{0W_{p}^{1-1/2p-\varepsilon/2}((0,\tau),L^{p}(\Gamma_{\text{in}}))\cap L^{p}((0,\tau),W_{p}^{2-1/p-\varepsilon}(\Gamma_{\text{in}}))} \\ &\leq C \|u\|_{\mathbb{U}_{p}^{\Omega}(T)} \|c_{i}\|_{0W_{p}^{1-\varepsilon/2}((0,\tau),L^{p}(\Omega))\cap L^{p}((0,\tau),W_{p}^{2-\varepsilon}(\Omega))} \\ &\leq C \tau^{\eta} \|c_{i}\|_{0\mathbb{E}_{\omega}^{\Omega}(\tau)} \qquad (c_{i} \in {}_{0}\mathbb{E}_{p}^{\Omega}(\tau)) \end{aligned}$$

with a constant C>0 and an exponent $\eta>0$ both being independent of τ . It follows that the linear operator

$$\mathcal{B}: {}_{0}\mathbb{E}_{p}^{\Omega}(\tau) \overset{N}{\times} {}_{0}\mathbb{E}_{p}^{\Sigma}(\tau) \overset{N}{\longrightarrow} {}_{0}\mathbb{F}_{p}^{\Omega,\Sigma}(\tau) \overset{N}{\longrightarrow} \\ \mathcal{B}(c,c^{\Sigma}) = ((u \cdot \nabla)c_{i}, 0, (u \cdot \nu)c_{i}, 0, 0, 0)_{i=1,\dots,N}$$

may be estimated by

$$\|\mathcal{B}(c, c^{\Sigma})\|_{0\mathbb{F}_{p}^{\Omega, \Sigma}(\tau)^{N}}$$

$$\leq C\tau^{\eta}\|(c, c^{\Sigma})\|_{0\mathbb{E}_{p}^{\Omega}(\tau)^{N} \times_{0}\mathbb{E}_{p}^{\Sigma}(\tau)^{N}}, \quad ((c, c^{\Sigma}) \in {}_{0}\mathbb{E}_{p}^{\Omega}(\tau)^{N} \times_{0}\mathbb{E}_{p}^{\Sigma}(\tau)^{N})$$

$$(14)$$

with a constant C > 0 and an exponent $\eta > 0$ both independent of $\tau < T$.

Step 2. We give the construction of the solution of (4) as a sum $c_i = \hat{c}_i + \bar{c}_i$, $c_i^{\Sigma} = \hat{c}_i^{\Sigma} + \bar{c}_i^{\Sigma}$. Let $(\hat{c}_i, \hat{c}_i^{\Sigma}) \in \mathbb{E}_p^{\Omega}(\tau) \times \mathbb{E}_p^{\Sigma}(\tau)$ be the solution to (12) with g_i^{in} replaced by some \hat{g}_i^{in} satisfying the compatibility condition $-d_i\partial_{\nu}c_{0,i} = \hat{g}_i^{\text{in}}(0)$ in case p > 3 and which exists according to Lemma 4.5. Next, we set

$$\bar{f}_i = -(u \cdot \nabla)\hat{c}_i, \qquad \bar{g}_i^{\text{in}} = g_i^{\text{in}} - \hat{g}_i^{\text{in}} - (u \cdot \nu)\hat{c}_i.$$

Note that then for p > 3 the compatibility condition

$$\bar{g}_i^{\rm in}|_{t=0} = g_i^{\rm in}|_{t=0} - \hat{g}_i^{\rm in}|_{t=0} - (u|_{t=0} \cdot \nu)c_{0,i} = 0$$

is satisfied by construction. Thus, the task is reduced to prove that for $0 < \tau \le T$ there exists a unique solution $(\bar{c}_i, \bar{c}_i^{\Sigma}) \in \mathbb{E}_p^{\Omega}(\tau) \times \mathbb{E}_p^{\Sigma}(\tau)$ of

$$\begin{cases}
\partial_{t}\bar{c}_{i} + (u \cdot \nabla)\bar{c}_{i} - d_{i}\Delta\bar{c}_{i} &= \bar{f}_{i} & \text{in } (0,\tau) \times \Omega, \\
\partial_{t}\bar{c}_{i}^{\Sigma} - d_{i}^{\Sigma}\Delta_{\Sigma}\bar{c}_{i}^{\Sigma} &= 0 & \text{on } (0,\tau) \times \Sigma, \\
(u \cdot \nu)\bar{c}_{i} - d_{i}\partial_{\nu}\bar{c}_{i} &= \bar{g}_{i}^{\text{in}} & \text{on } (0,\tau) \times \Gamma_{\text{in}}, \\
-d_{i}\partial_{\nu}\bar{c}_{i} &= 0 & \text{on } (0,\tau) \times \Sigma, \\
-d_{i}\partial_{\nu}\bar{c}_{i} &= 0 & \text{on } (0,\tau) \times \Gamma_{\text{out}}, \\
-d_{i}^{\Sigma}\partial_{\nu_{\Sigma}}\bar{c}_{i}^{\Sigma} &= 0 & \text{on } (0,\tau) \times \partial \Sigma, \\
\bar{c}_{i}|_{t=0} &= 0 & \text{in } \Omega, \\
\bar{c}_{i}^{\Sigma}|_{t=0} &= 0 & \text{in } \Sigma.
\end{cases} \tag{15}$$

This will be done in the final step.

Step 3. We show the unique solvability of (15) on some interval $(0,\tau)$. The proof will show that τ is independent of the data $f_i, f_i^{\Sigma}, g_i^{\text{in}}, g_i^{\Sigma}, g_i^{\text{out}}, c_{0,i}, c_{0,i}^{\Sigma}$. Due to the linearity of the system solvability then carries over to the whole time interval (0,T).

We apply a Neumann series argument to (15). To this end, let us reformulate (15) by means of the operators ${}_{0}\mathcal{L}_{\tau}$ induced by the left-hand side of (15) and \mathcal{B} . Let

$$\bar{F}_i = (\bar{f}_i, 0, \bar{g}_i^{\text{in}}, 0, 0, 0) \in {}_0\mathbb{F}_p^{\Omega, \Sigma}(\tau)^N, \qquad (i = 1, ..., N)$$

such that (15) is equivalent to

$${}_0\mathcal{L}_\tau(\bar{c},\bar{c}^{\scriptscriptstyle\Sigma}) + \mathcal{B}(\bar{c},\bar{c}^{\scriptscriptstyle\Sigma}) = \bar{F} \qquad ((\bar{c},\bar{c}^{\scriptscriptstyle\Sigma}) \in {}_0\mathbb{E}_p^\Omega(\tau) \overset{N}{\times} {}_0\mathbb{E}_p^\Sigma(\tau)^N).$$

Due to

$$_{0}\mathcal{L}_{\tau} + \mathcal{B} = (I + \mathcal{B}_{0}\mathcal{S}_{\tau})_{0}\mathcal{L}_{\tau}$$

with the solution operator ${}_0\mathcal{S}_{\tau} = {}_0\mathcal{L}_{\tau}^{-1}$ from Lemma 4.5, the invertibility of $(I + \mathcal{B}_0\mathcal{S}_{\tau})$ from ${}_0\mathbb{F}_p^{\Omega,\Sigma}(\tau)$ to ${}_0\mathbb{E}_p^{\Omega}(\tau)^N \times {}_0\mathbb{E}_p^{\Sigma}(\tau)^N$ and, in turn, of ${}_0\mathcal{L}_{\tau} + \mathcal{B}$ from ${}_0\mathbb{E}_p^{\Omega}(\tau)^N \times {}_0\mathbb{E}_p^{\Sigma}(\tau)^N$ to ${}_0\mathbb{F}_p^{\Omega,\Sigma}(\tau)$ readily follows from (14) if we choose τ so small that $C\tau^{\eta}M < 1$ with M from (13). Note that this is possible since M is independent of $\tau < T$. \square

5. Local well-posedness

In this section we derive unique solvability of

$$\begin{cases}
\partial_{t}c_{i} + (u \cdot \nabla)c_{i} - d_{i}\Delta c_{i} = f_{i} & \text{in } (0,T) \times \Omega, \\
\partial_{t}c_{i}^{\Sigma} - d_{i}^{\Sigma}\Delta_{\Sigma}c_{i}^{\Sigma} = r_{i}^{\text{sorp}}(c_{i}, c_{i}^{\Sigma}) + r_{i}^{\text{ch}}(c^{\Sigma}) & \text{on } (0,T) \times \Sigma, \\
(u \cdot \nu)c_{i} - d_{i}\partial_{\nu}c_{i} = g_{i}^{\text{in}} & \text{on } (0,T) \times \Gamma_{\text{in}}, \\
-d_{i}\partial_{\nu}c_{i} = r_{i}^{\text{sorp}}(c_{i}, c_{i}^{\Sigma}) & \text{on } (0,T) \times \Sigma, \\
-d_{i}\partial_{\nu}c_{i} = 0 & \text{on } (0,T) \times \Gamma_{\text{out}}, \\
-d_{i}^{\Sigma}\partial_{\nu_{\Sigma}}c_{i}^{\Sigma} = 0 & \text{on } (0,T) \times \partial\Sigma, \\
c_{i}|_{t=0} = c_{0,i} & \text{in } \Omega, \\
c_{i}^{\Sigma}|_{t=0} = c_{0,i}^{\Sigma} & \text{on } \Sigma,
\end{cases}$$

$$(16)$$

in the strong L^p -sense, locally in time. Recall when the index i is used we mean i=1,...,N, and have e.g. $c=(c_i)_{i=1,...,N}$ and $c^{\Sigma}=(c_i^{\Sigma})_{i=1,...,N}$.

Sorption and reaction terms are required to satisfy the following assumptions:

 $(\mathbf{A}_{\mathbf{F}}^{\mathbf{sorp}})$ The sorption rate acts as a function satisfying

$$r_i^{\text{sorp}} = r_i^{\text{sorp}}(c_i, c_i^{\Sigma}), \quad r_i^{\text{sorp}} \in C^2(\mathbb{R}^2), \quad \nabla r_i^{\text{sorp}} \in BC^1(\mathbb{R}^2, \mathbb{R}^2).$$

 $(\mathbf{A}_{\mathbf{M}}^{\mathbf{sorp}})$ The sorption rate increases monotonically in c_i and decreases monotonically in c_i .

 $(\mathbf{A_B^{sorp}})$ The sorption rate admits linear bounds

$$-k_i^{\text{de}} c_i^{\Sigma} \le r_i^{\text{sorp}}(c_i, c_i^{\Sigma}) \le k_i^{\text{ad}} c_i \qquad (c_i, c_i^{\Sigma} \ge 0)$$

for given adsorption and desorption constants $k_i^{\text{ad}}, k_i^{\text{de}} > 0$.

 $(\mathbf{A}_{\mathbf{F}}^{\mathbf{ch}})$ We assume that the chemical reactions fulfill

$$r_i^{\mathrm{ch}} = r_i^{\mathrm{ch}}(c^{\scriptscriptstyle \Sigma}), \quad r^{\mathrm{ch}} \in C^1([0,\infty)^N, \mathbb{R}^N), \qquad i = 1,..,N.$$

 $(\mathbf{A_N^{ch}})$ The reaction is supposed to be quasi-positive, i.e.

$$r_i^{\text{ch}}(y) \ge 0, \qquad (y \in [0, \infty)^N, \ y_i = 0).$$

(A^{ch}_P) The reaction admits polynomial growth, i.e. there exist a constant M>0 and an exponent $\gamma\in[1,\infty)$ if $p\in[2,\infty)$ and $\gamma=\frac{1}{1-p/2}\in[1,6]$ if $p\in(5/3,2)$, such that

$$|r^{\text{ch}}(y)| \le M (1 + |y|^{\gamma}) \qquad (y \in [0, \infty)^N).$$

Additionally, suppose the Jacobian fulfills

$$|(r^{\text{ch}})'(y)| \le M (1 + |y|^{\gamma - 1}) \qquad (y \in [0, \infty)^N).$$

Remark 5.1. Let us comment on the polynomial growth conditions of r^{ch} and $(r^{ch})'$.

a) In (A_P^{ch}) we restrict to $\gamma \in [1,6]$ if $p \in (5/3,2)$. This is due to the embedding

$$_{0}\mathbb{E}_{n}^{\Sigma}(T) \hookrightarrow L^{p\gamma}(\Sigma_{T})$$

for p > 5/3, cf. [2], which we employ in the proof of the local existence result. In case $p \ge 2$ only an arbitrary polynomial growth is required.

b) The growth rate γ of $r^{\rm ch}$ in $(A_{\rm P}^{\rm ch})$ yields that $r^{\rm ch}$ acts as a Nemytskij operator

$$r^{\mathrm{ch}}: L^{p\gamma}(\Sigma_T)^N \to L^p(\Sigma_T)^N,$$

cf. [3, Theorem 3.1]. Analogously the growth rate $\gamma - 1$ of $(r^{\text{ch}})'$ in (A_P^{ch}) yields

$$(r^{\mathrm{ch}})': L^{p\gamma}(\Sigma_T)^N \to L^{p\gamma/(\gamma-1)}(\Sigma_T)^{N\times N}.$$

In particular, $(r^{\text{ch}})'$ maps a ball \bar{B}_{δ} of radius $\delta > 0$ in $L^{p\gamma}(\Sigma_T)^N$ into a ball of radius $k(\delta) > 0$ in $L^{p\gamma/(\gamma-1)}(\Sigma_T)^{N\times N}$. Hence an application of the mean value theorem to the function r^{ch} and Hölder's inequality with

$$\frac{1}{p} = \frac{1}{p\gamma/(\gamma - 1)} + \frac{1}{p\gamma}$$

yields

$$\begin{aligned} & \|r^{\operatorname{ch}}(c^{\Sigma}) - r^{\operatorname{ch}}(z^{\Sigma})\|_{L^{p}(\Sigma_{T})^{N}} \\ & \leq \sup_{v^{\Sigma} \in \bar{B}_{\delta}} \|(r^{\operatorname{ch}})'(v^{\Sigma})\|_{L^{p\gamma/(\gamma-1)}(\Sigma_{T})^{N \times N}} \|c^{\Sigma} - z^{\Sigma}\|_{L^{p\gamma}(\Sigma_{T})^{N}} \\ & \leq k(\delta) \|c^{\Sigma} - z^{\Sigma}\|_{L^{p\gamma}(\Sigma_{T})^{N}} \qquad (c^{\Sigma}, z^{\Sigma} \in \bar{B}_{\delta}), \end{aligned}$$

i.e. $r^{\rm ch}$ acts as a locally Lipschitz continuous Nemytskij operator, cf. [3, Theorem 3.10].

For $y \in \mathbb{R}^N$ let us denote $y^+ = (y_i^+)_{i=1,\dots,N}$ where as before $y_i^+ = \max\{0, y_i\}$. Since we do not know a priori whether a corresponding solution (c, c^{Σ}) is nonnegative, we extend r^{ch} as

$$r_{i,+}^{\mathrm{ch}}:\mathbb{R}^N\to\mathbb{R}, \qquad r_{i,+}^{\mathrm{ch}}(y):=r_i^{\mathrm{ch}}(y^+) \qquad (y\in\mathbb{R}^N).$$

Then (16) remains meaningful even if c or c^{Σ} take negative values.

The outcome of this section is the proof of Theorem 2.1. It is based on maximal regularity of the linear system (4) proved in Section 4 and the contraction mapping principle. We start by proving the nonnegativity of c_i and of c_i^{Σ} .

5.1. Nonnegativity of concentrations

For nonnegative initial concentrations we have the following result.

Lemma 5.2. (Nonnegativity) Let 1 . Let <math>T > 0 and let $c_{0,i} \in \mathbb{I}_p^{\Omega,+}$, $c_{0,i}^{\Sigma} \in \mathbb{I}_p^{\Sigma,+}$ and $g_i^{in} \in \mathbb{G}_p^{in}(T)^-$ be given. Suppose u satisfies (A^{vel}) , r^{sorp} satisfies (A_F^{sorp}) , (A_M^{sorp}) , (A_B^{sorp}) and r^{ch} fulfills (A_F^{ch}) , (A_N^{ch}) , (A_P^{ch}) . Moreover, suppose $(c_i, c_i^{\Sigma}) \in \mathbb{E}_p^{\Omega}(T) \times \mathbb{E}_p^{\Sigma}(T)$ is a strong L^p -solution of (16). Then

$$c_i \geq 0$$
 a.e. in Ω_T , $c_i^{\Sigma} \geq 0$ a.e. on Σ_T

hold true.

Proof. Let $\phi_{\varepsilon} \in C^{\infty}(\mathbb{R})$ be a pointwise approximation of

$$\phi(r) = \left\{ \begin{array}{l} -r : r \leq 0 \\ 0 : r > 0 \end{array} \right.,$$

as $\varepsilon \to 0+$, which satisfies $\phi_{\varepsilon} \ge 0$, $\phi_{\varepsilon}' \le 0$ and $\phi_{\varepsilon}'' \ge 0$, e.g.

$$\phi_{\varepsilon}(r) := \left\{ \begin{array}{c} -re^{\varepsilon/r} : r \leq 0 \\ 0 : r > 0 \end{array} \right.$$

Then we have for $c_i^- = \max\{0, -c_i\}$ that

$$\phi_{\varepsilon}(c_i) \to c_i^-, \qquad c_i \phi_{\varepsilon}'(c_i) \to c_i \begin{cases} -1 : c_i < 0 \\ 0 : c_i \ge 0 \end{cases} = c_i^-$$

$$\tag{17}$$

as $\varepsilon \to 0+$. We show

$$\lim_{\varepsilon \to 0} \left[\int_{\Omega} \phi_{\varepsilon}(c_i) dx + \int_{\Sigma} \phi_{\varepsilon}(c_i^{\Sigma}) d\sigma \right] \le 0.$$

Applying ϕ_{ε} to c_i we obtain by partial integration

$$\frac{d}{dt} \int_{\Omega} \phi_{\varepsilon}(c_{i}) dx = \int_{\Omega} \phi'_{\varepsilon}(c_{i}) \partial_{t} c_{i} dx = \int_{\Omega} \phi'_{\varepsilon}(c_{i}) \operatorname{div} (d_{i} \nabla c_{i} - uc_{i}) dx$$

$$= \int_{\partial \Omega} \phi'_{\varepsilon}(c_{i}) (d_{i} \partial_{\nu} c_{i} - (u \cdot \nu) c_{i}) d\sigma - \int_{\Omega} \phi''_{\varepsilon}(c_{i}) \nabla c_{i} \cdot (d_{i} \nabla c_{i} - uc_{i}) dx$$

$$= -\int_{\Gamma_{\text{in}}} \phi'_{\varepsilon}(c_{i}) g_{i}^{\text{in}} d\sigma - \int_{\Sigma} \phi'_{\varepsilon}(c_{i}) r_{i}^{\text{sorp}} d\sigma - \int_{\Gamma_{\text{out}}} \phi'_{\varepsilon}(c_{i}) (u \cdot \nu) c_{i} d\sigma$$

$$- d_{i} \int_{\Omega} \phi''_{\varepsilon}(c_{i}) |\nabla c_{i}|^{2} dx + \int_{\Omega} \phi''_{\varepsilon}(c_{i}) \nabla c_{i} u c_{i} dx \tag{18}$$

due to the boundary conditions. In the same way we have

$$\frac{d}{dt} \int_{\Sigma} \phi_{\varepsilon}(c_{i}^{\Sigma}) d\sigma = \int_{\Sigma} \phi_{\varepsilon}'(c_{i}^{\Sigma}) \partial_{t} c_{i}^{\Sigma} d\sigma = \int_{\Sigma} \phi_{\varepsilon}'(c_{i}^{\Sigma}) (d_{i}^{\Sigma} \Delta_{\Sigma} c_{i}^{\Sigma} + r_{i}^{\text{sorp}}(c_{i}, c_{i}^{\Sigma}) + r_{i,+}^{\text{ch}}(c^{\Sigma})) d\sigma$$

$$= -d_{i}^{\Sigma} \int_{\Sigma} \phi_{\varepsilon}''(c_{i}^{\Sigma}) |\nabla_{\Sigma} c_{i}^{\Sigma}|^{2} d\sigma + \int_{\Sigma} \phi_{\varepsilon}'(c_{i}^{\Sigma}) (r_{i}^{\text{sorp}}(c_{i}, c_{i}^{\Sigma}) + r_{i,+}^{\text{ch}}(c^{\Sigma})) d\sigma. \tag{19}$$

Let us go through all the integrals appearing on the right-hand side of (18) and (19). The first and the fourth integrals on the right-hand side of (18) and the first integral on the right-hand side of (19) are negative or zero such that we may drop them. The remaining four integrals are treated as follows: We combine the sorption boundary integrals to

$$\int\limits_{\Sigma} (\phi_{\varepsilon}'(c_i^{\Sigma}) - \phi_{\varepsilon}'(c_i)) r_i^{\text{sorp}}(c_i, c_i^{\Sigma}) d\sigma$$

and split up this integral into three integrals on

$$\{\operatorname{sign}(c_i) = \operatorname{sign}(c_i^{\Sigma})\}, \qquad \{c_i \le 0 \le c_i^{\Sigma}\}, \qquad \{c_i^{\Sigma} \le 0 \le c_i\}.$$

When c_i and c_i^{Σ} have the same sign this integral tends to 0 as $\varepsilon \to 0+$. On $\{c_i \leq 0 \leq c_i^{\Sigma}\}$ we have

$$\int_{\Sigma} \mathbb{1}_{\{c_i \le 0 \le c_i^{\Sigma}\}} (\phi_{\varepsilon}'(c_i^{\Sigma}) - \phi_{\varepsilon}'(c_i)) r_i^{\text{sorp}}(c_i, c_i^{\Sigma}) d\sigma$$

$$\to \int_{\Sigma} r_i^{\text{sorp}}(c_i, c_i^{\Sigma}) d\sigma \le \int_{\Sigma} r_i^{\text{sorp}}(0, c_i^{\Sigma}) d\sigma \le 0$$

as $\varepsilon \to 0+$ by monotonicity of r_i^{sorp} and $r_i^{\text{sorp}}(0,c_i^{\Sigma}) \leq 0$. In the same way on $\{c_i^{\Sigma} \leq 0 \leq c_i\}$ we obtain

$$\int_{\Sigma} \mathbb{1}_{\{c_i^{\Sigma} \le 0 \le c_i\}} (\phi_{\varepsilon}'(c_i^{\Sigma}) - \phi_{\varepsilon}'(c_i)) r_i^{\text{sorp}}(c_i, c_i^{\Sigma}) d\sigma$$

$$\rightarrow -\int_{\Sigma} r_i^{\text{sorp}}(c_i, c_i^{\Sigma}) d\sigma \le -\int_{\Sigma} r_i^{\text{sorp}}(c_i, 0) d\sigma \le 0$$

as $\varepsilon \to 0+$ by monotonicity of r_i^{sorp} and $r_i^{\text{sorp}}(c_i,0) \ge 0$. Because of $u \cdot \nu \ge 0$ on Γ_{out} we see that

$$-\int_{\Gamma_{\text{out}}} \phi'_{\varepsilon}(c_i)(u \cdot \nu)c_i d\sigma \to -\int_{\Gamma_{\text{out}}} c_i^-(u \cdot \nu)d\sigma \le 0.$$

We treat the reaction boundary integral by the quasi-positivity of r^{ch} as follows. We show

$$\int\limits_{\Sigma} \phi_{\varepsilon}'(c_{i}^{\Sigma}) r_{i,+}^{\mathrm{ch}}(c^{\Sigma}) d\sigma \leq 0$$

through

$$\begin{split} \int\limits_{\Sigma} \phi_{\varepsilon}'(c_{i}^{\Sigma}) r_{i,+}^{\text{ch}}(c^{\Sigma}) d\sigma &= \int\limits_{\Sigma} \phi_{\varepsilon}'(c_{i}^{\Sigma}) \mathbb{1}_{\{c_{i}^{\Sigma} > 0\}} r_{i,+}^{\text{ch}}(c^{\Sigma}) d\sigma \\ &+ \int\limits_{\Sigma} \phi_{\varepsilon}'(c_{i}^{\Sigma}) \mathbb{1}_{\{c_{i}^{\Sigma} = 0\}} r_{i,+}^{\text{ch}}(c^{\Sigma}) d\sigma + \int\limits_{\Sigma} \phi_{\varepsilon}'(c_{i}^{\Sigma}) \mathbb{1}_{\{c_{i}^{\Sigma} < 0\}} r_{i,+}^{\text{ch}}(c^{\Sigma}) d\sigma. \end{split}$$

The first integral vanishes by the properties of ϕ'_{ε} , the second one is less than or equal to zero by quasipositivity and $\phi'_{\varepsilon}(0) \leq 0$ as $\varepsilon \to 0+$. The third one is less than or equal to zero by definition of the extension of r^{ch} to \mathbb{R}^N , i.e. because of $c_i^{\Sigma} < 0$ implies $c_i^{\Sigma^+} = 0$ and $r_i^{\text{ch}}(c^{\Sigma^+}) \geq 0$ by quasi-positivity. We turn to the remaining integral $\int_{\Omega} \phi''_{\varepsilon}(c_i) \nabla c_i u c_i dx$. Here we make use of $\nabla(\phi'_{\varepsilon}(c_i)) = \phi''_{\varepsilon}(c_i) \nabla c_i$ and in the same manner of $\nabla(\phi_{\varepsilon}(c_i)) = \phi'_{\varepsilon}(c_i) \nabla c_i$ and integrate by parts twice, such that

$$\int_{\Omega} \phi_{\varepsilon}''(c_{i}) \nabla c_{i} \cdot u c_{i} dx = \int_{\Omega} \nabla (\phi_{\varepsilon}'(c_{i})) \cdot u c_{i} dx$$

$$= \int_{\partial \Omega} \phi_{\varepsilon}'(c_{i}) (u \cdot \nu) c_{i} d\sigma - \int_{\Omega} \phi_{\varepsilon}'(c_{i}) \underbrace{\text{div } (u c_{i})}_{=u \cdot \nabla c_{i}} dx$$

$$= \int_{\partial\Omega} \phi_{\varepsilon}'(c_i)(u \cdot \nu)c_i d\sigma - \int_{\Omega} \nabla(\phi_{\varepsilon}(c_i))u dx$$
$$= \int_{\partial\Omega} \phi_{\varepsilon}'(c_i)(u \cdot \nu)c_i d\sigma - \int_{\partial\Omega} \phi_{\varepsilon}(c_i)(u \cdot \nu) d\sigma,$$

where in the second and in the last step we made use of div u = 0. Employing (17) we see

$$\int_{\Omega} \phi_{\varepsilon}''(c_i) \nabla c_i \cdot u c_i dx = \int_{\partial \Omega} (\phi_{\varepsilon}'(c_i) c_i - \phi_{\varepsilon}(c_i)) (u \cdot \nu) d\sigma \to 0$$

as $\varepsilon \to 0+$. Therefore summing up (18) and (19), integration in time over [0, t] and taking the limit $\varepsilon \to 0+$ yields

$$\int_{\Omega} c_i^-(t)dx + \int_{\Sigma} c_i^{\Sigma^-}(t)d\sigma = \int_{\Omega} \phi(c_i(t))dx + \int_{\Sigma} \phi(c_i^{\Sigma}(t))d\sigma$$

$$\leq \int_{\Omega} \phi(c_{0,i})dx + \int_{\Sigma} \phi(c_{0,i}^{\Sigma})d\sigma = \int_{\Omega} (c_{0,i})^- dx + \int_{\Sigma} (c_{0,i}^{\Sigma})^- d\sigma = 0$$

which in turn gives $c_i^- = 0$ a.e. in Ω_T , $c_i^{\Sigma^-} = 0$ a.e. on Σ_T and therefore $c_i \geq 0$ a.e. in Ω_T , $c_i^{\Sigma} \geq 0$ a.e. on Σ_T . Note that for $\varepsilon \to 0+$ we make use of Lebesgue's theorem on dominated convergence. \square

5.2. Existence of solutions

Let T' > 0 be given and $T \leq T'$. Assume a set of (fixed) data

$$(f_i, 0, g_i^{\text{in}}, 0, 0, 0, c_{0,i}, c_{0,i}^{\Sigma}) \in \mathbb{F}_{n,I}^{\Omega, \Sigma}(T')$$

to be given. We denote by

$$\mathcal{L}_{T,i}: \mathbb{E}_p^{\Omega}(T) \times \mathbb{E}_p^{\Sigma}(T) \to \mathbb{F}_p^{\Omega,\Sigma}(T)$$

the isomorphism induced by Proposition 4.1, that is, $\mathcal{L}_{T,i}$ is the full linear operator on the right-hand side of (4) (except for the time traces). The full nonlinear problem (1) then is reformulated as

$$\mathcal{L}_{T,i}(c_i, c_i^{\Sigma}) = (f_i, 0, g_i^{\text{in}}, 0, 0, 0) + \mathcal{N}_{T,i}(c, c^{\Sigma}),$$

$$c_i(0) = c_{0,i}, \quad c_i^{\Sigma}(0) = c_{i,0}^{\Sigma}, \quad i = 1, \dots, N,$$
(20)

where $\mathcal{N}_{T,i}$ includes the nonlinear sorption and reaction terms, i.e.,

$$\mathcal{N}_{T,i}(c,c^{\Sigma}) := \left(0, r_i^{\text{sorp}}(c_i, c_i^{\Sigma}) + r_{i,+}^{\text{ch}}(c^{\Sigma}), 0, r_i^{\text{sorp}}(c_i, c_i^{\Sigma}), 0, 0\right).$$

In order keep the constants resulting from the estimates below independent of T, we employ a suitable zero time trace splitting as described in the following.

First we take care of the compatibility condition arising from the nonlinear boundary condition on Σ . Taking time trace results in $r_i^{\text{sorp}}(c_{0,i},c_{0,i}^{\Sigma}) \in W_p^{1-3/p}(\Sigma)$, which will be extended to $\mathbb{G}_p^{\Sigma}(T)$ by setting

$$r_i^* := \mathcal{R}_{\Sigma} e^{t\Delta_{\Sigma(-\infty,\infty)}} \mathcal{E}_{\Sigma_{(-\infty,\infty)}} r_i^{\text{sorp}}(c_{0,i}, c_{0,i}^{\Sigma}).$$

Here $\mathcal{E}_{\Sigma_{(-\infty,\infty)}}$ denotes the extension operator from the lateral surface Σ to the surface of the infinite cylinder $\Sigma_{(-\infty,\infty)}$ and \mathcal{R}_{Σ} the corresponding restriction operator (note that both act as bounded operators on the function classes considered here, cf. [1]). Since $e^{t\Delta_{\Sigma_{(-\infty,\infty)}}}$ has the same regularizing properties as the Laplacian on the whole space \mathbb{R}^n , for which the desired regularity is well known [25], we see that $r_i^* \in \mathbb{G}_p^{\Sigma}(T)$.

Now we define the reference solution $(c_i^*, c_i^{\Sigma^*}) \in \mathbb{E}_p^{\Omega}(T') \times \mathbb{E}_p^{\Sigma}(T')$, existing according to Proposition 4.1, via

$$\mathcal{L}_{T,i}(c_i^*, c_i^{\Sigma^*}) = (f_i, 0, g_i^{\text{in}}, r_i^*, 0, 0), \quad c_i(0) = c_{0,i} \quad \text{in } \Omega, \quad c_i^{\Sigma}(0) = c_{0,i}^{\Sigma} \quad \text{on } \Sigma.$$
 (21)

Decomposing (c_i, c_i^{Σ}) as

$$c_i = \bar{c}_i + c_i^*, \quad c_i^{\Sigma} = \bar{c}_i^{\Sigma} + c_i^{\Sigma^*}$$

and subtracting (21) from (20), we end up with the reduced zero time trace problem

$${}_{0}\mathcal{L}_{T,i}(\bar{c},\bar{c}^{\Sigma}) = {}_{0}\mathcal{N}_{T,i}(\bar{c},\bar{c}^{\Sigma}) \qquad (i=1,...,N).$$

Here ${}_{0}\mathcal{L}_{T,i}$ denotes the restriction of $\mathcal{L}_{T,i}$ to ${}_{0}\mathbb{E}_{p}^{\Omega}(T) \times {}_{0}\mathbb{E}_{p}^{\Sigma}(T)$ and

$$_{0}\mathcal{N}_{T,i}(\bar{c},\bar{c}^{\Sigma}) := \mathcal{N}_{T,i}(\bar{c}_{i} + c_{i}^{*},\bar{c}_{i}^{\Sigma} + c_{i}^{\Sigma^{*}}) - (0,0,0,r_{i}^{*},0,0).$$

Next, we define $\Phi_T := (\Phi_{T,i})_{i=1,\ldots,N}$ through

$${}_{0}\Phi_{T}: {}_{0}\mathbb{E}_{p}^{\Omega}(T)^{N} \times {}_{0}\mathbb{E}_{p}^{\Sigma}(T)^{N} \to {}_{0}\mathbb{E}_{p}^{\Omega}(T)^{N} \times {}_{0}\mathbb{E}_{p}^{\Sigma}(T)^{N},$$

$${}_{0}\Phi_{T,i}(\bar{c},\bar{c}^{\Sigma}) := {}_{0}\mathcal{S}_{T,i} {}_{0}\mathcal{N}_{T,i}(\bar{c},\bar{c}^{\Sigma}), \qquad (i = 1, ..., N),$$

with the bounded linear inverse ${}_{0}\mathcal{S}_{T,i}$ of ${}_{0}\mathcal{L}_{T,i}$ given in Proposition 4.1.

Proof of Theorem 2.1. We apply the contraction mapping principle to ${}_0\Phi_T$, i.e., we show that there exists a $\delta > 0$, such that the mapping ${}_0\Phi_T$ constitutes a contraction on the closed ball $\bar{B}_{\delta}(0) \subset {}_0\mathbb{E}_p^{\Omega}(T)^N \times {}_0\mathbb{E}_p^{\Sigma}(T)^N$ and fulfills ${}_0\Phi_T : \bar{B}_{\delta}(0) \to \bar{B}_{\delta}(0)$.

(i) Contraction property: Let $(\bar{c}, \bar{c}^{\Sigma}), (\bar{z}, \bar{z}^{\Sigma}) \in \bar{B}_{\delta}(0)$. Then we have

$$\begin{split} &\|_{0}\Phi_{T}(\bar{c},\bar{c}^{\Sigma}) - {}_{0}\Phi_{T}(\bar{z},\bar{z}^{\Sigma})\|_{{}_{0}\mathbb{E}_{p}^{\Omega}(T)^{N}\!\times_{0}\mathbb{E}_{p}^{\Sigma}(T)^{N}} \\ &\leq C\|_{0}\mathcal{N}_{T}(\bar{c},\bar{c}^{\Sigma}) - {}_{0}\mathcal{N}_{T}(\bar{z},\bar{z}^{\Sigma})\|_{{}_{0}\mathbb{F}_{p}^{\Omega,\Sigma}(T)^{N}} \\ &= C\|r^{\text{sorp}}(\bar{c}+c^{*},\bar{c}^{\Sigma}+c^{\Sigma^{*}}) - r^{\text{sorp}}(\bar{z}+c^{*},\bar{z}^{\Sigma}+c^{\Sigma^{*}})\|_{(\mathbb{F}_{p}^{\Sigma}(T)\cap_{0}\mathbb{G}_{p}^{\Sigma}(T))^{N}} \\ &+ C\|r^{\text{ch}}((\bar{c}^{\Sigma}+c^{\Sigma^{*}})^{+}) - r^{\text{ch}}((\bar{z}^{\Sigma}+c^{\Sigma^{*}})^{+})\|_{\mathbb{F}_{n}^{\Sigma}(T)^{N}} \end{split} \tag{22}$$

with

$$C := \sup\{\|{}_{0}\mathcal{S}_{T}\|_{\mathscr{L}({}_{0}\mathbb{F}_{p}^{\Omega,\Sigma}(T)^{N}, {}_{0}\mathbb{E}_{p}^{\Omega}(T)^{N} \leftarrow 0\mathbb{E}_{p}^{\Sigma}(T)^{N})}: T \in (0,T']\}$$

independent of T, cf. Proposition 4.1. From Remark 5.1 we infer that

$$||r^{\text{ch}}((\bar{c}^{\Sigma} + c^{\Sigma^*})^+) - r^{\text{ch}}((\bar{z}^{\Sigma} + c^{\Sigma^*})^+)||_{L^p(\Sigma_T)^N} \le L||\bar{c}^{\Sigma} - \bar{z}^{\Sigma}||_{L^{p\gamma}(\Sigma_T)^N}$$

for a constant L > 0 depending on δ but not T and γ . Note in passing that we also used that $h \mapsto h^+$ is globally Lipschitz continuous from $L^{p\gamma}(\Sigma_T)$ to $L^{p\gamma}(\Sigma_T)$ with Lipschitz constant 1. By the fact that p > 5/3 we can estimate as

$$\|\bar{c}^{\Sigma} - \bar{z}^{\Sigma}\|_{L^{p\gamma}(\Sigma_T)^N} \le KT^{\eta} \|\bar{c}^{\Sigma} - \bar{z}^{\Sigma}\|_{0\mathbb{E}_p^{\Sigma}(T)^N}$$

with a constant K > 0 and an exponent $\eta > 0$ independent of T. We arrive at

$$||r^{\text{ch}}((\bar{c}^{\Sigma} + c^{\Sigma^*})^+) - r^{\text{ch}}((\bar{z}^{\Sigma} + c^{\Sigma^*})^+)||_{\mathbb{F}_p^{\Sigma}(T)^N} \le LKT^{\eta}||\bar{c}^{\Sigma} - \bar{z}^{\Sigma}||_{0}\mathbb{E}_p^{\Sigma}(T)^N.$$
(23)

We turn to the estimate of the sorption rate. By (A_F^{sorp}) , (A_B^{sorp}) the mapping r_i^{sorp} acts as a Nemytskij operator from $W_p^s(\Sigma_T) \times W_p^s(\Sigma_T)$ to $W_p^s(\Sigma_T)$ for $s \in (0,1)$, cf. Section 3.1 in [27]. Note that (A_B^{sorp}) implies $r_i^{sorp}(0,0) = 0$. An application of the mean value theorem to the function r_i^{sorp} and using $\nabla r_i^{sorp} \in BC^1(\mathbb{R}^2,\mathbb{R}^2)$ yields the global Lipschitz continuity of its induced Nemytskij operator. Hence we may employ

$$W_p^s(\Sigma_T) = W_p^s((0,T), L^p(\Sigma)) \cap L^p((0,T), W_p^s(\Sigma))$$

for $s \in (0,1)$ and

$${}_0\mathbb{H}^\Sigma_p(T)\times{}_0\mathbb{E}^\Sigma_p(T)\hookrightarrow \left({}_0W^{1-1/p+\varepsilon}_p((0,T),L^p(\Sigma))\cap L^p((0,T),W^{1-1/p+\varepsilon}_p(\Sigma))\right)^2$$

$$\xrightarrow{r_i^{\text{sorp}}}{}_0W^{1-1/p+\varepsilon}_p((0,T),L^p(\Sigma))\cap L^p((0,T),W^{1-1/p+\varepsilon}_p(\Sigma))\hookrightarrow{}_0\mathbb{G}^\Sigma_p(T),$$

for sufficiently small $\varepsilon > 0$, such that we obtain, similarly as for the estimate of the reaction term,

$$\|r_i^{\text{sorp}}(\bar{c} + c^*, \bar{c}^{\Sigma} + c^{\Sigma^*}) - r_i^{\text{sorp}}(\bar{z} + c^*, \bar{z}^{\Sigma} + c^{\Sigma^*})\|_{0\mathbb{G}_p^{\Sigma}(T)^N}$$

$$\leq L'K'T^{\eta}\left(\|\bar{c}_i - \bar{z}_i\|_{0\mathbb{E}_p^{\Omega}(T)} + \|\bar{c}_i^{\Sigma} - \bar{z}_i^{\Sigma}\|_{0\mathbb{E}_p^{\Sigma}(T)}\right)$$

$$(24)$$

with constants L', K' > 0 and an exponent $\eta > 0$ independent of T < T'. Combining (23) and (24) yields

for $(\bar{c}, \bar{c}^{\Sigma}), (\bar{z}, \bar{z}^{\Sigma}) \in \bar{B}_{\delta}(0)$. We choose T so small that

$$C(LK + L'K')T^{\eta} \le \frac{1}{2},\tag{25}$$

which is possible since all other constants appearing in (25) are independent of T < T'. Hence ${}_0\Phi_T$ is a contraction on $\bar{B}_{\delta}(0)$.

(ii) Self mapping property: Let $(\bar{c}, \bar{c}^{\Sigma}) \in \bar{B}_{\delta}(0)$. Then we have

$$\begin{split} &\|_{0}\Phi_{T}(\bar{c},\bar{c}^{\Sigma})\|_{{}_{0}\mathbb{E}_{p}^{\Omega}(T)^{N}\!\times_{0}\mathbb{E}_{p}^{\Sigma}(T)^{N}} \leq C\|_{0}\mathcal{N}_{T}(\bar{c}+c^{*},\bar{c}^{\Sigma}+c^{\Sigma^{*}})\|_{{}_{0}\mathbb{F}_{p}^{\Omega,\Sigma}(T)^{N}} \\ &\leq C\|r^{\mathrm{sorp}}(\bar{c}+c^{*},\bar{c}^{\Sigma}+c^{\Sigma^{*}})-r^{*}\|_{{}_{0}\mathbb{G}_{p}^{\Sigma}(T)^{N}}+C\|r^{\mathrm{ch}}((\bar{c}^{\Sigma}+c^{\Sigma^{*}})^{+})\|_{\mathbb{F}_{p}^{\Sigma}(T)^{N}}. \end{split}$$

Analogously to (i) we estimate the reaction term by

$$||r^{\operatorname{ch}}((\bar{c}^{\Sigma} + c^{\Sigma^*})^+)||_{L^p(\Sigma_T)^N}$$

$$\leq ||r^{\operatorname{ch}}((\bar{c}^{\Sigma} + c^{\Sigma^*})^+) - r^{\operatorname{ch}}((c^{\Sigma^*})^+)||_{L^p(\Sigma_T)^N} + ||r^{\operatorname{ch}}((c^{\Sigma^*})^+)||_{L^p(\Sigma_T)^N}$$

$$\leq L||\bar{c}^{\Sigma}||_{L^{p\gamma}(\Sigma_T)^N} + ||r^{\operatorname{ch}}((c^{\Sigma^*})^+)||_{L^p(\Sigma_T)^N}$$

$$\leq LKT^{\eta}||\bar{c}^{\Sigma}||_{{}_{0}\mathbb{E}^{\Sigma}_{2}(T)^N} + ||r^{\operatorname{ch}}((c^{\Sigma^*})^+)||_{L^p(\Sigma_T)^N}$$

$$(26)$$

with constants L, K > 0 being independent of T. In the same manner as in (i) for the sorption term we obtain

$$\begin{aligned} & \|r_{i}^{\text{sorp}}(\bar{c}_{i} + c_{i}^{*}, \bar{c}_{i}^{\Sigma} + c_{i}^{\Sigma^{*}}) - r_{i}^{*}\|_{0\mathbb{G}_{p}^{\Sigma}(T)} \\ & \leq \|r_{i}^{\text{sorp}}(\bar{c}_{i} + c_{i}^{*}, \bar{c}_{i}^{\Sigma} + c_{i}^{\Sigma^{*}}) - r_{i}^{\text{sorp}}(c_{i}^{*}, c_{i}^{\Sigma^{*}})\|_{0\mathbb{G}_{p}^{\Sigma}(T)} \\ & + \|r_{i}^{\text{sorp}}(c_{i}^{*}, c_{i}^{\Sigma^{*}})\|_{\mathbb{G}_{p}^{\Sigma}(T)} + \|r_{i}^{*}\|_{0\mathbb{G}_{p}^{\Sigma}(T)} \\ & \leq L'K'T^{\eta}\|(\bar{c}_{i}, \bar{c}_{i}^{\Sigma})\|_{0\mathbb{E}_{p}^{\Omega}(T)^{N_{\times 0}}\mathbb{E}_{p}^{\Sigma}(T)^{N}} + \|r_{i}^{\text{sorp}}(c_{i}^{*}, c_{i}^{\Sigma^{*}})\|_{\mathbb{G}_{p}^{\Sigma}(T)} + \|r_{i}^{*}\|_{0\mathbb{G}_{p}^{\Sigma}(T)} \end{aligned} \tag{27}$$

with constants L', K' > 0 and an exponent $\eta > 0$ being all independent of T. Putting together (26) and (27) yields

Since c_i^* , $c_i^{\Sigma^*}$, and r_i^* are fixed functions, notice that the latter three terms can be made small by choosing T > 0 small. Thus, by choosing T so small that the sum of those three terms is less than $\delta/2$ and such that (25) is satisfied, we arrive at

The proof of Theorem 2.1 is now complete. \Box

6. Global well-posedness

In this section we show that the local-in-time strong solutions obtained in Theorem 2.1 in fact exist globally, provided the reaction rates satisfy some structural condition that allows for the derivation of a priori estimates. We suppose in addition to (A_F^{ch}) , (A_N^{ch}) and (A_F^{ch}) the following assumption on the structure of the reaction rates.

(A_S^{ch}) There exists a lower triangular matrix $Q=(q_{ij})_{1\leq i,j\leq N}\in\mathbb{R}^{N\times N}$ with strictly positive diagonal entries such that

$$Qr^{\mathrm{ch}}(y) \le C\left(1 + \sum_{j=1}^{N} y_j\right)v, \qquad y \in [0, \infty)^N$$
(28)

for some constant C > 0 and v = (1, ..., 1).

Reaction—diffusion systems with this triangular condition have been widely studied by several authors; see [24] and the references cited therein. When proving global existence results, condition (28) allows for an iteration scheme which has been applied successfully in many situations. A major objective of this section is

to generalize this iteration scheme for standard reaction–diffusion systems subject to (A_S^{ch}) to heterogeneous catalysis. The main difference compared to standard systems lies in the fact that the reaction takes place on the boundary instead inside the bulk and that we also have to deal with terms arising from sorption processes.

In this setting, we obtain the global-in-time well-posedness result, given by Theorem 2.3. The proof of Theorem 2.3 requires, besides the maximal regularity estimates obtained in Section 4, also some comparison principles and some weak-type estimates, which are provided by the following results.

Lemma 6.1. Let T > 0 and let $1 . Let <math>\alpha, \beta > 0$.

a) Assume $f \in \mathbb{F}_p^{\Omega}(T)$, $g^{in} \in \mathbb{G}_p^{in}(T)$, $g^{\Sigma} \in \mathbb{G}_p^{\Sigma}(T)$ and $v_0 \in \mathbb{F}_p^{\Omega}$ with

$$(u \cdot \nu)v_0 - \beta \partial_{\nu}v_0 = g^{in}(0)$$
 on Γ_{in} , $-\beta \partial_{\nu}v_0 = g^{\Sigma}(0)$ on Σ , $\Omega_p^{\Omega}(T)$ be a strong solution to

if p > 3. Let $v \in \mathbb{E}_p^{\Omega}(T)$ be a strong solution to

$$\begin{cases} \alpha \partial_t v + (u \cdot \nabla)v - \beta \Delta v = f & in \ (0, T) \times \Omega, \\ (u \cdot \nu)v - \beta \partial_\nu v = g^{in} & on \ (0, T) \times \Gamma_{in}, \\ -\beta \partial_\nu v = g^{\Sigma} & on \ (0, T) \times \Sigma, \\ -\beta \partial_\nu v = 0 & on \ (0, T) \times \Gamma_{out}, \\ v(0) = v_0 & in \ \Omega. \end{cases}$$

If $f, v_0 > 0$ and $g^{in}, g^{\Sigma} < 0$, then v > 0.

b) Assume $f \in \mathbb{F}_p^{\Sigma}(T)$ and $v_0 \in \mathbb{F}_p^{\Sigma}$ with $-\beta \partial_{\nu_{\Sigma}} v_0 = 0$ on $\partial \Sigma$, if p > 3. Let $v \in \mathbb{F}_p^{\Sigma}(T)$ be a strong solution

$$\begin{cases} \alpha \partial_t v - \beta \Delta_{\Sigma} v = f & on (0, T) \times \Sigma, \\ -\beta \partial_{\nu_{\Sigma}} v = 0 & on (0, T) \times \partial \Sigma, \\ v(0) = v_0 & on \Sigma. \end{cases}$$

If $f, v_0 \geq 0$, then $v \geq 0$.

Proof. The proof follows the lines of the proof of Lemma 5.2, except that here we deal with a linear problem, only. \Box

Lemma 6.2. Let $T^*>0$ and let $2\leq p<\infty$ and $2\leq q\leq\infty$. Let $\alpha,\,\beta>0$. Assume $f\in\mathbb{F}_p^\Omega(T),\,g^{in}\in\mathbb{G}_p^{in}(T),$ $g^{\Sigma} \in \mathbb{G}_p^{\Sigma}(T) \text{ and } v_0 \in \mathbb{I}_p^{\Omega} \cap BC(\Omega) \text{ with }$

$$(u \cdot \nu)v_0 - \beta \partial_{\nu}v_0 = g^{in}(0)$$
 on Γ_{in} , $-\beta \partial_{\nu}v_0 = g^{\Sigma}(0)$ on Σ , T be a strong solution to

if p > 3. Let $v \in \mathbb{E}_p^{\Omega}(T)$ be a strong solution to

$$\begin{cases} \alpha \partial_t v + (u \cdot \nabla)v - \beta \Delta v = f & in \ (0, T) \times \Omega, \\ (u \cdot \nu)v - \beta \partial_\nu v = g^{in} & on \ (0, T) \times \Gamma_{in}, \\ -\beta \partial_\nu v = g^{\Sigma} & on \ (0, T) \times \Sigma, \\ -\beta \partial_\nu v = 0 & on \ (0, T) \times \Gamma_{out}, \\ v(0) = v_0 & in \ \Omega. \end{cases}$$

Then

$$||v||_{L^q(\Omega_T)} + ||v||_{L^q(\Sigma_T)} \le C \Big(||f||_{L^q(\Omega_T)} + ||g^{in}||_{L^q(\Gamma_{in,T})} + ||g^{\Sigma}||_{L^q(\Sigma_T)} + ||v_0||_{L^q(\Omega)} \Big)$$

for some constant C > 0 that is independent of p and $0 < T < T^*$.

Proof. If the right-hand side is infinite, nothing has to be proved. So, we now assume that the data admits finite L^q -norm. We distinguish two cases:

Case 1: p > 5/2, such that $\mathbb{E}_p^{\Omega}(T) \hookrightarrow BC(\overline{\Omega_T})$, cf. [2].

q=2: standard (multiply the equation by v to obtain L^2 -estimates).

 $q = \infty$: Let \mathcal{L}_T denote the operator given by the left-hand side of the system under consideration, such that $\mathcal{L}_T v = (f, g^{\text{in}}, g^{\Sigma}, 0, v_0)$. For fixed 5/2 < r < 3 let $\phi \in \mathbb{E}_r^{\Omega}(T)$ denote the solution of

$$\mathcal{L}_T \phi = (1, -1, -1, 0, 1) =: F.$$

Note that $\mathbb{E}_r^{\Omega}(T) \hookrightarrow BC(\overline{\Omega_T})$ and that no compatibility conditions for F occur. Let

$$\delta := \|f\|_{L^{\infty}(\Omega_T)} + \|g^{\mathrm{in}}\|_{L^{\infty}(\Gamma_{\mathrm{in},T})} + \|g^{\Sigma}\|_{L^{\infty}(\Sigma_T)} + \|v_0\|_{L^{\infty}(\Omega)}$$

and let $\bar{v} \in \mathbb{E}_p^{\Omega}(T)$ be given through $\mathcal{L}_T \bar{v} = \delta F$. Then the comparison principle (Lemma 6.1 a)) applied to $\bar{v} - v$ yields $v \leq \bar{v}$ on $\overline{\Omega_T}$. Since $\bar{v} = \delta \phi$ and $\|\bar{v}\|_{BC} = \delta \|\phi\|_{BC}$ it follows that

$$v \le \delta \left(\|\phi\|_{L^{\infty}(\Omega_T)} + \|\phi\|_{L^{\infty}(\Sigma_T)} \right) =: \delta C$$

on $\overline{\Omega_T}$. Analogously $-\delta C \leq v$, such that $\|v\|_{L^{\infty}(\Omega_T)} \leq C\delta$. Since in this case $v \in BC(\overline{\Omega_T})$, we trivially have

$$||v||_{L^{\infty}(\Sigma_T)} \le ||v||_{L^{\infty}(\Omega_T)} \le C\delta,$$

hence that

$$||v||_{L^{\infty}(\Sigma_T)} + ||v||_{L^{\infty}(\Omega_T)} \le C\delta.$$

For $2 \le q \le \infty$ we set

$$X_q := L^q(\Omega_T) \times L^q(\Sigma_T), \quad Y_q := L^q(\Omega_T) \times L^q(\Gamma_{\mathrm{in},T}) \times L^q(\Sigma_T) \times \{0\} \times L^q(\Omega).$$

Let S_T denote the system's solution operator. By the L^2 - and the L^{∞} -estimates obtained above we have

$$S_T \in \mathcal{L}(Y_2, X_2) \cap \mathcal{L}(Y_\infty, X_\infty).$$

By interpolation $S_T \in \mathcal{L}(Y_q, X_q)$ which yields the assertion for Case 1.

Case 2: $2 \le p \le 5/2$.

For convenience set $F:=(f,0,g^{\mathrm{in}},g^{\Sigma},0,0,v_0,0)\in\mathbb{F}_{p,I}^{\Omega,\Sigma}(T)$. Let r>5/2 and choose

$$F_n := (f_n, 0, g_n^{\text{in}}, g_n^{\Sigma}, 0, 0, v_{0,n}, 0) \in \mathbb{Y}_{r,q} := \mathbb{F}_{r,I}^{\Omega,\Sigma}(T) \cap (L^q(\Omega_T) \times \{0\} \times L^q(\Gamma_{\text{in},T}) \times L^q(\Sigma_T) \times \{0\} \times \{0\} \times L^q(\Omega) \times \{0\})$$

for $n \in \mathbb{N}$, such that $F_n \to F$ in $\mathbb{Y}_{p,q}$ (e.g. by extension of F to $\mathbb{R} \times \mathbb{R}^3$ respectively $\mathbb{R} \times (\partial A \times \mathbb{R})$ and mollification). Let $v_n \in \mathbb{E}_r^{\Omega}(T)$ denote the corresponding solution of $\mathcal{L}_T v_n = F_n$. Then Case 1 applies to F_n, v_n and there exists a C > 0 independent of $n \in \mathbb{N}$, such that

$$||v_n||_{L^q(\Omega_T)} + ||v_n||_{L^q(\Sigma_T)} \le C||F_n||_{L^q}.$$
(29)

Obviously $(v_n)_n$ is a Cauchy sequence in

$$\mathbb{X}_{p,q} := \mathbb{E}_p^{\Omega}(T) \cap (L^q(\Omega_T) \times L^q(\Sigma_T)),$$

such that we may pass to the limit $n \to \infty$ in (29). Hence we obtain $v_n \to v$ in $\mathbb{X}_{p,q}$ with v being the solution of $\mathcal{L}_T v = F$ and

$$||v||_{L^q(\Omega_T)} + ||v||_{L^q(\Sigma_T)} \le C||F||_{L^q}.$$

The next Lemma is standard for equations on standard domains Ω , cf. [24, Lemma 3.4]. Here we give a proof since we employ it on Σ .

Lemma 6.3. Let $T^* > 0$ and let $1 < p, q < \infty$. Let $\mu > 0$ and let arbitrary coefficients $\alpha_1, \ldots, \alpha_N, \beta_1, \ldots, \beta_N \in \mathbb{R}$ be given. Assume the data satisfies $f, g_1, \ldots, g_n \in \mathbb{F}_p^{\Sigma}(T)$ and $u_0, v_0^1, \ldots, v_0^N \in \mathbb{F}_p^{\Sigma} \cap BC(\Sigma)$ with $\partial_{\nu_{\Sigma}} u_0 = \partial_{\nu_{\Sigma}} v_0^j = 0$ on $\partial \Sigma$, if p > 3. Furthermore, let $u, v_1, \ldots, v_N \in \mathbb{F}_p^{\Sigma}(T)$ be strong solutions to

$$\begin{cases}
\partial_t u - \mu \Delta_{\Sigma} u = f, & \alpha_j \partial_t v_j - \beta_j \Delta_{\Sigma} v_j = g_j & on (0, T) \times \Sigma, \\
-\mu \partial_{\nu_{\Sigma}} u = 0, & -\beta_j \partial_{\nu_{\Sigma}} v_j = 0 & on (0, T) \times \partial \Sigma, \\
u(0) = u_0, & v_j(0) = v_0^j & on \Sigma
\end{cases}$$
(30)

for some $0 < T < T^*$. If $f \leq \sum_{j=1}^{N} g_j$, then

$$||u^+||_{L^q(\Sigma_T)} \le C \left(1 + \sum_{j=1}^N ||v_j||_{L^q(\Sigma_T)}\right)$$

for some constant $C = C(\|u_0\|_{BC(\Sigma)}, \|v_0^1\|_{BC(\Sigma)}, \dots, \|v_0^N\|_{BC(\Sigma)}) > 0$ that is independent of $0 < T < T^*$.

Proof. Let $\theta \in L^{q'}(\Sigma_T)^+$ with 1/q + 1/q' = 1 and let $0 < \tau < T$. By the transformation $t \mapsto \tau - t$ applied to the equation for u in (30) with data $f = \theta$, $u_0 = 0$ we obtain the backward heat equation,

$$\begin{cases}
-[\partial_t \phi + \mu \Delta_{\Sigma} \phi] = \theta(\tau - t) \text{ on } (0, \tau) \times \Sigma, \\
-\mu \partial_{\nu_{\Sigma}} \phi = 0 & \text{ on } (0, \tau) \times \partial \Sigma, \\
\phi(\tau) = 0 & \text{ on } \Sigma,
\end{cases}$$
(31)

which admits maximal $L^{q'}$ -regularity, cf. Section 4. In particular, we have $\phi \in BUC([0,\tau],L^{q'}(\Sigma))$ and

$$\sup_{s \in [0,\tau]} \|\phi(s)\|_{L^{q'}(\Sigma)} + \|\partial_t \phi\|_{L^{q'}(\Sigma_{\tau})} + \mu \|\Delta_{\Sigma} \phi\|_{L^{q'}(\Sigma_{\tau})} \le C \|\theta\|_{L^{q'}(\Sigma_{\tau})}$$

for a constant C > 0 which is independent of τ . Observe that $\phi \ge 0$, since $\theta \ge 0$. Hence by plugging in the first line of (31) and multiple partial integrations we have

$$\int_{\Sigma_{\tau}} u\theta d(t,\sigma) = \int_{\Sigma} u_0 \phi(0) d\sigma + \int_{\Sigma_{\tau}} \phi(\partial_t u - \mu \Delta_{\Sigma} u) d(t,\sigma)
\leq \int_{\Sigma} u_0 \phi(0) d\sigma + \sum_{j=1}^{N} \int_{\Sigma_{\tau}} \phi(\alpha_j \partial_t v_j - \beta_j \Delta_{\Sigma} v_j) d(t,\sigma)$$

$$= \int\limits_{\Sigma} (u_0 - \sum_{j=1}^N \alpha_j v_0^j) \phi(0) d\sigma + \sum_{j=1}^N \int\limits_{\Sigma_{\tau}} (-\alpha_j \partial_t \phi - \beta_j \Delta_{\Sigma} \phi) v_j d(t, \sigma)$$

$$\leq C \|\theta\|_{L^{q'}(\Sigma_{\tau})} \left(1 + \sum_{j=1}^N \|v_j\|_{L^q(\Sigma_{\tau})} \right).$$

Employing

$$\|u^+\|_{L^q(\Sigma_T)} = \sup \left\{ \int_0^T \int_{\Sigma} u\theta \, d\sigma(x) \, dt : \theta \in L^{q'}(\Sigma_T)^+, \|\theta\|_{L^{q'}(\Sigma_T)} < 1 \right\}$$

the assertion follows. \Box

Now, we are in position to prove Theorem 2.3.

Proof of Theorem 2.3. Since according to Theorem 2.1 the (local-in-time) solutions to (1) generate a local semi-flow in the phase space $\mathbb{I}_2^{\Omega} \times \mathbb{I}_2^{\Sigma} = H^1(\Omega) \times H^1(\Sigma)$, we may assume $T^* < \infty$ and show that the solution stays bounded in $H^1(\Omega) \times H^1(\Sigma)$ on $(0, T^*)$ in order to obtain a contradiction. It is sufficient to establish L^{∞} -bounds for the solution in order to obtain boundedness in the phase space. Then the H^1 -boundedness of solutions follows, as it is shown in the last part of this proof.

We will now derive L^{∞} -bounds, which requires several steps. Note that we may use the fact that c_i , $c_i^{\Sigma} \geq 0$ on $(0, T^*)$ thanks to Lemma 5.2.

Step 1. We have $c_i \in \mathbb{E}_2^{\Omega}(T) = H^1((0, T), L^2(\Omega)) \cap L^2((0, T), H^2(\Omega))$ and

$$\begin{cases}
\partial_t c_i + (u \cdot \nabla) c_i - d_i \Delta c_i = f_i & \text{in } (0, T) \times \Omega, \\
(u \cdot \nu) c_i - d_i \partial_\nu c_i = g_i^{\text{in}} & \text{on } (0, T) \times \Gamma_{\text{in}}, \\
-d_i \partial_\nu c_i = r_i^{\text{sorp}}(c_i, c_i^{\Sigma}) & \text{on } (0, T) \times \Sigma, \\
-d_i \partial_\nu c_i = 0 & \text{on } (0, T) \times \Gamma_{\text{out}}, \\
c_i(0) = c_{0,i} & \text{in } \Omega
\end{cases}$$
(32)

for $i=1,\ldots,N$ and all $0< T< T^*$. Now, $r_i^{\text{sorp}}(c_i,c_i^{\Sigma})\geq -k_i^{\text{de}}c_i^{\Sigma}$. Thus, by Lemma 6.1 a) we have $0\leq c_i\leq z_i$ for the unique maximal regular solution z_i to

$$\begin{cases} \partial_t z_i + (u \cdot \nabla) z_i - d_i \Delta z_i = f_i & \text{in } (0, T) \times \Omega, \\ (u \cdot \nu) z_i - d_i \partial_\nu z_i = g_i^{\text{in}} & \text{on } (0, T) \times \Gamma_{\text{in}}, \\ -d_i \partial_\nu z_i = -C c_i^{\Sigma} & \text{on } (0, T) \times \Sigma, \\ -d_i \partial_\nu z_i = 0 & \text{on } (0, T) \times \Gamma_{\text{out}}, \\ z_i(0) = c_{0,i} & \text{in } \Omega \end{cases}$$

with some appropriate constant $C = C((k_j^{\text{de}})_{j=1,\dots,N}) > 0$. Note that this problem allows for a strong solution in the L^2 -setting without any compatibility conditions between the right hand sides of the boundary conditions and the initial value. Since z_i is a solution to a linear problem, we may write $z_i = u_i + v_i$, where u_i solves

$$\begin{cases} \partial_t u_i + (u \cdot \nabla) u_i - d_i \Delta u_i = 0 & \text{in } (0, T) \times \Omega, \\ (u \cdot \nu) u_i - d_i \partial_\nu u_i = 0 & \text{on } (0, T) \times \Gamma_{\text{in}}, \\ -d_i \partial_\nu u_i = -C c_i^{\Sigma} & \text{on } (0, T) \times \Sigma, \\ -d_i \partial_\nu u_i = 0 & \text{on } (0, T) \times \Gamma_{\text{out}}, \\ u_i(0) = 0 & \text{in } \Omega \end{cases}$$

and v_i solves

$$\begin{cases} \partial_t v_i + (u \cdot \nabla) v_i - d_i \Delta v_i = f_i & \text{in } (0, T) \times \Omega, \\ (u \cdot \nu) v_i - d_i \partial_\nu v_i = g_i^{\text{in}} & \text{on } (0, T) \times \Gamma_{\text{in}}, \\ -d_i \partial_\nu v_i = 0 & \text{on } (0, T) \times \Sigma, \\ -d_i \partial_\nu v_i = 0 & \text{on } (0, T) \times \Gamma_{\text{out}}, \\ v_i(0) = c_{0,i} & \text{in } \Omega \end{cases}$$

For these solutions we have

$$||u_i||_{L^q(\Omega_T)} + ||u_i||_{L^q(\Sigma_T)} \le C' ||c_i^{\Sigma}||_{L^q(\Sigma_T)}$$

as well as

$$||v_i||_{L^q(\Omega_T)} + ||v_i||_{L^q(\Sigma_T)} \le A_i$$

provided that $2 \leq q \leq \infty$. Here, we employed Lemma 6.2 to obtain constants C' = C'(C, q) > 0 and $A_i = A_i(\|f_i\|_{L^q(\Omega_T)}, \|g_i^{\text{in}}\|_{L^q(\Gamma_{\text{in},T})}, \|c_{0,i}\|_{L^q(\Omega)}, q) > 0$ that are independent of $0 < T < T^*$. Since $\|c_i\|_{L^q} \leq \|c_i\|_{L^q} \leq \|c_i\|_{L^q} + \|c_i\|_{L^q}$ we may sum up the above estimates to obtain

$$\sum_{i=1}^{N} \|c_i\|_{L^q(\Omega_T)} + \sum_{i=1}^{N} \|c_i\|_{L^q(\Sigma_T)} \le C^* \left(1 + \sum_{j=1}^{N} \|c_j^{\Sigma}\|_{L^q(\Sigma_T)} \right), \tag{33}$$

where $C^* = C^*(C', (A_j)_{j=1,...,N}) > 0$ denotes a constant that is independent of $0 < T < T^*$. Note that this estimate is available for all $2 \le q \le \infty$.

Step 2. We have $c_i^{\Sigma} \in \mathbb{E}_2^{\Sigma}(T) = H^1((0,T),L^2(\Sigma)) \cap L^2((0,T),H^2(\Sigma))$ and

$$\begin{cases}
\partial_t c_i^{\Sigma} - d_i^{\Sigma} \Delta_{\Sigma} c_i^{\Sigma} = r_i^{\text{sorp}}(c_i, c_i^{\Sigma}) + r_i^{\text{ch}}(c^{\Sigma}) & \text{on } (0, T) \times \Sigma, \\
-d_i^{\Sigma} \partial_{\nu_{\Sigma}} c_i^{\Sigma} = 0 & \text{on } (0, T) \times \partial \Sigma, \\
c_i^{\Sigma}(0) = c_{0,i}^{\Sigma} & \text{on } \Sigma
\end{cases}$$
(34)

for i = 1, ..., N and all $0 < T < T^*$. Now we use the triangular structure of the reaction rates that is guaranteed by (A_S^{ch}) to treat the cases i = 1 and i = 2, ..., N separately.

Step 2.1. According to assumption (A_B^{sorp}) we have $r_1^{\text{sorp}}(c_1, c_1^{\Sigma}) \leq k_1^{\text{ad}}c_1$ and according to (A_S^{ch}) we have $q_{11}r_1^{\text{ch}}(c^{\Sigma}) \leq C(1 + \sum_{j=1}^N c_j^{\Sigma})$ for some fixed constant C > 0. Thus, by Lemma 6.1 b) we have $0 \leq c_1^{\Sigma} \leq z_1^{\Sigma}$ for the unique solution z_1^{Σ} to

$$\begin{cases} \partial_t z_1^\Sigma - d_1^\Sigma \Delta_\Sigma z_1^\Sigma = C' \Big(1 + c_1 + \sum_{j=1}^N c_j^\Sigma \Big) & \text{ on } (0, T) \times \Sigma, \\ -d_1^\Sigma \partial_{\nu_\Sigma} z_1^\Sigma = 0 & \text{ on } (0, T) \times \partial \Sigma, \\ z_1^\Sigma(0) = c_{0,1}^\Sigma & \text{ on } \Sigma \end{cases}$$

with some appropriate constant $C' = C'(C, k_1^{\text{ad}}, q_{11}) > 0$. Since z_1^{Σ} is a solution to a linear problem, we may write $z_1^{\Sigma} = u_1^{\Sigma} + v_1^{\Sigma}$, where u_1^{Σ} solves

$$\begin{cases} \partial_t u_1^\Sigma - d_1^\Sigma \Delta_\Sigma u_1^\Sigma = C' \Big(1 + c_1 + \sum_{j=1}^N c_j^\Sigma \Big) & \text{ on } (0, T) \times \Sigma, \\ -d_1^\Sigma \partial_{\nu_\Sigma} u_1^\Sigma = 0 & \text{ on } (0, T) \times \partial \Sigma, \\ u_1^\Sigma(0) = 0 & \text{ on } \Sigma \end{cases}$$

and v_1^{Σ} solves

$$\begin{cases} \partial_t v_1^\Sigma - d_1^\Sigma \Delta_\Sigma v_1^\Sigma = 0 & \text{ on } (0, T) \times \Sigma, \\ -d_1^\Sigma \partial_{\nu_\Sigma} v_1^\Sigma = 0 & \text{ on } (0, T) \times \partial \Sigma, \\ v_1^\Sigma(0) = c_{0,1}^\Sigma & \text{ on } \Sigma. \end{cases}$$

For these solutions we have

$$||u_1^{\Sigma}||_{L^q(\Sigma_T)} \le M||u_1^{\Sigma}||_{0\mathbb{E}_p^{\Sigma}(T)} \le LMC'\left(\Theta + ||c_1||_{L^p(\Sigma_T)} + \sum_{j=1}^N ||c_j^{\Sigma}||_{L^p(\Sigma_T)}\right)$$

with $\Theta = \Theta(p) = ||1||_{L^p(\Sigma_{T^*})}$ as well as

$$||(v_1^{\Sigma})^+||_{L^q(\Sigma_T)} \le A_1,$$

provided that $2 \leq p, \ q < \infty$. Here, $M = M(p, \ q) > 0$ denotes the norm of the embedding ${}_0\mathbb{E}_p^\Sigma(T) \hookrightarrow L^q(\Sigma_T)$, cf. [2] and L = L(p) > 0 denotes the norm of the solution operator in the L^p -setting for the time interval (0, T), which are both independent of $0 < T < T^*$ thanks to the homogeneous initial condition. Furthermore, $A_1 = A_1(\|c_{0,1}^\Sigma\|_{BC(\Sigma)}, q) > 0$ denotes the constant delivered by Lemma 6.3, which is also independent of $0 < T < T^*$. Observe, that a standard maximum principle could have been applied here, too. Note that $v_1^\Sigma \leq (v_1^\Sigma)^+$ implies $0 \leq c_1^\Sigma \leq z_1^\Sigma \leq u_1^\Sigma + (v_1^\Sigma)^+$. Therefore,

$$||c_1^{\Sigma}||_{L^q(\Sigma_T)} \leq ||z_1^{\Sigma}||_{L^q(\Sigma_T)} \leq ||u_1^{\Sigma}||_{L^q(\Sigma_T)} + ||(v_1^{\Sigma})^+||_{L^q(\Sigma_T)}$$

$$\leq C'' \left(1 + ||c_1||_{L^p(\Sigma_T)} + \sum_{j=1}^N ||c_j^{\Sigma}||_{L^p(\Sigma_T)}\right),$$

provided that $2 \le p, q < \infty$. Here $C'' = C''(LMC', \Theta, A_1) > 0$ is independent of $0 < T < T^*$.

Step 2.2. Now fix $i \in \{2, ..., N\}$. By (34) we obtain

$$q_{ii}(\partial_t c_i^{\Sigma} - d_i^{\Sigma} \Delta_{\Sigma} c_i^{\Sigma}) + \sum_{j < i} q_{ij} \left(\partial_t c_j^{\Sigma} - d_j^{\Sigma} \Delta_{\Sigma} c_j^{\Sigma} \right) = [Qr^{\text{sorp}}(c, c^{\Sigma})]_i + [Qr^{\text{ch}}(c^{\Sigma})]_i$$

on $(0, T) \times \Sigma$. According to assumption (A_B^{sorp}) we have

$$[Qr^{\mathrm{sorp}}(c,\,c^{\Sigma})]_i = \sum_{j \leq i} q_{ij} r_j^{\mathrm{sorp}}(c_j,\,c_j^{\Sigma}) \leq \sum_{j \leq i,q_{ij} > 0} q_{ij} k_j^{\mathrm{ad}} c_j - \sum_{j \leq i,q_{ij} < 0} q_{ij} k_j^{\mathrm{de}} c_j^{\Sigma}$$

and according to (A_S^{ch}) we have $[Qr^{\text{ch}}(c^{\Sigma})]_i \leq C(1+\sum_{j=1}^N c_j^{\Sigma})$ for some fixed constant C>0. Thus, by Lemma 6.1 b) we have $0\leq c_i^{\Sigma}\leq z_i^{\Sigma}$ for the unique maximal regular solution z_i^{Σ} to

$$\begin{cases} \partial_t z_i^\Sigma - d_i^\Sigma \Delta_\Sigma z_i^\Sigma = C' \Big(1 + \sum_{j \le i} c_j + \sum_{j=1}^N c_j^\Sigma \Big) \\ - \sum_{j < i} r_{ij} \Big(\partial_t c_j^\Sigma - d_j^\Sigma \Delta_\Sigma c_j^\Sigma \Big) & \text{in } (0, T) \times \Sigma, \\ - d_i^\Sigma \partial_{\nu_\Sigma} z_i^\Sigma = 0 & \text{on } (0, T) \times \partial \Sigma, \\ z_i^\Sigma (0) = c_{0,i}^\Sigma & \text{on } \Sigma \end{cases}$$

with some appropriate constant $C' = C'(C, k_j^{\text{ad}}, k_j^{\text{de}}, q_{ij}) > 0$ and $r_{ij} = q_{ij}/q_{ii}$. Since z_i^{Σ} is a solution to a linear problem, we may write $z_i^{\Sigma} = u_i^{\Sigma} + v_i^{\Sigma}$, where u_i^{Σ} solves

$$\begin{cases} \partial_t u_i^\Sigma - d_i^\Sigma \Delta_\Sigma u_i^\Sigma = C' \Big(1 + \sum_{j \leq i} c_j + \sum_{j=1}^N c_j^\Sigma \Big) & \text{in } (0, T) \times \Sigma, \\ - d_i^\Sigma \partial_{\nu_\Sigma} u_i^\Sigma = 0 & \text{on } (0, T) \times \partial \Sigma, \\ u_i^\Sigma(0) = 0 & \text{on } \Sigma \end{cases}$$

and v_i^{Σ} solves

$$\begin{cases} \partial_t v_i^{\scriptscriptstyle \Sigma} - d_i^{\scriptscriptstyle \Sigma} \Delta_{\scriptscriptstyle \Sigma} v_i^{\scriptscriptstyle \Sigma} = -\sum_{j < i} r_{ij} \big(\partial_t c_j^{\scriptscriptstyle \Sigma} - d_j^{\scriptscriptstyle \Sigma} \Delta_{\scriptscriptstyle \Sigma} c_j^{\scriptscriptstyle \Sigma} \big) & \text{in } (0, T) \times \Sigma, \\ - d_i^{\scriptscriptstyle \Sigma} \partial_{\nu_{\scriptscriptstyle \Sigma}} v_i^{\scriptscriptstyle \Sigma} = 0 & \text{on } (0, T) \times \partial \Sigma, \\ v_i^{\scriptscriptstyle \Sigma}(0) = c_{0,i}^{\scriptscriptstyle \Sigma} & \text{on } \Sigma. \end{cases}$$

For these solutions we have

$$\|u_i^{\Sigma}\|_{L^q(\Sigma_T)} \le M\|u_i^{\Sigma}\|_{0\mathbb{E}_p^{\Sigma}(T)} \le LMC'\left(\Theta + \sum_{j \le i} \|c_j\|_{L^p(\Sigma_T)} + \sum_{j=1}^N \|c_j^{\Sigma}\|_{L^p(\Sigma_T)}\right)$$

with $\Theta = \Theta(p) = ||1||_{L^p(\Sigma_{T^*})}$ as well as

$$\|(v_i^{\Sigma})^+\|_{L^q(\Sigma_T)} \le A_i \left(1 + \sum_{j < i} \|c_j^{\Sigma}\|_{L^q(\Sigma_T)}\right),$$

provided that $2 \leq p, q < \infty$. Here, M = M(p, q) > 0 denotes the norm of the embedding ${}_0\mathbb{E}_p^\Sigma(T) \hookrightarrow L^q(\Sigma_T)$ and L = L(p) > 0 denotes the norm of the solution operator in the L^p -setting for the time interval (0, T), which are both independent of $0 < T < T^*$ thanks to the homogeneous initial condition as in Step 2.1. Furthermore, $A_i = A_i((\|c_{0,j}^\Sigma\|_{BC(\Sigma)})_{j=1,\dots,i}, q) > 0$ denotes the constant delivered by Lemma 6.3, which is also independent of $0 < T < T^*$. We again have $c_i^\Sigma \leq z_i^\Sigma \leq u_i^\Sigma + (v_i^\Sigma)^+$. Therefore,

$$\begin{aligned} \|c_i^{\Sigma}\|_{L^q(\Sigma_T)} &\leq \|z_i^{\Sigma}\|_{L^q(\Sigma_T)} \leq \|u_i^{\Sigma}\|_{L^q(\Sigma_T)} + \|(v_i^{\Sigma})^+\|_{L^q(\Sigma_T)} \\ &\leq C'' \left(1 + \sum_{j \leq i} \|c_j\|_{L^p(\Sigma_T)} + \sum_{j=1}^N \|c_j^{\Sigma}\|_{L^p(\Sigma_T)}\right) + A_i \sum_{j < i} \|c_j^{\Sigma}\|_{L^q(\Sigma_T)}, \end{aligned}$$

provided that $2 \le p$, $q < \infty$. Here $C'' = C''(LMC', A_i, \Theta) > 0$ is independent of $0 < T < T^*$.

Step 2.3. Now we may combine the estimates obtained in Steps 2.1 and 2.2, recursively, and infer that

$$\sum_{i=1}^{N} \|c_i^{\Sigma}\|_{L^q(\Sigma_T)} \le C^* \left(1 + \sum_{j=1}^{N} \|c_j\|_{L^p(\Sigma_T)} + \sum_{j=1}^{N} \|c_j^{\Sigma}\|_{L^p(\Sigma_T)} \right)$$
(35)

provided that $2 \le p$, $q < \infty$. Here, $C^* = C^*(C'', (A_j)_{j=1,\dots,N}) > 0$ is independent of $0 < T < T^*$.

Step 3. Now we combine estimates (33) and (35) to obtain

$$\sum_{i=1}^{N} \|c_{i}\|_{L^{q}(\Omega_{T})} + \sum_{i=1}^{N} \|c_{i}\|_{L^{q}(\Sigma_{T})} + \sum_{i=1}^{N} \|c_{i}^{\Sigma}\|_{L^{q}(\Sigma_{T})}$$

$$\leq C^{*} \left(1 + \sum_{j=1}^{N} \|c_{j}\|_{L^{p}(\Omega_{T})} + \sum_{j=1}^{N} \|c_{j}\|_{L^{p}(\Sigma_{T})} + \sum_{j=1}^{N} \|c_{j}^{\Sigma}\|_{L^{p}(\Sigma_{T})} \right), \tag{36}$$

provided that $2 \le p$, $q < \infty$. Here, $C^* > 0$ is independent of $0 < T < T^*$. Using this inequality for p = 2 we may in particular obtain $L^q - L^2$ -estimates for arbitrary $2 \le q < \infty$.

Step 4.1. The surface concentrations satisfy

$$\begin{cases} \partial_t c_i^{\Sigma} - d_i^{\Sigma} \Delta_{\Sigma} c_i^{\Sigma} = f_i^{\Sigma} & \text{on } (0, T) \times \Sigma, \\ -d_i^{\Sigma} \partial_{\nu_{\Sigma}} c_i^{\Sigma} = 0 & \text{on } (0, T) \times \partial \Sigma, \\ c_i^{\Sigma}|_{t=0} = c_{0,i}^{\Sigma} & \text{on } \Sigma, \end{cases}$$

with $f_i^{\Sigma} := r_i^{\text{sorp}}(c_i, c_i^{\Sigma}) + r_i^{\text{ch}}(c^{\Sigma})$. Due to the polynomial growth of the nonlinearities we may estimate f^{Σ} in terms of c and c^{Σ} and employ (36) to obtain

$$\sum_{i=1}^{N} \|f_i^{\Sigma}\|_{L^r(\Sigma_T)} \le C^* \left(1 + \sum_{j=1}^{N} \|c_j\|_{L^p(\Omega_T)} + \sum_{j=1}^{N} \|c_j\|_{L^p(\Sigma_T)} + \sum_{j=1}^{N} \|c_j^{\Sigma}\|_{L^p(\Sigma_T)} \right),$$

where $2 \le p, r < \infty$ and $C^* > 0$ is independent of $0 < T < T^*$. Thus, for given $2 \le q \le \infty$ we may use this estimate for sufficiently large $2 \le r < \infty$ together with a classical result from [15], which yields the estimate

$$\sum_{i=1}^{N} \|c_i^{\Sigma}\|_{L^q(\Sigma_T)} \le C^* \left(1 + \sum_{j=1}^{N} \|c_j\|_{L^p(\Omega_T)} + \sum_{j=1}^{N} \|c_j\|_{L^p(\Sigma_T)} + \sum_{j=1}^{N} \|c_j^{\Sigma}\|_{L^p(\Sigma_T)} \right). \tag{37}$$

Let us note that [15, Theorem III.7.1] is stated for Dirichlet boundary conditions, but the result remains true in the Neumann case; see [6, Theorem 4], whose proof carries over to smooth manifolds as Σ . Note that in contrast to (35) obtained in the second step, the estimate (37) is available for all $2 \le p < \infty$ and all $2 \le q \le \infty$, while $C^* > 0$ is still independent of $0 < T < T^*$.

Step 4.2. Now we combine estimates (33) and (37) to obtain

$$\sum_{i=1}^{N} \|c_{i}\|_{L^{q}(\Omega_{T})} + \sum_{i=1}^{N} \|c_{i}\|_{L^{q}(\Sigma_{T})} + \sum_{i=1}^{N} \|c_{i}^{\Sigma}\|_{L^{q}(\Sigma_{T})}$$

$$\leq C^{*} \left(1 + \sum_{j=1}^{N} \|c_{j}\|_{L^{p}(\Omega_{T})} + \sum_{j=1}^{N} \|c_{j}\|_{L^{p}(\Sigma_{T})} + \sum_{j=1}^{N} \|c_{j}^{\Sigma}\|_{L^{p}(\Sigma_{T})} \right), \tag{38}$$

provided that $2 \le p < \infty$ and $2 \le q \le \infty$. Here, $C^* > 0$ is independent of $0 < T < T^*$. Using this inequality for p = 2 we may in particular obtain $L^q - L^2$ -estimates for arbitrary $2 \le q \le \infty$.

Step 5.1. The estimate (38) applied, for p=2 and $q=\infty$, yields

$$\begin{split} \sum_{i=1}^{N} \|c_i(t)\|_{L^2(\Omega)}^2 + \sum_{i=1}^{N} \|c_i(t)\|_{L^2(\Sigma)}^2 + \sum_{i=1}^{N} \|c_i^{\Sigma}(t)\|_{L^2(\Sigma)}^2 \\ & \leq C^* \left(1 + \sum_{j=1}^{N} \int\limits_0^t \|c_j(s)\|_{L^2(\Omega)}^2 \, \mathrm{d}s + \sum_{j=1}^{N} \int\limits_0^t \|c_j(s)\|_{L^2(\Sigma)}^2 \, \mathrm{d}s + \sum_{j=1}^{N} \int\limits_0^t \|c_j^{\Sigma}(s)\|_{L^2(\Sigma)}^2 \, \mathrm{d}s \right) \end{split}$$

for all $0 < t < T < T^*$, where $C^* > 0$ is independent of $0 < T < T^*$. Thus, a standard Gronwall argument implies

$$||c_i||_{L^2(\Omega_T)}, ||c_i||_{L^2(\Sigma_T)}, ||c_i^{\Sigma}||_{L^2(\Sigma_T)} \le Me^{\omega T}, \qquad 0 < T < T^*,$$
 (39)

with some constants $M, \omega > 0$, which are independent of $0 < T < T^*$.

Step 5.2. The estimate (38) again applied for p=2 and $q=\infty$ together with (39) yields

$$||c_i||_{L^{\infty}(\Omega_T)}, ||c_i||_{L^{\infty}(\Sigma_T)}, ||c_i^{\Sigma}||_{L^{\infty}(\Sigma_T)} \le Me^{\omega T}, \qquad 0 < T < T^*,$$
 (40)

with some constants $M, \omega > 0$, which are independent of $0 < T < T^*$.

Step 6. Now the obtained a priori estimates (40) carry over from L^{∞} to $H^1(\Omega)$ and $H^1(\Sigma)$. This may be seen by the following argument: Due to the L^{∞} -estimates, the L^2 -solution of (1) is contained in $\mathbb{E}_p^{\Omega}(T) \times \mathbb{E}_p^{\Sigma}(T)$ for each $1 with <math>p \neq 3$ by bootstrapping. Here the crucial estimate is

$$||r_i^{\text{sorp}}(c_i, c_i^{\Sigma})||_{\mathbb{G}_p^{\Sigma}(T)} \le C||(c_i, c_i^{\Sigma})||_{\mathbb{E}_q^{\Omega}(T) \times \mathbb{E}_q^{\Sigma}(T)}$$

with suitable q < p; see the proof of Theorem 2.1. This, in turn, yields

$$c_i \in BC(\overline{\Omega_T}), \qquad c_i^{\Sigma} \in BC(\overline{\Sigma_T}), \qquad (T < T^*).$$

Hence, by plugging in c, c^{Σ} into r^{sorp} , r^{ch} , we may consider (1) again, as a linear problem, this time for data being continuous in time. More precisely, we consider (4) for data

$$\begin{split} f_i &\in BC([0,T^*), L^2(\Omega)), \qquad f_i^{\Sigma} \in BC([0,T^*), L^2(\Sigma)), \\ g_i^{\text{in}} &\in BC([0,T^*), L^2(\Gamma_{\text{in}})), \quad g_i^{\Sigma} \in BC([0,T^*), L^2(\Sigma)), \\ g_i^{\text{out}} &\in BC([0,T^*), L^2(\Gamma_{\text{out}})), \end{split}$$

and

$$c_{0,i} \in H^1(\Omega), \qquad c_{0,i}^{\Sigma} \in H^1(\Sigma).$$

Following the strategy of the proof of our linear result in Section 4, in particular by transferring Lemma 4.4 to handle the inhomogeneous boundary data g_i^{in} , g_i^{out} and g_i^{Σ} and by employing a semigroup representation, we obtain that the unique solution of (1) satisfies

$$c_i \in BC([0, T^*), H^1(\Omega)), \qquad c_i^{\Sigma} \in BC([0, T^*), H^1(\Sigma)),$$

and the corresponding a priori estimates

$$\sup_{t \in [0,T]} (\|c_i(t)\|_{H^1(\Omega)} + \|c_i^{\Sigma}(t)\|_{H^1(\Sigma)}) \le M e^{\omega T} \sup_{t \in [0,T]} (\|f_i(t)\|_{L^2(\Omega)} + \|f_i^{\Sigma}(t)\|_{L^2(\Sigma)}
+ \|g_i^{\text{in}}(t)\|_{L^2(\Gamma_{\text{in}})} + \|g_i^{\Sigma}(t)\|_{L^2(\Sigma)} + \|g_i^{\text{out}}(t)\|_{L^2(\Gamma_{\text{out}})}), \qquad (T < T^*)$$
(41)

for constants $M, \omega > 0$ independent of T. Hence, we may pass to the limit $T \to T^*$ and see that both sides of (41) stay finite. The proof of Theorem 2.3 is now complete. \Box

Remark 6.4 (Sorption and reaction examples). A few remarks on the examples given in the introduction are in order here concerning the assumptions $(A_{\rm F}^{\rm sorp}), (A_{\rm M}^{\rm sorp}), (A_{\rm B}^{\rm sorp}), (A_{\rm F}^{\rm ch}), (A_{\rm N}^{\rm ch}), (A_{\rm P}^{\rm ch}), (A_{\rm S}^{\rm ch})$ stated in Section 5 and Section 6. Evidently, Henry's law (S1) satisfies all of our assumptions. However, Langmuir's law (S2) needs to be modified in order to meet all assumptions on the sorption. To this end we introduce ζ^+ , a smooth cut-off function, which approximates $(\cdot)^+$ pointwise and ζ^B a smooth and bounded function with bounded derivatives up to order 2 which is monotonically increasing. In addition suppose $\zeta^+(0) = 0$ and $\zeta^B(0) = 0$. Then we consider

$$\tilde{r}_{L,i}^{\text{sorp}}(c_i, c_i^{\Sigma}) = k_i^{\text{ad}} \zeta^{\text{B}}(c_i) \zeta^+ \left(1 - \frac{\zeta^+(c_i^{\Sigma})}{c_{\infty,i}^{\Sigma}} \right) - k_i^{\text{de}} c_i^{\Sigma},$$

which indeed satisfies $(A_{\rm F}^{\rm sorp}), (A_{\rm M}^{\rm sorp}), (A_{\rm B}^{\rm sorp})$ and therefore is covered by our main results. This modification is only necessary due to technical reasons, since

- $(A_{\mathrm{F}}^{\mathrm{sorp}})$ is violated due to $\nabla r_{L,i}^{\mathrm{sorp}} \notin BC^1(\mathbb{R}^2, \mathbb{R}^2)$.
- $(A_{\mathrm{M}}^{\mathrm{Forp}})$ is not guaranteed since $c_{i}^{\Sigma} \leq c_{\infty,i}^{\Sigma}$ is not postulated.
- For the same reason $(A_{\rm B}^{\rm sorp})$ is not satisfied since $(1-c_i^{\Sigma}/c_{\infty,i}^{\Sigma})$ could be negative.

For the time being it is not clear to the authors whether there are still global solutions in case we omit the cut-off functions. Observe that in our model there is no maximal capacity on the active surface, which is required in the original Langmuir law (S2) to gain nonnegativity of concentrations. Nonnegativity in turn is employed in the proof of the global existence result.

The reaction rate r^{ch} given in (R1) satisfies $(A_{\text{F}}^{\text{ch}})$, admits quadratic growth $(A_{\text{P}}^{\text{ch}})$, is quasi-positive $(A_{\text{N}}^{\text{ch}})$ and respects the triangular structure $(A_{\text{S}}^{\text{ch}})$ with corresponding matrix

$$Q = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 1 & 0 & 1 \end{pmatrix}.$$

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