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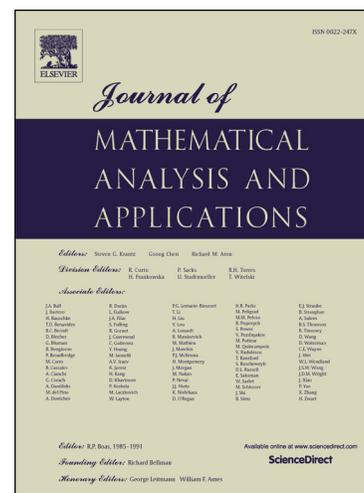
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THE BOLZANO MEAN-VALUE THEOREM AND PARTIAL DIFFERENTIAL EQUATIONS

WOJCIECH KRYSZEWSKI AND JAKUB SIEMIANOWSKI

ABSTRACT. We study the existence of solutions to abstract equations of the form $0 = Au + F(u)$, $u \in K \subset E$, where A is an abstract differential operator acting in a Banach space E , K is a closed convex set of constraints being invariant with respect to resolvents of A and a perturbation F satisfies a certain tangency condition. Such problems are closely related to the so-called Poincaré-Miranda theorem, being the multi-dimensional counterpart of the celebrated Bolzano intermediate value theorem. In fact our main results should be regarded as infinite-dimensional variants of Bolzano and Miranda-Poincaré theorems. Along with single-valued problems we deal with set-valued ones, yielding the existence of the so-called constrained equilibria of set-valued maps. The abstract results are applied to show existence of (strong) steady state solutions to some weakly coupled systems of drift reaction-diffusion equations or differential inclusions of this type. In particular we get the existence of strong solutions to the Dirichlet, Neumann and periodic boundary problems for elliptic partial differential inclusions under the presence of state constraints of different type. Certain aspects of the Bernstein theory for bvp for second order ODE are studied, too. No assumptions concerning structural coupling (monotonicity, cooperativity) are undertaken.

1. INTRODUCTION

The purpose of this paper is twofold. On the one hand, being motivated by some concrete applications (see subsection 2.1), we want to establish appropriate topological tools to study the existence of solutions to systems of N partial differential equations or N -dimensional partial differential inclusions subject to various boundary conditions and under state constraints. The presence of such constraints is justified and explained below. This is closely related to the method of the so-called ‘moving rectangles’ (see e.g. [50]) and corresponding techniques used for the study of long time behavior of evolution systems. We, however, leave aside questions concerning existence, stability and invariance of solutions of parabolic evolution equations, but in this paper we confine ourselves to elliptic equations rather and their solutions, i.e. steady state or stationary solutions to related evolutions problems. Nevertheless the ‘evolution’ origin of the studied steady state problems is of great importance.

On the other hand the proposed topological methods are closely related to problems of the existence of constrained equilibria or fixed points of abstract single- or set-valued maps, having origins in the Bolzano mean-value theorem (see Subsection 2.2). This celebrated result is perhaps the most important topological device when studying one-dimensional equations of the form $f(x) = 0$. This fact was extensively used and generalized by numerous authors for almost 150 years (see [39] and [12]) and various important results were established. One of the best known statements in this direction is

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the Poincaré-Miranda theorem, which is a direct N -dimensional version of the Bolzano theorem. We develop the infinite-dimensional counterparts of Poincaré-Miranda theorem, show their relation with different branches of research concerning e.g. viability theory for differential inclusions and, finally, apply them in the context of constrained PDE.

The notation used throughout the paper is standard. In particular $x \cdot y$ is the scalar product of $x, y \in \mathbb{R}^N$ and $|x| = \sqrt{x \cdot x}$ stands for the norm of x . The use of function spaces (L^p , Sobolev etc.), linear (unbounded in general) operators in Banach spaces, C_0 semigroups is standard. In the paper, for the sake of generality, we deal mostly with set-valued maps (the terminology in set-valued analysis is taken after [4]: the symbol \dashrightarrow denotes a set-valued map with at least closed values). It is however important to observe that results we propose are, to the best of our knowledge, new in the single-valued case, too.

The paper is organized as follows: in Section 2 we discuss origins of problems and motivations of main assumptions; in Section 3 we establish main abstract results, while Section 4 is devoted to applications. Subsection 3.2 contains a discussion of invariance issues playing an important role in the paper.

2. THE MOTIVATION

2.1. Drift reaction-diffusion equations. When dealing with an evolving in time multicomponent active continuous substance, whose components interact via certain reaction mechanism, such as e.g. predator-prey, activator-inhibitor, competition, reaction kinetics etc., and they all diffuse with different (in general) diffusive constants and are subject to drift or advection, i.e. a passive transfer caused by, for instance, the moving ambient media, such as gas or fluid, then the adequate model is provided by the so-called systems of drift reaction-diffusion equations (see e.g. [42]). Such systems in general are of the form

$$(2.1) \quad \partial_t u_i = \operatorname{div}(d_i \nabla u_i) + f_i(t, x, u, \nabla u), \quad i = 1, \dots, N,$$

or, shortly,

$$u_t = Lu + f(t, x, u, \nabla u), \quad \text{where } Lu = (v_1, \dots, v_N), \quad v_i = \operatorname{div}(d_i \nabla u_i) \quad \text{for } i = 1, \dots, N,$$

along with initial and the (Neumann or Dirichlet) boundary conditions on $\partial\Omega$. Here the unknown $u = (u_1, \dots, u_N)$ depends on spatial variables $x = (x_1, \dots, x_M) \in \Omega \subset \mathbb{R}^M$, Ω is an open smooth domain, and time $t \in [0, T]$, $T > 0$, $\nabla u = \left[\frac{\partial u_i}{\partial x_j} \right]$ is the derivative of u . The diffusive coefficients $d_i \in C^1(\bar{\Omega})$, $i = 1, \dots, N$, and $f = (f_1, \dots, f_N)$ is the source/advection term depending on $(t, x, u, \nabla u)$. Sometimes it is convenient to distinguish the advection term

$$(2.2) \quad f_i(t, x, u, \nabla u) = g_i(t, x, u) - \gamma^i \cdot \nabla u_i, \quad i = 1, \dots, N,$$

or, shortly,

$$f(t, x, u, \nabla u) = g(t, x, u) - \Gamma u, \quad \text{where } \Gamma u = (v_1, \dots, v_N), \quad v_i = \gamma^i \cdot \nabla u_i \quad \text{for } i = 1, \dots, N,$$

where drift vectors $\gamma^i := [\gamma_1^i, \dots, \gamma_M^i] \in L^\infty(\Omega, \mathbb{R}^M)$, $i = 1, \dots, N$; the second summand is responsible for the drift of the system.

Our interest is mainly focused on ecological or chemical systems, where $u_i(x, t)$ is the concentration at $x \in \Omega$ and time $t \in [0, T]$ of the i -th reactant, $i = 1, \dots, N$, contained in a bounded stirred up

vessel (or reactor) Ω . Clearly the initial state $u(\cdot, 0) \geq 0$ on Ω and the natural expectation is that $u_i(x, t) \geq 0$ since the concentration cannot be negative. On the other hand there is a threshold value $R_i > 0$ beyond which the i -th component is saturated or the model is not adequate. In a similar manner the implicit threshold value of concentrations may follow from mass conservation: the total mass of reactants, say R , must be constant. Therefore it makes sense to look for solutions $u(x, t)$ in the rectangle $\{u \in \mathbb{R}^N \mid 0 \leq u_i \leq R_i, i = 1, \dots, N\}$ or in the simplex $\{u \in \mathbb{R}^N \mid \sum_{i=1}^N u_i = R, u_i \geq 0, i = 1, \dots, N\}$. In general equations of the form (2.1) should be therefore considered under the presence of state constraints: $u(x, t) \in C$ for $x \in \Omega, t \in [0, T]$, where C is a given closed subset of the phase space \mathbb{R}^N .

This requirement $u(x, t) \in C$ has an impact on the source term f or g . If, for instance, we assume the mass conservation, i.e. demand that solution stay in the simplex $C = \{u \in \mathbb{R}^N \mid \sum_{i=1}^N u_i = R, u_i \geq 0, i = 1, \dots, N\}$, then the natural hypotheses onto g is that $\sum_{i=1}^N g_i(t, x, u) = 0$ (the total mass neither increases nor decreases during the process) and if $u_i = 0$, then $g_i(t, x, u) \geq 0$ (if a component vanishes, it can only be recreated). When carefully studying this and other examples, we gather that under the presence of state constraints in a closed convex $C \subset \mathbb{R}^N$ the correct restriction onto the source term is the *tangency condition*: $f(t, x, u) \in T_C(u)$ for $t \in [0, T], x \in \Omega$ and $u \in C$, where $T_C(u)$ stands for the cone tangent to C at u (see (2.7) and Remark 2.4).

In what follows we admit also discontinuous nonlinearities f (or g). This appears when, for instance, system data are determined by measurements or are subject to phase transition phenomena and is motivated by numerous applications of systems with hysteresis (see e.g. [52], [11]). The typical situation concerns (2.1) with $N = M = 1$ and is of the form

$$u_t = u_{xx} + H(u),$$

where H is the hysteresis operator – see [54], [34]. In the simplest case H is driven by the Heaviside function and may be described via the related Nemytskii operator: given a threshold value $\alpha > 0$, an input function $u : [0, T] \rightarrow \mathbb{R}$ with $u(0) \leq \alpha$, then the output $H(u)(t) = 0$ if $u(s) < \alpha$ for all $s \in [0, t]$ and $H(u)(t) = 1$ if $u(s) = \alpha$ for some $s \in [0, t]$. For some other instances of the problem – see [8], [16] and numerous examples in [10]. The common way to overcome this obstacle is to replace the discontinuous f , or g in (2.2), by an appropriate set-valued regularization F or G (introduced e.g. by Fillipov or Krasovski – see [27, Sect. 2.7] or [3, p. 101]) and instead of (2.1) consider a problem

$$(2.3) \quad u_t \in Lu + F(t, x, u, \nabla u)$$

or

$$(2.4) \quad u_t \in Lu - \Gamma u + G(t, x, u)$$

subject to initial and boundary conditions, where $F : [0, T] \times \Omega \times C \times \mathbb{R}^{MN} \rightarrow \mathbb{R}^N$ (or $G : [0, T] \times \Omega \times C \rightarrow \mathbb{R}^N$) is an upper semicontinuous set-valued map with compact convex values satisfying the appropriate tangency conditions with respect to the constraint set $C \subset \mathbb{R}^N$.

2.2. Zeros of set-valued maps. In the present paper we shall deal with the existence of steady state (stationary solutions) of state constrained autonomous problems related to (2.3) or (2.4). This leads to the second objective of the present paper.

Since 1941, when Kakutani showed that every upper semicontinuous set-valued self map φ of the closed ball D in \mathbb{R}^n admitting closed convex values has a fixed point, a lot of attention has been paid

to the different aspects of the fixed point theory for set-valued maps (see e.g. [29]). In one direction the development has led to substantial weakening in the assumption that the values of the mapping are subsets of its domain. The idea is well-illustrated by the classical (single-valued) mean value theorem of Bolzano.

Theorem 2.1. *If $f : C := [a, b] \rightarrow \mathbb{R}$ is continuous, $f(a)f(b) \leq 0$ (for instance $f(a) \geq 0$ and $f(b) \leq 0$), then there is $\bar{x} \in C$ such that $f(\bar{x}) = 0$.*

This important observation has been generalized by Poincaré in 1883 in his famous conjecture proved by Miranda [43].

Theorem 2.2. (Poincaré-Miranda) *Let $C = \Pi_{k=1}^n [a_k, b_k]$ be an n -dimensional cube and let $F_k^- := \{x \in C \mid x_k = a_k\}$, $F_k^+ := \{x \in C \mid x_k = b_k\}$, $k = 1, 2, \dots, n$, denote the k -th face of C . Let $f = (f_1, \dots, f_n) : C \rightarrow \mathbb{R}^n$ be continuous and suppose that for all $k = 1, \dots, n$*

$$(2.5) \quad f_k(x) \begin{cases} \geq 0 & \text{for every } x \in F_k^- \\ \leq 0 & \text{for every } x \in F_k^+. \end{cases}$$

Then f has a zero, i.e., there is $\bar{x} \in C$ such that $f(\bar{x}) = 0$.

Quite a complicated history of this and other related results is well-described by Mawhin [39] and [12] (see also [40], [41]). In the spirit of the above we have (see [41], [48])

Theorem 2.3. *Let $C = \{x \in \ell^2 \mid |x_k| \leq k^{-1}\}$ be the Hilbert cube, let $f : C \rightarrow \ell^2$ be continuous and such that for all $k \in \mathbb{N}$*

$$(2.6) \quad f_k(x_1, \dots, x_{k-1}, -\frac{1}{k}, x_{k+1}, \dots) \geq 0, \quad f_k(x_1, \dots, x_{k-1}, \frac{1}{k}, x_{k+1}, \dots) \leq 0,$$

then f has a zero.

In order to understand the nature of assumptions of these results we need to recall the following definition. Let E be a Banach space, $K \subset E$ be a closed set and $x \in K$. The contingent (or Bouligand) cone $T_K(x)$ is defined by

$$T_K(x) = \{v \in E \mid \liminf_{h \rightarrow 0^+} d_K(x + hv)/h = 0\}$$

and the Clarke tangent cone is defined by

$$C_K(x) = \{v \in E \mid \lim_{h \rightarrow 0^+, y \rightarrow x, y \in K} d_K(y + hv)/h = 0\},$$

where $d_K(u) := \inf_{y \in K} \|y - u\|$. $T_K(x)$ and $C_K(x)$ are closed cones; additionally $C_K(x)$ is convex. In general $C_K(x) \subset T_K(x)$ and if K is convex, then

$$(2.7) \quad T_K(x) = C_K(x) = \text{cl } S_K(x)$$

where

$$S_K(x) := \bigcup_{h > 0} h(K - x).$$

Observe that if x belongs to the interior of K , then $T_K(x) = E$. For examples and a detailed discussion see [4].

Remark 2.4. (1) If $C \subset \mathbb{R}^n$ is a cube (as in Theorem 2.2) and $x \in C$ is a boundary point, then $v \in T_C(x)$ if and only if $v_k \geq 0$ when $x_k = a_k$ and $v_k \leq 0$ when $x_k = b_k$. This implies that assumption (2.5) holds true if and only if $f(x) \in T_C(x)$, $x \in C$. Similarly, in the context of Theorem 2.3, assumption (2.6) is satisfied if and only if $f(x) \in T_C(x)$ for all $x \in C$.

(2) Let $C = \{x \in \mathbb{R}^n \mid \sum_{i=1}^n x_i = R, x_i \geq 0, i = 1, \dots, n\}$ and $x \in C$. If $x_i > 0$ for all i , then $v \in T_C(x)$ if and only if $\sum_{i=1}^n v_i = 0$; if $x_j = 0$, then $v \in T_C(x)$ if and only if $\sum_{i=1}^n v_i = 0$ and $v_j \geq 0$.

(3) If $K = D(0, r)$ is a closed ball in a Hilbert space E , then $T_K(x) = \{v \in E \mid \langle x, v \rangle \leq 0\}$ for x in the boundary of K .

In the apparently independent stream of research, the best known equilibrium result is the following pioneering result of Browder [13] (with some modification due to Halpern and Bergman [31, 32]) being, in the opinion of Aubin and Cellina (see [3, p. 213, Chapter 5.2] and the discussion therein), ‘one of the most powerful theorems of nonlinear analysis’.

Theorem 2.5. *Assume that $K \subset E$ is compact convex and $\varphi : K \multimap E$ is upper semicontinuous with closed convex values. If φ satisfies the **weak tangency** condition with respect to K , i.e.*

$$(2.8) \quad \forall x \in K \quad \varphi(x) \cap T_K(x) \neq \emptyset,$$

*then φ has an equilibrium: there is $\bar{x} \in K$ such that $0 \in \varphi(\bar{x})$. If φ satisfies the so-called the **weak inwardness (or outwardness)** condition, i.e.*

$$(2.9) \quad \forall x \in K \quad \varphi(x) \cap (x + T_K(x)) \neq \emptyset$$

(or $\varphi(x) \cap (x - T_K(x)) \neq \emptyset$ for $x \in K$), then φ has a fixed point.

In view of Remark 2.4 it is evident that Theorem 2.5 provides a far reaching generalization of Theorems 2.2 and 2.3.

Remark 2.6. It is easy to see that if K is convex and $0 \in K$, then $T_K(x) \subset x + T_K(x)$ for all $x \in K$. Hence in Theorem 2.3 (and in 2.2 if $a_k \leq 0 \leq b_k$, $k = 1, \dots, n$) f has a fixed point since $f(x) \in x + T_C(x)$, $x \in C$ and, therefore $f(x) - x = 0$ for some $x \in C$. Similarly in Theorem 2.5 if $0 \in K$, then (2.8) implies (2.9) and φ has fixed points.

Two drawbacks of this result has to be pointed out. In order to get a decent tool to study existence of equilibria one needs to get rid of compactness and convexity in Theorem 2.5. The best known result in this direction is due to Deimling – see [22, Th.11.5], [24].

Theorem 2.7. *Let K be a closed bounded convex subset of a Banach space E and let an upper semicontinuous map $\varphi : K \multimap E$ with compact convex values be condensing with respect to the Kuratowski or Hausdorff measure of noncompactness. If φ is weakly inward, then φ has a fixed point.*

It is interesting to observe that in Deimling’s theorem there is no way to replace inwardness by outwardness condition although it was possible in Theorem 2.5.

In order to discuss a nonconvex version of Theorem 2.5 one needs to understand which property of a set is a suitable substitute for convexity and what should be a suitable counterpart of tangency. This problem was addressed in [7] and discussed in [35].

After [7] we say that a closed $K \subset E$ is an \mathcal{L} -retract if there is $\varepsilon > 0$, a continuous $r : U = B(K, \eta) \rightarrow K$, where $B(K, \eta) := \{x \in E \mid d_K(x) < \eta\}$ ⁽¹⁾ and a constant $L \geq 1$ such that

$$r(x) = x \text{ for } x \in K \text{ and } \|r(x) - x\| \leq Ld_K(x) \text{ for } x \in U.$$

Therefore K is an \mathcal{L} -retract whenever K is a neighborhood retract in E with retraction r such that distance of $x \in U$ from $r(x) \in K$ may be controlled by the distance $d_K(x)$. The class of \mathcal{L} -retracts is large. Closed convex sets (in this case one can define r on E with $L = 1 + \varepsilon$, where $\varepsilon > 0$ is arbitrary), compact sets being bi-Lipschitz homeomorphic with closed convex sets, the so-called proximate retracts, Lipschitz retracts and epi-Lipschitz sets (in the sense of Rockafellar [46]) are \mathcal{L} -retracts.

Remark 2.8. (1) If X is a topological space of finite type, i.e., such that the (singular with rational coefficients) cohomology groups $H^k(X; \mathbb{Q})$, $k \geq 0$, groups are finitely generated and vanish above some dimension, then the Euler characteristic $\chi(X) := \sum_{k=0}^{\infty} (-1)^k \dim H^k(X, \mathbb{Q})$ of X is well-defined.

(2) If X is a neighborhood retract in E and $f : X \rightarrow X$ is compact, then f is a Lefschetz map, i.e. the homomorphism $H^*(f)$ is a Leray endomorphism of $H^*(X, \mathbb{Q})$ and the generalized Lefschetz number $\Lambda(f)$ of f is well-defined – see [25, Def. V.(2.1), (3.1). Th. (5.1)]. If f is homotopic to the identity I_X on X , then $H^*(f) = H^*(I_X)$ is the identity $H^*(X)$. This implies that I_X is a Lefschetz map; hence $H^*(X)$ is of finite type and the Euler characteristic $\chi(X)$ is well-defined. Moreover, in this case $\Lambda(f)$ is equal to the ordinary Lefschetz number $\lambda(f) = \lambda(I_X) = \chi(X)$ (for details concerning these notions see also e.g. [14]). In particular if $\chi(X) \neq 0$, then f has a fixed point.

In view of the above if K is a compact \mathcal{L} -retract, then its Euler characteristic $\chi(K)$ is well-defined. Note that if K is additionally convex, then $\chi(K) = 1$. After [7] (see also [19]) we have the following result.

Theorem 2.9. *Let $K \subset E$ be a compact \mathcal{L} -retract with $\chi(K) \neq 0$. If $\varphi : K \multimap E$ is upper semicontinuous with closed convex values and weakly tangent to K in the sense of Clarke, i.e.*

$$(2.10) \quad \forall x \in K \quad \varphi(x) \cap C_K(x) \neq \emptyset,$$

then φ has an equilibrium.

Note that in condition (2.10) the Bouligand cone has been replaced by the Clarke cone; there are examples showing that (2.8) is not sufficient (see [35]); however if $\varphi = f$ is a single-valued map, then (2.8) implies (2.10). It is also evident that the weak inwardness in the sense of Clarke cone implies the existence of fixed points.

There is no direct generalization of the equilibrium problem from Theorem 2.9 in the noncompact setting, although there were some partial answers have been discussed in [19] and [20], since we have the following example showing a compact tangent map without zeros.

Example 2.10. Let $E = \ell^2$ be the classical Hilbert space. Let $D_1 = \{x \in E \mid \|x\| \leq 1/\sqrt{2}\}$ and let $f : D_1 \rightarrow D_1$ be the variant of the famous Kakutani map: $f(x) = (\sqrt{1/2 - \|x\|^2}, x_1, x_2, \dots)$ for $x = (x_i) \in D_1$. Then f is continuous, it has neither zeros nor fixed points and $\|f(x)\| = 1/\sqrt{2}$ for all

¹Clearly $B(K, \eta) = K + \eta B$, where $B = B(0, 1)$ is the open unit ball in E .

$x \in D_1$. Let D be the unit closed ball in E and $r : D \rightarrow D_1$ the radial retraction. Define $g : D \rightarrow D$ by

$$g(x) := \begin{cases} -f(x) & \text{for } x \in D_1; \\ (2\|x\|^2 - 2)f(r(x)) + (1 - 2\|x\|^2)x & \text{for } x \in D \setminus D_1. \end{cases}$$

One can see that g is well-defined, continuous and $g(x) = -x$ whenever $x \in \partial D$; an easy argument yields $g(x) \neq 0$ for every $x \in D$. Finally, define $\kappa : E \rightarrow E$ by

$$\kappa(x_1, x_2, \dots, x_n, \dots) := (x_1, 2^{-1}x_2, \dots, n^{-1}x_n, \dots).$$

Clearly κ is an injective compact linear map. Thus $G := \kappa \circ g : D \rightarrow D$ is compact and $G(x) \neq 0$ for $x \in D$. If $x \in \partial D$, then $\langle G(x), x \rangle = -\sum_{n=1}^{\infty} \frac{1}{n} x_n^2 \leq 0$; by Remark 2.4 (3), G is tangent to D .

For examples, further generalizations and a deeper discussion of issues surveyed above the reader can see [35].

The main aim of the present paper is to show a result in this direction with applications to constrained steady state problems related to (2.3) or (2.4).

3. EXISTENCE RESULTS

3.1. The setting and results. In order to study the existence of steady states of autonomous problem (2.3) we shall take an appropriate abstract setting and consider the following coincidence problem

$$(3.1) \quad 0 \in Au + \Phi(u), \quad u \in K \subset E,$$

where:

- (A₁) $(E, \|\cdot\|)$ is a Banach space, $K \subset E$ is a closed convex set;
- (A₂) $A : D(A) \rightarrow E$ is a densely defined linear operator such that, for some $\omega \in \mathbb{R}$, $(\omega, \infty) \subset \rho(A)$;
- (A₃) there is another Banach space $(E_0, \|\cdot\|_0)$ and a closed convex $K_0 \subset E_0$ such that $D(A) \subset E_0 \subset E$, $K_0 \subset K$ and the identities $D(A) \rightarrow E_0$, $j : E_0 \rightarrow E$ are continuous ⁽²⁾;
- (A₄) $\Phi : K_0 \rightarrow E$ is H -upper semicontinuous, bounded and has convex weakly compact values;
- (A₅) for all $u \in K_0$, $\Phi(u) \cap T_K(j(u)) \neq \emptyset$.

Remark 3.1. (1) Let us recall the so-called Lions construction, see [21], [2]. Let V be a reflexive Banach space which is dense in a (real) Hilbert space H and suppose that the identity $V \rightarrow H$ is continuous. Suppose a bilinear continuous form $a : V \times V \rightarrow \mathbb{R}$ is such that

$$(3.2) \quad \forall v \in V \quad a(v, v) + \omega \|v\|_H^2 \geq \alpha \|v\|_V^2,$$

where $\omega \in \mathbb{R}$ and $\alpha > 0$. Let $D(A)$ be the set of all $u \in V$ such that the function $V \ni v \mapsto a(u, v)$ is continuous on V with the H -norm, i.e. such that there is $\beta_u \geq 0$ with $|a(u, v)| \leq \beta_u \|v\|_H$ for all $v \in V$. For any $u \in D(A)$ there is a unique $Au \in H$ such that

$$a(u, v) = -\langle Au, v \rangle_H, \quad v \in V.$$

²On $D(A)$ the graph norm $\|\cdot\|_A$ is considered:

$$\|u\|_A := \|u\| + \|Au\|, \quad u \in D(A).$$

This defines a linear $A : D(A) \rightarrow H$. According to e.g. [49, Prop. 4.1], $D(A)$ is dense in H , A is closed and $(\omega, +\infty) \subset \rho(A)$. Moreover, A is maximal ω -dissipative and, in view of [2, Theorem 4.2], A is the generator of a holomorphic C_0 semigroup $\{S(t)\}_{t \geq 0}$ of linear operators on H with growth bound equal to ω . Putting $E := H$, $E_0 := V$ and taking a closed convex bounded $K \subset E$ we see that assumption $(A_1) - (A_3)$ are satisfied provided $K_0 := K \cap E_0$. The situation described in this example is very typical in various applications.

(2) Assumption (A_4) is motivated by applications. H -upper semicontinuity (where H stands for ‘Hausdorff’) in (A_4) means that for each $x \in K_0$ and $\varepsilon > 0$ there is $\delta > 0$ such that $\Phi(y) \in B(\Phi(x), \varepsilon)$ if $y \in K_0$ and $\|y - x\|_0 < \delta$. It is well known that (see e.g. [10, prop. 2.3]) that Φ , being H -upper semicontinuous with weakly compact values, is upper semicontinuous when E is endowed with the weak topology. This, in turn, implies, that given sequences $(x_n) \subset K_0$ and $y_n \in \Phi(x_n)$, if $x_n \rightarrow x_0 \in K_0$, there is a subsequence y_{n_k} such that $y_{n_k} \rightarrow y_0 \in \Phi(x_0)$ (\rightarrow denotes the weak convergence). Obviously if Φ is upper semicontinuous with closed compact values, then (A_4) is satisfied, too.

(3) Let $h > 0$ and $h\omega < 1$, then $h^{-1} \in \rho(A)$ and the resolvent

$$J_h := (I - hA)^{-1} : E \rightarrow E$$

is well defined and continuous. Observe that $J_h(u) \in D(A) \subset E_0$ for any $u \in E$. Moreover for any $u \in E$

$$\|J_h(u)\|_A = \|J_h(u)\| + \|AJ_h(u)\| \leq ((1 + h^{-1})\|J_h\| + h^{-1})\|u\|,$$

This, together with (A_3) , shows that $J_h : E \rightarrow E_0$ is continuous, too. Observe that J_h is compact for some $h > 0$ with $h\omega < 1$ if and only if the identity $D(A) \rightarrow E$ is compact. If $D(A) \rightarrow E_0$ is compact, then J_h , as a map from E to E_0 is compact for all $h > 0$ with $h\omega < 1$. It is worth to note that in the situation described in part (1), $\|J_h\| \leq (1 - h\omega)^{-1}$.

Let us recall a version of Lemma 17 from [6]; for the reader’s convenience we give an independent proof.

Lemma 3.2. *For any $\varepsilon > 0$, there exists a locally Lipschitz map $f = f_\varepsilon : K_0 \rightarrow E$ being an ε -graph-approximation ⁽³⁾ of Φ and, for all $u \in K_0$,*

$$(3.3) \quad f(u) \in T_K(j(u)).$$

PROOF: Take $\varepsilon > 0$ and $u \in K_0$. By (A_5) and (2.7) there is $v(u) \in E$ such that

$$v(u) \in B(\Phi(u), \varepsilon/4) \cap S_K(j(u)).$$

Hence, there is $\alpha(u) > 0$ such that

$$j(u) + \alpha(u)v(u) \in K.$$

By the H -upper semicontinuity choose a number $\gamma(u)$, $0 < \gamma(u) < \varepsilon/4$ such that $\Phi(B_0(u, 2\gamma(u)) \cap K_0) \subset B(\Phi(u), \varepsilon/2)$ and a number $0 < \delta(u) < \min\{\gamma(u)/C, \gamma(u)/\alpha(u)\}$, where $C := \|j\|$.

Let $\{\lambda_s\}_{s \in S}$ be a locally finite locally Lipschitzian partition of unity refining the open cover $\{B_0(u, \delta(u)\alpha(u)) \cap K_0\}_{u \in K_0}$. For any $s \in S$, there is $u_s \in K$ such that the support $\text{supp } \lambda_s \subset$

³The graph of f is contained in the ε -neighborhood of the graph of Φ , i.e. $f(x) \in \Phi(B_0(x, \varepsilon) \cap K_0) + \varepsilon B$, where B is the unit open ball in E , for any $x \in K_0$; in particular f is bounded. Here and below we write $B(x, r)$, $x \in E$, to denote a ball in E and $B_0(x, r)$, $x \in E_0$, to denote a ball in E_0 .

$B_0(u_s, \delta_s \alpha_s) \cap K_0$ where we have put $\delta_s := \delta(u_s)$ and $\alpha_s := \alpha(u_s)$. Additionally let us set $v_s := v(u_s)$ and $\gamma_s := \gamma(u_s)$.

For any $s \in S$, we define a map $f_s : K_0 \rightarrow E$ by the formula

$$f_s(u) := \frac{1}{\alpha_s}(j(u_s) - j(u)) + v_s, \quad u \in K_0.$$

Observe, that for $s \in S$, $u \in K_0$, $j(u) + \alpha_s f_s(u) = j(u_s) + \alpha_s v_s \in K$. Hence, for all $u \in K_0$,

$$f_s(u) \in S_K(j(u)) \subset T_K(j(u)).$$

It is clear that f_s , $s \in S$, is Lipschitz continuous.

Now we define $f : K_0 \rightarrow E$ by the formula

$$f(u) := \sum_{s \in S} \lambda_s(u) f_s(u), \quad u \in K_0.$$

Observe that f is locally Lipschitz because so are all functions λ_s , f_s for $s \in S$, and the covering $\{\text{supp } \lambda_s\}_{s \in S}$ is locally finite. Moreover, since, for $u \in K_0$, $f(u)$ is a (finite) convex combination of vectors $f_s(u) \in T_K(j(u))$ and since $T_K(j(u))$ is convex, we see that $f(u) \in T_K(j(u))$ for all $u \in K_0$.

Take $u \in K_0$ and let $S(u) = \{s \in S \mid u \in \text{supp } \lambda_s\}$. It is clear that $S(u)$ is a finite set and

$$f(u) = \sum_{s \in S(u)} \lambda_s(u) f_s(u).$$

For any $s \in S(u)$, we have $u \in \text{supp } \lambda_s \subset B_0(u_s, \delta_s \alpha_s) \cap K_0$, i.e.

$$\|u - u_s\|_0 < \delta_s \alpha_s < \gamma_s \quad \text{and} \quad \|f_s(u) - v_s\| < \delta_s C < \gamma_s.$$

There is $s_0 \in S(u)$ such that $\gamma_{s_0} = \max_{s \in S(u)} \gamma_s$. If $s \in S(u)$, then

$$\|u_s - u_{s_0}\|_0 \leq \|u_s - u\|_0 + \|u_{s_0} - u\|_0 < \gamma_s + \gamma_{s_0} \leq 2\gamma_{s_0}.$$

Therefore, for any $s \in S(u)$,

$$\begin{aligned} f_s(u) &\in B(v_s, \gamma_{s_0}) \subset \Phi(u_s) + (\varepsilon/4 + \gamma_{s_0})B \subset \Phi(B_0(u_{s_0}, 2\gamma_{s_0}) \cap K_0) + (\varepsilon/4 + \gamma_{s_0})B \\ &\subset B(\Phi(u_{s_0}), \varepsilon/4 + \varepsilon/2 + \gamma_{s_0}) \subset \Phi(u_{s_0}) + \varepsilon B. \end{aligned}$$

Hence, by convexity of $\Phi(u_{s_0}) + \varepsilon B$,

$$f(u) \in \Phi(u_{s_0}) + \varepsilon B \subset \Phi(B_0(u, \gamma_{s_0}) \cap K_0) + \varepsilon B \subset \Phi(B_0(u, \varepsilon) \cap K_0) + \varepsilon B. \quad \square$$

Lemma 3.3. *For every $u \in K_0$, we have*

$$(3.4) \quad \lim_{h \rightarrow 0^+, v \rightarrow u, v \in K_0} \frac{1}{h} d_K(j(v) + hf(v)) = 0,$$

where f comes from Lemma 3.2.

PROOF: Choose $u \in K_0$ and $\varepsilon > 0$. Taking into account (3.3), (2.7) and the continuity of f there is $\delta > 0$ such that if $v \in K_0$, $\|v - u\|_0 < \delta$ and $0 < h < \delta$ then

$$d_K(j(v) + hf(u)) < \varepsilon h/2 \quad \text{and} \quad \|f(v) - f(u)\| < \varepsilon/2.$$

Thus, for such v and h we have

$$d_K(j(v) + hf(v)) \leq d_K(j(v) + hf(u)) + h\|f(u) - f(v)\| < \varepsilon h. \quad \square$$

Now we are ready to prove

Theorem 3.4. *In addition to $(A_1) - (A_5)$ above, let us assume that K is bounded and for all sufficiently small $h > 0$:*

- (A_6) $J_h(K) \subset K_0$;
- (A_7) $J_h : E \rightarrow E_0$ is compact.

Then there is $u \in K_0 \cap D(A)$ such that $0 \in Au + \Phi(u)$.

PROOF: Choose $\varepsilon > 0$ and $f = f_\varepsilon$ according to Lemma 3.2. Denote by $r : E \rightarrow K$ an \mathcal{L} -retraction onto K , i.e. $\|r(u) - u\| \leq Ld_K(u)$ for $u \in E$ for some $L \geq 1$. For every sufficiently small $h > 0$ with $h\omega < 1$, the map $J_h \circ r(j + hf) : K_0 \rightarrow K_0$ is well-defined due to (A_6) . Moreover $J_h \circ r(j + hf)$ is continuous and compact, since K and Φ (and so does f) are bounded and J_h is compact. Then, by the Schauder fixed point theorem, for large $n \geq 1$ (precisely for $n > \omega$), there is u_n in K_0 such that

$$u_n = J_{1/n} \circ r \left(j(u_n) + \frac{1}{n}f(u_n) \right),$$

so $u_n \in D(A)$ and

$$j(u_n) - \frac{1}{n}Au_n = r \left(j(u_n) + \frac{1}{n}f(u_n) \right).$$

As a result, we have

$$Au_n + f(u_n) = n \left(j(u_n) + \frac{1}{n}f(u_n) - r \left(j(u_n) + \frac{1}{n}f(u_n) \right) \right).$$

Therefore

$$(3.5) \quad \|Au_n + f(u_n)\| \leq nLd_K \left(j(u_n) + \frac{1}{n}f(u_n) \right) \leq L\|f(u_n)\|.$$

Hence $\{Au_n\}_{n \geq 1}$ is bounded in E . Fix $h > 0$ with $h\omega < 1$ and note that

$$\{u_n\}_{n \geq 1} = J_h \left(\{j(u_n) - hAu_n\}_{n \geq 1} \right).$$

Since $\{j(u_n) - hAu_n\}_{n \geq 1}$ is bounded in E , the above equality yields $\{u_n\}_{n \geq 1}$ is relatively compact in E_0 . Passing to a subsequence if necessary, we can assume that $u_n \rightarrow u_\varepsilon$ in E_0 and $u_\varepsilon \in K_0$. In view of (3.5) and Lemma 3.3 we have

$$(3.6) \quad \|Au_n + f(u_n)\| = nd_K \left(j(u_n) + \frac{1}{n}f(u_n) \right) \rightarrow 0, \quad \text{as } n \rightarrow \infty.$$

Hence $Au_n \rightarrow -f(u_\varepsilon)$ in E . The closedness of A yields $u_\varepsilon \in D(A) \cap K_0$ and $-Au_\varepsilon = f(u_\varepsilon)$.

Arguing as above, we may assume without loss of generality that $u_\varepsilon \rightarrow u_0 \in K_0$ as $\varepsilon \rightarrow 0$. Let $v_\varepsilon := -Au_\varepsilon$. Since

$$v_\varepsilon := f(u_\varepsilon) \in \Phi(B_0(u_\varepsilon, \varepsilon) \cap K_0) + \varepsilon B,$$

there is $u'_\varepsilon \in K_0$ and $v'_\varepsilon \in \Phi(u'_\varepsilon)$ such that $\|u_\varepsilon - u'_\varepsilon\|_0 < \varepsilon$ and $\|v_\varepsilon - v'_\varepsilon\| < \varepsilon$. Clearly $(u'_\varepsilon, v'_\varepsilon) \in \text{Gr}(\Phi)$ and $u'_\varepsilon \rightarrow u_0$; in view of Remark 3.1 (2) we gather that, after passing to a subsequence if necessary, $v'_\varepsilon \rightarrow v_0 \in \Phi(u_0)$. Since $v_\varepsilon \rightarrow v_0$, too, and the graph of A , being closed and convex, is also weakly closed, we see that $u_0 \in D(A)$ and $-Au_0 = v_0 \in \Phi(u_0)$. \square

Now we are going to establish a counterpart of Theorem 3.4 for \mathcal{L} -retracts. In this case the choice of E_0 is immaterial since we shall assume that Φ is defined on K . In addition let us assume that:

- (B_1) E is a Banach space and $K \subset E$ is a bounded \mathcal{L} -retract;

- (B₂) $A : D(A) \rightarrow E$ is a densely defined linear operator such that for some $\omega \in \mathbb{R}$, $(\omega, +\infty) \subset \rho(A)$ and $\|(A - \lambda I)^{-1}\| \leq (\lambda - \omega)^{-1}$ for $\lambda > \max\{0, \omega\}$ ⁽⁴⁾;
- (B₃) $\Phi : K \rightarrow E$ is bounded H -upper semicontinuous with convex weakly compact values;
- (B₄) For $h > 0$ with $h\omega < 1$, $J_h : E \rightarrow E$ is compact and K is resolvent invariant, i.e. $J_h(K) \subset K$ for sufficiently small h .

Remark 3.5. (1) In view (B₂) and (B₄), the Euler characteristic $\chi(K)$ is well-defined. Indeed, by Remark 2.8, it is sufficient to prove that the identity I_K is homotopic to a compact map. To this end fix $h > 0$ with $h\omega < 1$, such that J_h is compact and consider $h : K \times [0, 1] \rightarrow K$ given by the formula

$$h(x, t) = \begin{cases} J_{th}(x) & \text{if } t \in (0, 1], \\ x & \text{if } t = 0, \end{cases} \quad x \in K.$$

Assumption (B₂) implies that $\lim_{t \rightarrow 0^+} J_{th}(x) = x$ for all $x \in E$. Moreover the map $E \times (0, 1] \ni (x, t) \mapsto J_{th}(x) \in E$ is continuous. Thus h is continuous and provides a homotopy joining the identity to the compact map J_h . As a consequence if $\chi(K) \neq 0$, then any compact map $g : K \rightarrow K$ homotopic to the identity has fixed points.

(2) If A is given as in Remark 3.1 (with $E = H$), then assumption (B₂) is satisfied. Moreover, in this case the resolvent invariance holds for a convex K if and only if K is semigroup invariant, i.e. $S(t)K \subset K$ for any $t \geq 0$. Indeed if K is resolvent invariant then by the Post-Widder formula (see [26, Corollary 5.5, 5.6]) for each $x \in K$ and $t > 0$

$$S(t)x = \lim_{n \rightarrow \infty} J_{t/n}^n x \in K.$$

Conversely, if K is semigroup invariant, then by [26, Th. 1.10], for any $h > 0$ with $h\omega < 1$ and $x \in K$,

$$J_h x = \frac{1}{h} \int_0^\infty e^{-t/h} S(t)x dt \in K.$$

Theorem 3.6. *Under assumptions (B₁) – (B₄), the problem (3.1) has a solution if the Euler characteristic $\chi(K) \neq 0$ and Φ satisfies the weak tangency condition in the sense of Clarke cones, i.e.*

$$(3.7) \quad \forall u \in K \quad \Phi(u) \cap C_K(u) \neq \emptyset.$$

First we need a result, similar to Lemma 3.2, which may be of interest on its own.

Lemma 3.7. *Suppose $X \subset E$ is closed and that $\Phi : X \rightarrow E$ is H -upper semicontinuous with convex values. Let a function $\xi : X \times E \rightarrow \mathbb{R}$ be such that for each $z \in E$, $\xi(\cdot, z)$ is upper semicontinuous (as a real function) and for each $x \in X$, $\xi(x, \cdot)$ is convex. If for all $x \in X$, $\inf_{z \in \Phi(x)} \xi(x, z) \leq 0$, then for any $\varepsilon > 0$ there exists a locally Lipschitz ε -graph-approximation $f = f_\varepsilon : X \rightarrow E$ of Φ such $\xi(x, f(x)) < \varepsilon$ for all $x \in X$.*

PROOF: For any $z \in X$ choose $0 < \delta_z < \varepsilon$ such that $\Phi(B(z, \delta_z) \cap X) \subset \Phi(z) + \varepsilon B$ and let an open covering \mathcal{U} of X be a star refinement of the covering $\{B(z, \delta_z) \cap X\}_{z \in X}$ of X .

For each $x \in X$ choose $z_x \in \Phi(x)$ such that $\xi(x, z_x) < \varepsilon$. Given $U \in \mathcal{U}$ and $x \in U$ let

$$V_U(x) := \{y \in U \mid \xi(y, z_x) < \varepsilon\}.$$

⁴I.e. A is maximal ω -dissipative.

Clearly $x \in V_U(x)$. Hence $\mathcal{V} := \{V_U(x)\}_{U \in \mathcal{U}, x \in U}$ is an open cover of X . Let $\{\lambda_s\}_{s \in S}$ be a locally Lipschitz partition of unity subordinated to \mathcal{V} , i.e. for any $s \in S$, there is $U_s \in \mathcal{U}$, $x_s \in U_s$ such that $\text{supp } \lambda_s \subset V_s := V_{U_s}(x_s)$. Let

$$f(x) := \sum_{s \in S} \lambda_s(x) z_s, \quad x \in X,$$

where $z_s := z_{x_s}$. Then f is well-defined and locally Lipschitz.

Take $x \in X$ and let $S(x) := \{s \in S \mid \lambda_s(x) \neq 0\}$. If $s \in S(x)$, then $x \in V_s$, i.e. $\xi(x, z_s) < \varepsilon$. By the convexity of $\xi(x, \cdot)$ we gather that $\xi(x, f(x)) < \varepsilon$. Since $x_s \in U_s$ and \mathcal{U} is a star refinement of $\{B(z, \delta_z) \cap X\}$ we get that for all $s \in S(x)$, x, x_s belong to the star of x with respect to \mathcal{U} :

$$x, x_s \in \bigcup_{\{U \in \mathcal{U} \mid x \in U\}} U \subset B(z, \delta_z) \cap X$$

for some $z \in Z$. Hence $z \in B(x, \varepsilon)$ and for $s \in S(x)$, $z_s \in \Phi(x_s) \subset \Phi(z) + \varepsilon B$. This together with the convexity of $\Phi(z)$ shows that

$$f(x) \in \Phi(z) + \varepsilon B \subset \Phi(B(x, \varepsilon)) + \varepsilon B. \quad \square$$

PROOF of Theorem 3.6: Let

$$\xi(u, v) := \sup_{p \in \partial d_K(u)} \langle p, v \rangle \in \mathbb{R}, \quad u \in K, v \in E,$$

where $\partial d_K(u) \subset E^*$ denotes the generalized Clarke gradient at $u \in K$ of the (locally Lipschitz) function d_K . It is clear that

$$\xi(u, v) = d_K^\circ(u; v) = \limsup_{y \rightarrow u, y \in K, t \rightarrow 0^+} \frac{d_K(y + tv)}{t}$$

is the Clarke directional derivative of d_K at u in the direction v . Then $\xi : K \times E \rightarrow \mathbb{R}$ is upper semicontinuous and, for each $u \in K$, $\xi(u, \cdot)$ is convex.

Observe now that

$$C_K(u) = [\partial d_K(u)]^- := \{v \in E \mid \xi(u, v) \leq 0\}.$$

Condition (3.7) together with with (B_3) show that all assumptions of Lemma 3.7 are satisfied.

Take an \mathcal{L} -retraction $r : B(K, \eta) \rightarrow K$ with constant L . Since Φ is bounded, there are $\lambda_0 > 0$ and $\varepsilon_0 > 0$ such that

$$\forall u \in K \quad u + \lambda f(u) \in B(K, \eta),$$

for any $f : K \rightarrow E$ being an ε -graph-approximation of Φ with $0 < \varepsilon < \varepsilon_0$ and $0 < \lambda \leq \lambda_0$.

Suppose now to the contrary that there are no solutions to (3.1). We claim that there is $0 < \varepsilon < \varepsilon_0$ such that if $u \in K \cap D(A)$ and f is an ε -graph-approximation of Φ , then

$$(3.8) \quad \|Au + f(u)\| \geq (L + 1)\varepsilon.$$

If not then there are sequences $\varepsilon_n > \varepsilon_n \rightarrow 0^+$, $u_n \in K$ and an ε_n -approximation $f_n : K \rightarrow E$ of Φ such that

$$\|Au_n - f_n(u_n)\| < (L + 1)\varepsilon_n, \quad n \in \mathbb{N}.$$

This implies that the sequence (Au_n) is bounded; hence by the same argument as in the proof of Theorem 3.4, we gather that, passing to a subsequence if necessary, $u_n \rightarrow u_0 \in K$.

Since $f_n(u_n) \in \Phi(B(u_n, \varepsilon_n)) + \varepsilon_n B$, we find $u'_n \in B(u_n, \varepsilon_n)$ and $v'_n \in \Phi(u'_n)$ such that $\|f_n(u_n) - v'_n\| < \varepsilon_n$. By Remark 3.1 (2) we may assume that $v'_n \rightarrow v_0 \in \Phi(u_0)$. This implies that $f_n(u_n) \rightarrow v_0$

and, thus $-Au_n \rightarrow v_0$, too. Hence $v_0 = -Au_0$, i.e. $0 \in Au_0 + \Phi(u_0)$: a contradiction.

Now take $\varepsilon > 0$ provided above and, using Lemma 3.7, let $f : K \rightarrow \mathbb{E}$ be an ε -graph-approximation of Φ such that $\xi(u, f(u)) < \varepsilon$ for all $u \in K$. Take a decreasing sequence $h_n \rightarrow 0^+$ with $h_1 < \lambda$. Since f is bounded the map

$$K \ni u \mapsto g_n(u) := J_{h_n} \circ r(u + h_n f(u)) \in K$$

is well-defined a compact. Moreover $h : K \times [0, 1] \rightarrow K$ given by

$$h(u, t) := \begin{cases} J_{th_n} \circ r(u + th_n f(u)) & \text{if } t \in (0, 1]; \\ u & \text{if } t = 0, \end{cases} \quad u \in K$$

provides a (continuous) homotopy joining the identity on K with g_n . In view of Remark 3.5, $g_n(u_n) = u_n$ for some $u_n \in K \cap D(A)$. This means that

$$(3.9) \quad Au_n + f(u_n) = h_n^{-1}(u_n + h_n f(u_n) - r(u_n + h_n f(u_n))).$$

Similarly as before we may suppose that $u_n \rightarrow u_0 \in K$; therefore $f(u_n) \rightarrow f(u_0)$. By (3.8) and (3.9), for all $n \in \mathbb{N}$,

$$(L + 1)\varepsilon \leq \|Au_n + f(u_n)\| \leq Lh_n^{-1}d_K(u_n + h_n f(u_n)) \leq Lh_n^{-1}d_K(u_n + h_n f(u_0)) + L\|f(u_n) - f(u_0)\|.$$

Passing to lim sup and remembering that $\xi(u_0, f(u_0)) < \varepsilon$ we get

$$(L + 1)\varepsilon \leq \limsup_{n \rightarrow \infty} \|Au_n + f(u_n)\| \leq L\varepsilon,$$

a contradiction. This completes the proof. \square

3.2. Invariance and viability. A central role among assumptions of Theorems 3.4 and 3.6 is played by the resolvent invariance of the set K and the tangency condition. Let us consider conditions (B_1) – (B_4) and (3.7) and let K be an arbitrary closed subset of E . The Hille-Yosida Theorem implies that in this case A is the generator of a C_0 semigroup $\{S(t)\}_{t \geq 0}$. It is not difficult to show that (B_4) and (3.7) imply that

$$\forall u \in K \quad \liminf_{v \rightarrow u, v \in K, h \rightarrow 0^+} \frac{\text{dist}(J_h(v + h\Phi(v)), K)}{h} = 0 \quad (5).$$

This condition implies that

$$(3.10) \quad \forall u \in K \quad \Phi(u) \cap T_K^A(u) \neq \emptyset,$$

where

$$T_K^A(u) = \left\{ v \in E \mid \liminf_{h \rightarrow 0^+} \frac{d_K(S_v(h)u)}{h} = 0 \right\}$$

and $[0, +\infty) \ni t \mapsto S_v(t)u$ is the mild solution to the problem $\dot{x} = Ax + v$, $x(0) = u$. Finally (3.10) is equivalent to the following:

$$(3.11) \quad \begin{aligned} & \text{the problem } \dot{u} \in Au + \Phi(u), \quad u(0) = x \in M \text{ has} \\ & \text{a mild solution } u : [0, +\infty) \rightarrow E \text{ such that } u(t) \in M \text{ for } t \geq 0, \end{aligned}$$

because the semigroup $\{S(t)\}_{t \geq 0}$ is immediately compact. This and related results are thoroughly discussed in [56] and [10]. We thus see that our conditions imply the **invariance** of K (sometimes

⁵Here dist stands for the distance between sets, i.e. $\text{dist}(X, Y) := \inf\{\|x - y\| \mid x \in X, y \in Y\}$.

called viability) with respect to the ‘heat flow’ generated by A , i.e. condition (3.11). Conversely condition

$$(3.12) \quad \forall u \in K \quad 0 \in T_K^A(u)$$

implies the semigroup invariance and, in case of a convex K , resolvent invariance (B_4). The point is that, in concrete situations of differential problems, condition (3.12) needs to be verified. In most cases this can be done via an appropriate use of the maximum principles. In the next section we shall encounter examples of such arguments.

The problem of invariance of systems of parabolic PDE was studied in numerous papers [1], [38], [50], [6], [10], [57] (and references therein), [58] (the so-called Müller conditions important in various applications). The most general, often necessary and sufficient, abstract results are presented in [56]. The invariance problem of parabolic problem from (3.11) will be studied in the forthcoming paper [36]. In particular we shall study the topological structure of the set of all viable (i.e. ‘surviving’ in K) solutions and show its relation with the existence of steady states, i.e. solutions to (3.1).

4. APPLICATIONS

4.1. The Neumann problem I. We now study the existence of steady state solutions to (2.3) and consider the problem

$$(4.1) \quad \begin{cases} -Lu \in F(x, u, \nabla u), & u(x) \in C \text{ a.e. on } \Omega, \\ \frac{\partial}{\partial \mathbf{n}} u = 0 & \text{on } \partial\Omega, \end{cases}$$

where $C \subset \mathbb{R}^N$ is a compact and convex set, $\frac{\partial}{\partial \mathbf{n}} u = \left(\frac{\partial u_1}{\partial \mathbf{n}}, \dots, \frac{\partial u_N}{\partial \mathbf{n}} \right)$ denotes the outward normal derivative of u . We are going to find a *strong solutions*: a function $u \in H^2(\Omega, \mathbb{R}^N)$ such that $-Lu(x) \in F(x, u(x), \nabla u(x))$ for a.a. $x \in \Omega$ and $\frac{\partial u_i}{\partial \mathbf{n}} \Big|_{\partial\Omega} = 0$, $i = 1, \dots, N$, in the sense of trace.

Let us make the following assumptions:

- (D) For all $i = 1, \dots, N$, $d_i = d \in C^1(\bar{\Omega})$ and $d > 0$;
- (F₁) $F : \Omega \times C \times \mathbb{R}^{MN} \rightarrow \mathbb{R}^N$ is upper semicontinuous with compact convex values;
- (F₂) there is a nonnegative $b \in L^2(\Omega)$ such that $\sup_{y \in F(x, u, v)} |y| \leq b(x)$ for a.a. $x \in \Omega$, $u \in C$ and $v \in \mathbb{R}^{MN}$;
- (F₃) F is weakly tangent to C with respect to the second variable, i.e. $F(x, u, v) \cap T_C(u) \neq \emptyset$ for all $x \in \Omega$, $u \in C$ and all $v \in \mathbb{R}^{MN}$.

We now put

$$\begin{aligned} E &:= L^2(\Omega, \mathbb{R}^N), \quad K := \{u \in E \mid u(x) \in C \text{ a.e. on } \Omega\}; \\ E_0 &:= H^1(\Omega, \mathbb{R}^N), \quad K_0 := \{u \in E_0 \mid u(x) \in C \text{ a.e. on } \Omega\}; \\ D(A) &:= \left\{ u \in H^2(\Omega, \mathbb{R}^N) \mid \frac{\partial u}{\partial \mathbf{n}} \Big|_{\partial\Omega} = 0 \right\}, \quad Au := Lu, \quad u \in D(A). \end{aligned}$$

Clearly assumptions (A₁), (A₂) (with $\omega = 0$) and (A₃) are satisfied, K is closed convex and bounded; condition (A₇) holds true since the embedding $D(A) \rightarrow E_0$ is compact.

For any $u \in K_0$, let

$$(4.2) \quad \Phi(u) := \{v \in E \mid v(x) \in F(x, u(x), \nabla u(x)) \text{ for a.a. } x \in \Omega\}.$$

Evidently values $\Phi : K_0 \rightarrow E$ are nonempty and convex.

Proposition 4.1. *The map Φ satisfies conditions (A_4) and (A_5) .*

PROOF: It is straightforward to show that $\Phi(u)$ is weakly compact (we work in a Hilbert space, thus closed convex and bounded sets are convex weakly compact). Below we shall prove a slight generalization of Proposition 6.2 from [10]. It implies immediately that Φ is H -upper semicontinuous.

Lemma 4.2. *If $\psi : \Omega \times \mathbb{R}^d \rightarrow \mathbb{R}^N$ is upper semicontinuous with convex compact values and $\sup_{y \in \psi(x,u)} |y| \leq b(x) + a|u|$ for all $u \in \mathbb{R}^d$ and a.a. $x \in \Omega$, where $b \in L^2(\Omega)$ and $a > 0$, then the Nemytskii operator*

$$\Psi : \mathcal{E} := L^2(\Omega, \mathbb{R}^d) \rightarrow E, \quad \Psi(u) := \{v \in E \mid v(x) \in \psi(x, u(x)) \text{ for a.a. } x \in \Omega\}, \quad u \in \mathcal{E},$$

is H -upper semicontinuous.

PROOF: Suppose it is not the case: there are $\varepsilon_0 > 0$, a sequences $u_n \rightarrow u_0$ in \mathcal{E} and $v_n \in \Psi(u_n)$ such that

$$(4.3) \quad v_n \notin \Psi(u_0) + B_E(0, \varepsilon_0), \quad n \geq 1.$$

Up to a subsequence $(u_n)_{n \geq 1}$ converges a.e. on Ω to u_0 and there is $h \in L^2(\Omega, \mathbb{R})$ such that $|u_n(x)| \leq h(x)$ for a.e. $x \in \Omega$ and every $n \geq 0$. By assumption

$$|v_n(x)| \leq b(x) + a|u_n(x)| \leq b(x) + ah(x) \text{ for } n \geq 0 \text{ and a.e. } x \in \Omega.$$

There is $\eta > 0$ such that for $A \subset \Omega$ with Lebesgue measure $\mu(A) < \eta$

$$(4.4) \quad \int_A 4(b(x) + ah(x))^2 dx < \varepsilon_0^2/2.$$

For each $n \geq 0$, the set-valued map $H_n := \psi(\cdot, u_n(\cdot)) : \Omega \rightarrow \mathbb{R}^N$ is measurable and if $w : \Omega \rightarrow \mathbb{R}^N$ is a measurable selection of H_n , then $w \in E$ since

$$(4.5) \quad |w(x)| \leq b(x) + ah(x) \quad \text{for a.e. } x \in \Omega.$$

By the Egorov and Lusin theorems (see [5, Th. 1] for a multivalued version of the Lusin theorem) there is a compact $\Omega_\eta \subset \Omega$ such that $\mu(\Omega \setminus \Omega_\eta) < \eta$, $u_n \rightarrow u_0$ uniformly on Ω_η , the restriction $u_0|_{\Omega_\eta} : \Omega_\eta \rightarrow \mathbb{R}^d$ is continuous and $H_0|_{\Omega_\eta} : \Omega_\eta \rightarrow \mathbb{R}^N$ is H -lower semicontinuous.

Let $\delta := \varepsilon_0/\sqrt{2\mu(\Omega)}$. We will show that there is n_0 such that if $n \geq n_0$ and $x \in \Omega_\eta$, then

$$H_n(x) \subset H_0(x) + B_{\mathbb{R}^N}(0, \delta).$$

Suppose to the contrary that there is a subsequence $(n_j)_{j \geq 1}$ and a sequence $(x_j)_{j \geq 1}$ in Ω_η such that

$$(4.6) \quad H_{n_j}(x_j) \not\subset H_0(x_j) + B_{\mathbb{R}^N}(0, \delta).$$

We can assume that $x_j \rightarrow x_0 \in \Omega_\eta$, since Ω_η is compact. The continuity of $u_0|_{\Omega_\eta}$ and the uniform convergence $u_n \rightarrow u_0$ on Ω_η imply that $u_{n_j}(x_j) \rightarrow u_0(x_0)$ and thus $(x_j, u_{n_j}(x_j)) \rightarrow (x_0, u_0(x_0))$ as $j \rightarrow \infty$. The upper semicontinuity of ψ together with the H -lower semicontinuity of H_0 on Ω_η show that $H_{n_j}(x_j) \subset H_0(x_j) + B_{\mathbb{R}^N}(0, \delta)$ for sufficiently large j , which contradicts (4.6).

Let us fix $n \geq n_0$. For a.e. $x \in \Omega_\eta$ we have

$$(4.7) \quad v_n(x) \in H_n(x) \subset H_0(x) + B_{\mathbb{R}^N}(0, \delta).$$

Observe that the map $\Omega_\eta \ni x \mapsto B_{\mathbb{R}^N}(v_n(x), \delta) \cap H_0(x)$ is measurable and has nonempty values for a.e. $x \in \Omega_\eta$. By the Kuratowski-Ryll-Nardzewski theorem, there is a measurable selection

$v : \Omega_\eta \rightarrow \mathbb{R}^N$, i.e. $v(x) \in B_{\mathbb{R}^N}(v_n(x), \delta) \cap H_0(x)$ for a.e. $x \in \Omega_\eta$. Clearly $v \in L^2(\Omega_\eta, \mathbb{R}^N)$ and for a.e. $x \in \Omega_\eta$, $|v_n(x) - v(x)| < \delta$. Thus

$$\int_{\Omega_\eta} |v_n(x) - v(x)|^2 dx < \delta^2 \mu(\Omega_\eta) < \varepsilon_0^2/2.$$

Take an arbitrary selection $w : \Omega \rightarrow \mathbb{R}^N$ of H_0 , i.e. $w(x) \in H_0(x)$ for a.e. $x \in \Omega$. Let $\chi = \chi_{\Omega_\eta}$ be the indicator of Ω_η . Notice that $\chi v + (1 - \chi)w : \Omega \rightarrow \mathbb{R}^N$ is a square-integrable selection of H_0 (we identify $v : \Omega_\eta \rightarrow \mathbb{R}^N$ with the function $v : \Omega \rightarrow \mathbb{R}^N$ putting $v \equiv 0$ on $\Omega \setminus \Omega_\eta$). By (4.5)

$$|v_n(x) - w(x)| \leq |v_n(x)| + |w(x)| \leq 2(b(x) + ah(x)) \quad \text{for a.e. } x \in \Omega \setminus \Omega_\eta.$$

Recall that $\mu(\Omega \setminus \Omega_\eta) < \eta$, hence and by (4.4)

$$\begin{aligned} \|v_n - \chi v + (1 - \chi)w\|^2 &= \int_{\Omega_\eta} |v_n(x) - v(x)|^2 dx + \int_{\Omega \setminus \Omega_\eta} |v_n(x) - w(x)|^2 dx \\ &< \varepsilon_0^2/2 + \int_{\Omega \setminus \Omega_\eta} 4(\alpha(x) + h(x))^2 dx < \varepsilon_0^2. \end{aligned}$$

Thus, contrary to (4.3), $v_n \in \Psi(t_0, u_0) + B_{L^2(\Omega, \mathbb{R}^N)}(0, \varepsilon_0)$ for infinitely many $n \geq 1$. \square

In order to get the weak tangency (A₅) fix $u \in K_0$ and define $G, H : \Omega \rightarrow \mathbb{R}^N$, by

$$G(x) := F(x, u(x), \nabla u(x)), \quad H(x) := T_C(u(x)) \quad \text{for } x \in \Omega.$$

The map $T_C(\cdot) : C \rightarrow \mathbb{R}^N$ is lower semicontinuous (see [4, Th. 4.2.2]), G is measurable; hence $\Omega \ni x \rightarrow G(x) \cap H(x) \subset \mathbb{R}^N$ is measurable with nonempty values. By the Kuratowski–Ryll–Nardzewski theorem, there is a measurable $v : \Omega \rightarrow \mathbb{R}^N$ such that $v(x) \in G(x) \cap H(x)$ for a.e. $x \in \Omega$. Clearly $v \in E$ and $v \in T_K(u) \cap F(t, u)$ since in view of [4, Cor. 8.5.2] $T_K(u) = \{v \in E \mid v(x) \in T_C(u(x)) \text{ a.e. in } \Omega\}$. \square

Proposition 4.3. *For any $h > 0$ the resolvent J_h maps K into K_0 .*

PROOF: In view of [9, Cor. 7.49], C is an intersection of countably many closed half-spaces containing it, i.e. $C = \bigcap_{n \geq 1} C_n$, where $C_n := \{x \in \mathbb{R}^N \mid p_n \cdot x \leq a_n\}$ for some $p_n \in \mathbb{R}^N$ and $a_n \in \mathbb{R}$. Thus, it is enough to show that $J_h(K_n) \subset K_n$, where

$$K_n = \{u \in E \mid u(x) \in C_n \text{ for a.a. } x \in \Omega\}$$

for every $n \geq 1$, since then

$$J_h(K) = J_h\left(\bigcap_{n \geq 1} K_n\right) \subset \bigcap_{n \geq 1} J_h(K_n) \subset \bigcap_{n \geq 1} K_n = K \text{ } ^{(6)}.$$

But $J_h(K) \subset E_0$ so, eventually, $J_h(K) \subset K_0$.

Without loss of generality, we assume that $C = \{x \in \mathbb{R}^N \mid p \cdot x \leq a\}$ for some $p \in \mathbb{R}^N$ and $a \in \mathbb{R}$ and $K = \{u \in E \mid p \cdot u(x) \leq a \text{ for a.a. } x \in \Omega\}$. Take $f \in K$ and put $u = J_h(f)$. By definition $u \in D(A)$ and

$$u - hAu = f.$$

⁶Observe that in order to have that $\bigcap K_n \subset K$ one needs the *countable* collection of supporting functionals.

Define $\bar{f}(x) := p \cdot f(x)$, $\bar{u}(x) := p \cdot u(x)$ for $x \in \Omega$. Observe that $\bar{f} \leq a$ a.e., $\bar{u} \in H^2(\Omega)$ and for every $\xi \in H^1(\Omega)$

$$\int_{\Omega} \bar{u}(x)\xi(x) dx + h \int_{\Omega} d(x)\nabla\bar{u}(x)\nabla\xi(x) dx = \int_{\Omega} \bar{f}(x)\xi(x) dx$$

what yields

$$\int_{\Omega} (\bar{u}(x) - a)\xi(x) dx = \int_{\Omega} (\bar{f}(x) - a)\xi(x) dx - h \int_{\Omega} d(x)\nabla(\bar{u} - a)(x)\nabla\xi(x) dx.$$

Taking $\xi = (\bar{u} - a)_+ := \max\{0, \bar{u} - a\}$, we have $\xi \in H^1(\Omega)$ and $\nabla\xi = \chi\nabla(\bar{u} - a)$ by [17, Cor. 1.3.6], where $\chi = \chi_{\{\bar{u} > a\}}$. Therefore, for such ξ :

$$0 \leq \int_{\Omega} (\bar{u} - a)_+^2(x) dx = \int_{\Omega} (\bar{f}(x) - a)(\bar{u} - a)_+(x) dx - h \int_{\{\bar{u} > a\}} d(x)|\nabla(\bar{u} - a)_+(x)|^2 dx \leq 0.$$

As a result $\bar{u} \leq a$ a.e., that is $J_h(f) \in K$. \square

In view of Propositions 4.1 and 4.3 we get

Theorem 4.4. *If assumption (D), (F₁) – (F₃) are satisfied, the problem (4.1) has a solution.* \square

4.2. The Neumann problem II. We will establish the existence of steady state solutions to problem (2.4), i.e.

$$(4.8) \quad \begin{cases} -Lu + \Gamma u \in G(x, u) & u(x) \in C, \text{ a.e. on } \Omega, \\ \frac{\partial}{\partial \mathbf{n}} u = 0 & \text{on } \partial\Omega, \end{cases}$$

where $C \subset \mathbb{R}^N$ is compact and convex. We assume (D) and

- (Γ) for all $i = 1, \dots, N$, $\gamma_i = \gamma = (\gamma_1, \dots, \gamma_M) \in L^\infty(\Omega, \mathbb{R}^M)$; for any $i = 1, \dots, N$;
- (G_1) $G : \bar{\Omega} \times C \rightarrow \mathbb{R}^N$ is upper semicontinuous with compact convex values ⁽⁷⁾;
- (G_2) G is weakly tangent to C , i.e. $G(x, u) \cap T_C(u) \neq \emptyset$ for all $x \in \Omega$ and all $u \in C$.

Similarly as before we put

$$E = L^2(\Omega, \mathbb{R}^N), \quad K := \{u \in E \mid u(x) \in C \text{ for a.a. } x \in \Omega\}.$$

Thus (B_1) is satisfied. Let us define a continuous bilinear form

$$a(u, v) = \int_{\Omega} (d\nabla u \cdot \nabla v + \Gamma u \cdot v) dx, \quad u, v \in H^1(\Omega, \mathbb{R}^N),$$

where $\nabla u \cdot \nabla v$ is the Frobenius product of derivatives ⁽⁸⁾ and $\Gamma u \cdot v = \sum_{i=1}^N (\gamma \cdot \nabla u_i) v_i$. Observe that for any $v \in H^1(\Omega, \mathbb{R}^N)$ and $\varepsilon > 0$,

$$d_0 \|\nabla v\|_{L^2}^2 \leq a(v, v) - \int_{\Omega} \Gamma v \cdot v dx \leq a(v, v) + \|\gamma\|_{L^\infty} \left(\varepsilon \|\nabla v\|_{L^2}^2 + \frac{1}{4\varepsilon} \|v\|_{L^2}^2 \right),$$

where $d_0 = \inf_{x \in \bar{\Omega}} |d(x)|$, in view of the so-called ε -Cauchy inequality. Taking $0 < \varepsilon < d_0/(2\|\gamma\|_{L^\infty})$ we get

$$(4.9) \quad c \|\nabla v\|_{L^2}^2 \leq a(v, v) + C \|v\|_{L^2}^2$$

⁷Note that G is bounded, i.e. $\sup\{|y| \mid y \in G(x, u), x \in \bar{\Omega}, u \in C\} < \infty$.

⁸Recall that if $A = [a_{ij}]$, $B = [b_{ij}]$ are $N \times M$ matrices, then $A \cdot B := \sum_{i=1}^N \sum_{j=1}^M a_{ij} b_{ij}$.

for some positive constants c, C . Therefore for all $v \in H^1(\Omega, \mathbb{R}^N)$

$$c\|v\|_{H^1}^2 \leq a(v, v) + (c + C)\|v\|_{L^2}^2.$$

This implies that we are back in the situation of Remark 3.1 (1), (3) and putting

$$D(A) = \{u \in H^1(\Omega, \mathbb{R}^N) \mid \forall v \in H^1(\Omega, \mathbb{R}^N) \quad a(u, v) = \langle f, v \rangle_{L^2} \text{ for some } f \in L^2(\Omega, \mathbb{R}^N)\}$$

the formula

$$Au = -f, \quad u \in D(A),$$

where f corresponds to u as in the definition of $D(A)$, well defines a closed densely defined linear operator satisfying assumption (B_2) . Moreover A is the generator of a holomorphic C_0 semigroup $\{S(t)\}_{t \geq 0}$. The smoothness of the boundary $\partial\Omega$ and the standard regularity arguments imply that

$$D(A) = \left\{ u \in H^2(\Omega, \mathbb{R}^N) \mid \frac{\partial u}{\partial \mathbf{n}} \Big|_{\partial\Omega} = 0 \right\} \quad \text{and} \quad Au := Lu - \Gamma u, \quad u \in D(A).$$

Now, for any $u \in K$, we put

$$\Phi(u) := \{v \in E \mid v(x) \in G(x, u(x)) \text{ for a.a. } x \in \Omega\}.$$

Arguing as in the proof of Proposition 4.1 we easily get that Φ has properties (B_3) and (3.7).

In order to apply Theorem 3.6 we need

Proposition 4.5. *Condition (B_4) is satisfied.*

PROOF: In view of Remark 3.1 (3) and since $D(A)$ embeds compactly in E we need to check the resolvent invariance of K only. To this end we will use the C_0 semigroup structure. In view of Remark 3.5 (2) we need to show that K is semigroup invariant i.e. $S(t)u_0 \in K$ for all $t \geq 0$ and $u_0 \in K$. It is well known that

$$[0, +\infty) \ni t \mapsto u(t) := S(t)u_0$$

is the unique *mild* solution to the Cauchy initial value problem

$$(4.10) \quad \begin{cases} u' = Au, & u \in E, \quad t > 0 \\ u(0) = u_0. \end{cases}$$

Observe that u is a strong solution, too ⁽⁹⁾. This holds true since the semigroup $\{S(t)\}$ is holomorphic and, thus, immediately differentiable – see [26, Def. II.4.13]; that $u \in W_{loc}^{1,1}((0, +\infty), E)$ follows from Theorem II.4.6 in [26].

Exactly as in the proof of Proposition 4.3, using supporting functionals, we may assume that $C = \{x \in \mathbb{R}^N \mid p \cdot x \leq a\}$ for some $p \in \mathbb{R}^n$ and $a \in \mathbb{R}$. Put $\bar{u} := p \cdot u$, i.e., $\bar{u}(t)(x) = p \cdot u(t)(x)$ and let $v := \bar{u} - a$. Then $v(0) = p \cdot u_0 - a$ and $v(t) \in H^1(\Omega)$ for $t > 0$. It is easy to see that $v : [0, +\infty) \rightarrow L^2(\Omega)$ is a strong solution to the following problem

$$v'(t) = \bar{A}v(t),$$

where

$$\langle \bar{A}\xi, \zeta \rangle_{L^2(\Omega)} = -\bar{a}(\xi, \zeta) := - \int_{\Omega} d\nabla\xi \cdot \nabla\zeta \, dx - \int_{\Omega} (\gamma \cdot \nabla\xi)\zeta \, dx, \quad \xi, \zeta \in H^1(\Omega).$$

⁹Recall that a function $v : [0, +\infty) \rightarrow E$ is a *strong* solution to (4.10) if $v(0) = u_0$, $v(t) \in D(A)$ for $t \in (0, +\infty)$, $v \in W_{loc}^{1,1}((0, +\infty), E)$ and $v'(t) = Av(t)$ for a.a. $t \in (0, +\infty)$ ($v'(t)$ denotes the strong derivative of v which exists a.a. since $v \in W_{loc}^{1,1}$). It is clear that each strong solution is a mild solution.

Therefore $v(t) \in D(\bar{A})$ for $t > 0$ (as in the case of A , the domain of \bar{A} consists of functions in $H^2(\Omega)$ whose normal outward derivative vanishes on $\partial\Omega$ in the sense of trace) and $v \in W_{loc}^{1,1}((0, +\infty), L^2(\Omega))$. For any $\xi \in H^1(\Omega)$ and $t > 0$

$$\langle v'(t), \xi \rangle_{L^2} = -\bar{a}(v(t), \xi).$$

Now let us take $w(t) = v(t)_+$, i.e. $w(t) := \max\{\bar{u}(t) - a, 0\}$, for $t \geq 0$. Then $w(t) \in H^1(\Omega)$; $t > 0$, by [49, Prop. III.1.2] for $t > 0$,

$$\frac{1}{2} \frac{d}{dt} \|w(t)\|_{L^2}^2 = \langle w'(t), w(t) \rangle_{L^2} = \langle v'(t), w(t) \rangle_{L^2} = -\bar{a}(v(t), w(t)) = -\bar{a}(w(t), w(t)).$$

Using a counterpart of estimate (4.9) valid for \bar{a} we see that for

$$\frac{1}{2} \frac{d}{dt} \|w(t)\|_{L^2}^2 \leq -c \|\nabla w(t)\|_{L^2}^2 + C \|w(t)\|_{L^2}^2 \leq C \|w(t)\|_{L^2}^2.$$

By the Gronwall inequality we infer that $\|w(t)\|_{L^2} = 0$ for all $t > 0$, since $w(0) = 0$ in $L^2(\Omega)$ (note that $\bar{u}(0) = p \cdot u_0 \leq a$ since $u_0(x) \in C$ for a.a. $x \in \Omega$). In other words $u(t) \in K$ for all $t > 0$. \square

Theorem 4.6. *If assumptions (D), (Γ) , (G_1) and (G_2) are satisfied, then problem (4.8) has a solution.* \square

Remark 4.7. (1) Observe that theorems 4.4 and 4.6 are true if 0 belongs to the constraint set C and problems (4.1), (4.8) are subject to the Dirichlet condition. The only difference in proof is to see that if $0 \in C$ and a functional p supports C , i.e. $C \subset \{x \in \mathbb{R}^N \mid p \cdot x \leq a\}$, then $a \geq 0$. Hence $(p \cdot u - a)_+ \in H_0^1(\Omega)$.

(2) Let C be convex compact subset of \mathbb{R}^N and $Q = (0, l)^M \subset \mathbb{R}^M$ be an open cube. We will look for solutions to the problem (4.8) with the periodic boundary conditions, i.e.

$$(4.11) \quad \begin{cases} -Lu + \Gamma u \in G(x, u) & u(x) \in C \text{ a.e. on } Q, \\ u|_{\{x_i=0\}} = u|_{\{x_i=l\}}, & \frac{\partial}{\partial x_i} u|_{\{x_i=0\}} = \frac{\partial}{\partial x_i} u|_{\{x_i=l\}}. \end{cases}$$

Assume (D), (Γ) (with $\Omega = Q$) and $d \in C_p^1(Q)$ ⁽¹⁰⁾ and

- (P₁) $G : Q \times C \rightarrow \mathbb{R}^n$ is upper semicontinuous with compact convex values;
- (P₂) G is weakly tangent to C , i.e. $G(x, u) \cap T_C(x) \neq \emptyset$ for all $x \in Q$ and $u \in C$.

Let us put ⁽¹¹⁾:

$$\begin{aligned} E &:= L^2(Q, \mathbb{R}^N), & K &:= \{u \in E \mid u(x) \in C \text{ a.e. on } Q\}, \\ D(A) &:= H_p^2(Q)^N & \text{and } Au &:= Lu - \Gamma u, \quad u \in D(A). \end{aligned}$$

We show, exactly as in Section 4.2, that conditions $(B_1) - (B_4)$ are satisfied. Thus, Theorem 3.6 yields the existence of solutions to (4.7).

¹⁰This symbol stands for the restrictions to Q of functions from $C^1(\mathbb{R}^M)$ which are l -periodic in each direction.

¹¹By $H_p^k(Q)$ we denote the Sobolev space of l -periodic functions on the M -dimensional domain Q (k positive integer); see [45, Chapter 5.10] for definitions and properties of $H_p^k(Q)$.

4.3. Some remarks to the Bernstein theory. In a series of results in [30] authors presented a modern approach to the so-called Bernstein theory for boundary value problems for second order ordinary differential equations (see also [28], [25, II.7.4]; for a numerous research afterwards see e.g. [53] and bibliography therein). For example we formulate a model result [30, Theorem 1.7].

Theorem 4.8. *Suppose that $f : [0, T] \times \mathbb{R} \times \mathbb{R} \rightarrow \mathbb{R}$ is continuous such that*

- (i) *(Sign condition) there is $R \geq 0$ such that $f(t, u, 0)u \leq 0$ for $|u| > R$ and $t \in [0, T]$;*
- (ii) *there are $a, b \geq 0$ so that $|f(t, u, v)| \leq av^2 + b$ for $|u| \leq R$, $t \in [0, T]$ and $v \in \mathbb{R}$.*

Then the problem

$$(4.12) \quad -u'' = f(t, u, u')$$

subject to Dirichlet ($u(0) = u(T) = 0$), Neumann ($u'(0) = u'(T) = 0$) or periodic ($u(0) = u(T)$, $u'(0) = u'(T)$) boundary conditions has a solution in $C^2([0, T], \mathbb{R})$ such that $|u(t)| \leq R$ for $0 \leq t \leq T$.

To illustrate our approach we will stay on the level of an ordinary differential inclusion and study the Dirichlet problem (the Neumann and periodic problems may be studied analogously) for a Sturm-Liouville operator of the form

$$(4.13) \quad -(pu')' + qu \in \varphi(t, u, u'),$$

where

- (φ_0) $p, q : [0, T] \rightarrow \mathbb{R}$ are continuous, $p > 0$ on $[0, T]$ and p is C^1 -smooth;
- (φ_1) $\varphi : [0, T] \times \mathbb{R}^N \times \mathbb{R}^n \rightarrow \mathbb{R}^N$ is upper continuous with compact convex values;
- (φ_2) there are $R \geq 0$ and $c \in \mathbb{R}$ such that $c \leq q(t)$ for $t \in [0, T]$ and $\min_{y \in \varphi(t, u, v)} y \cdot u \leq cR^2$ whenever $t \in [0, T]$, $\|u\| = R$, $v \in \mathbb{R}^N$ with $u \cdot v = 0$;
- (φ_3) φ is bounded on the strip $[0, T] \times C \times \mathbb{R}^N$, where $C := \{y \in \mathbb{R}^N \mid |y| \leq R\}$.

In order to apply Theorem 3.4 let us put

$$\begin{aligned} E &:= L^2((0, T), \mathbb{R}^N), \quad K := \{u \in E \mid |u(x)| \leq R \text{ a.e. on } (0, T)\}; \\ E_0 &:= C_0^1([0, T], \mathbb{R}^N) = \{u \in C^1([0, T], \mathbb{R}^N) \mid u(0) = u(T) = 0\}, \quad K_0 = K \cap E_0; \\ Au &:= u'' + (c - q)u \quad \text{for } u \in D(A) := H^2((0, T), \mathbb{R}^N) \cap C_0([0, T], \mathbb{R}^N). \end{aligned}$$

Moreover define $\Phi : K_0 \rightarrow E$ by

$$\Phi(u) = \{v \in E \mid v(t) \in \varphi(t, u(t), u'(t)) - cu(t) \text{ a.e. on } [0, T]\}, \quad u \in K_0.$$

Within this setting we see that conditions (A_1) , (A_3) and (A_4) are satisfied. As concerns (A_2) note that for $u \in D(A)$ and $v \in V := H_0^1((0, T), \mathbb{R}^N)$

$$-\langle Au, v \rangle_{L^2} = a(u, v)$$

where the bilinear form $a : V \times V \rightarrow \mathbb{R}$ is given by

$$a(u, v) = \int_0^T u' \cdot v' dt + \int_0^T (q(t) - c)u \cdot v dt, \quad u, v \in V.$$

Thus $a(v, v) \geq \|v\|_{H_0^1}^2$ for any $v \in V$, where $\|\cdot\|_{H_0^1}$ is the ‘short’ norm in V . Hence we see that A is the generator of the C_0 semigroup of linear operators on E and conditions (A_2) , (A_7) hold true, since

the inclusion $D(A) \rightarrow E_0$ is compact. Condition (A_6) may be shown as in Proposition 4.3. Therefore we only need

Proposition 4.9. *Condition (A_5) is satisfied.*

PROOF: Let $u \in K_0$, i.e. $u \in C^1([0, T], \mathbb{R}^N)$, $u(0) = u(T) = 0$, $|u(t)| \leq R$ for all $t \in [0, T]$. Let

$$X := \{t \in [0, T] \mid |u(t)| = R\}.$$

It is clear that $X \subset (0, T)$. If $t \in X$, then $u(t) \cdot u'(t) = 0$ since the function $(0, T) \ni s \mapsto |x(s)|^2$ takes maximum at t . Therefore there is $z \in \varphi(t, u(t), u'(t))$ such that $z \cdot u(t) \leq cR^2 = c|u(t)|^2$, i.e. $(z - cu(t)) \cdot u(t) \leq 0$. Hence and by Remark 2.4 (3)

$$[\varphi(t, u(t), u'(t)) - cu(t)] \cap T_C(u(t)) \neq \emptyset.$$

If $t \notin X$, then $T_C(u(t)) = \mathbb{R}^N$ and so

$$\forall t \in [0, T] \quad [\varphi(t, u(t), u'(t)) - cu(t)] \cap T_C(u(t)) \neq \emptyset.$$

Arguing as in the last part of the proof of Proposition 4.1 we produce a measurable $v \in \Phi(u)$ such that $v(t) \in T_C(u(t))$ for all $t \in [0, T]$. Then again by [4, Cor. 8.5.2], $v \in T_K(j(u))$. \square

According to Theorem 3.4, there is $u \in D(A)$ such that $-Au \in \Phi(u)$. This implies

Theorem 4.10. *Under the above assumptions the problem (4.13) has a strong solution $u \in K$.* \square

Corollary 4.11. *Suppose that $f : [0, T] \times \mathbb{R}^N \times \mathbb{R}^N \rightarrow \mathbb{R}^N$ is continuous such that*

- (i) *there are $c \in \mathbb{R}$ and $R \geq 0$ such that $u \cdot f(t, u, v) \leq cR^2$ for $t \in [0, T]$, $|u| = R$ and $v \in \mathbb{R}^N$ with $u \cdot v = 0$ ⁽¹²⁾;*
- (ii) *f is bounded on the strip of the form $[0, T] \times D(0, M) \times \mathbb{R}^N$;*
- (iii) *$p : [0, T] \rightarrow \mathbb{R}$ is C^1 -smooth and positive, $q : [0, T] \rightarrow \mathbb{R}$ is continuous and $q \geq c$.*

Then the problem $-(pu')' - q(t)u = f(t, u, u')$ subject to Dirichlet ($u(0) = u(T) = 0$), Neumann ($u'(0) = u'(T) = 0$) or periodic ($u(0) = u(T)$, $u'(0) = u'(T)$) boundary conditions has a solution in $C^2([0, T], \mathbb{R})$ such that $|u(t)| \leq R$ for $0 \leq t \leq T$. \square

The reader will easily formulate analogous results for elliptic PDE or partial differential inclusions. For instance one can get the generalization of the classical concerning the existence of steady states of the heat equation $u_t - \Delta u = g(u)$ subject to the Dirichlet boundary condition, where a continuous g is such that for some positive K, C one has $ug(u) \leq C|u|^2$ for $|u| \leq K$.

4.4. Sub- and superharmonics; moving rectangles. In this section we will discuss the Dirichlet problem

$$(4.14) \quad -\Delta u \in H(x, u, \nabla u), \quad u|_{\partial\Omega} = 0.$$

Now we assume

¹²Observe that if $N = 1$, then this means that $u \cdot f(t, u, 0) \leq cR^2$.

(H₁) There are $\alpha, \beta \in C^1(\Omega, \mathbb{R}^N) \cap C(\bar{\Omega}, \mathbb{R}^N)$ such that $\alpha \leq \beta$ ⁽¹³⁾, α is a weakly sub- and β a weakly superharmonic, i.e. for all $\varphi \in C_0^\infty(\Omega)$, $\varphi \geq 0$,

$$\int_{\Omega} \Delta \varphi \alpha \, dx \geq 0, \quad \alpha|_{\partial\Omega} \leq 0 \quad \text{and} \quad \int_{\Omega} \Delta \varphi \beta \leq 0, \quad \beta|_{\partial\Omega} \geq 0;$$

(H₂) $H : C \times \mathbb{R}^{MN} \rightarrow \mathbb{R}^N$ is a bounded upper semicontinuous map with convex compact values, where

$$C := \{(x, u) \in \Omega \times \mathbb{R}^N \mid \alpha(x) \leq u \leq \beta(x)\}.$$

Note that, for any $x \in \Omega$, the section $C(x) := \{u \in \mathbb{R}^n \mid (x, u) \in C\}$ is a cube in \mathbb{R}^n , hence the cone $T_{C(x)}(u)$ is determined in Remark 2.4 (1). The set C maybe viewed as the graph of the *moving rectangles* $\Omega \ni x \mapsto C(x)$.

For any $i = 1, \dots, N$ let us introduce the lower and upper i -th ‘faces’ of C :

$$C_i(x) := \{u \in \mathbb{R}^n \mid (x, u) \in C, u_i = \alpha_i(x)\}, \quad C^i(x) := \{u \in \mathbb{R}^n \mid (x, u) \in C, u_i = \beta_i(x)\}.$$

Further we assume that for all $i = 1, \dots, N$ and $x \in \Omega$

(H₃) if $u \in C_i(x)$ and $v = (v_1, \dots, v_N) \in (\mathbb{R}^M)^N$ with $v_i = \nabla \alpha_i(x)$, then there is $y \in H(x, u, v)$ with $y_i \geq 0$;

(H₄) if $u \in C^i(x)$ and $v = (v_1, \dots, v_N) \in (\mathbb{R}^M)^N$ with $v_i = \nabla \beta_i(x)$, then there is $y \in H(x, u, v)$ with $y_i \leq 0$.

Let $p > M$ and

$$E := L^p(\Omega, \mathbb{R}^N), \quad E_0 = C^1(\Omega, \mathbb{R}^N) \cap C_0(\bar{\Omega}, \mathbb{R}^N);$$

$$K := \{u \in E \mid \alpha(x) \leq u(x) \leq \beta(x) \text{ a.e. on } \Omega\}, \quad K_0 = \{u \in E_0 \mid \alpha(x) \leq \beta(x) \text{ on } \Omega\};$$

$$D(A) = W_0^{1,p}(\Omega, \mathbb{R}^N) \cap W^{2,p}(\Omega, \mathbb{R}^N), \quad Au = \Delta u, \quad u \in D(A).$$

Theorem 4.12. *Under hypotheses (H₁) – (H₄) problem (4.14) has a strong solution in K .*

PROOF: It is clear that assumptions (A₁) – (A₃) (with $\omega = 0$) from section 3.1 are satisfied. Let us define $\Phi : K_0 \rightarrow E$ by

$$\Phi(u) = \{v \in E \mid v(x) \in H(x, u(x), \nabla u(x)) \text{ for } x \in \Omega\}.$$

As in Proposition 4.1 we check that assumption (A₄) is also verified. Moreover (A₇) holds true since $p > M$ and, thus, the inclusion $D(A) \rightarrow E_0$ is compact. We will check that (A₅) and (A₆) are true.

Take $u \in K_0$ and let $X_i := \{x \in \Omega \mid u(x) \in C_i(x)\}$, $X^i := \{x \in \Omega \mid u(x) \in C^i(x)\}$, $i = 1, \dots, N$. If $x \in X_i$ for some i , then $u_i(x) = \alpha_i(x)$ and $\nabla u_i(x) = \nabla \alpha_i(x)$ since $u_i - \alpha_i$ attains a minimum at x ; similarly if $x \in X^i$, then $u_i = \beta_i(x)$ and $\nabla u_i(x) = \nabla \beta_i(x)$. Hence, by (H₃) and (H₄), if $x \in \bigcup_{i=1}^N (X_i \cup X^i)$, then $H(x, u(x), \nabla u(x)) \cap T_{C(x)}(u(x)) \neq \emptyset$. Otherwise, if $\alpha_i < u(x) < \beta_i(x)$ for all $i = 1, \dots, N$, then $T_{C(x)}(u(x)) = \mathbb{R}^N$. Taking into account [4, Cor. 8.5.2] we see that $w \in T_K(j(u))$ if and only $w(x) \in T_{C(x)}(u(x))$ for a.a. $x \in \Omega$. Arguing as in the last part of the proof of Proposition 4.1 we see that (A₅) is satisfied.

Condition (A₆) follows implicitly from [38, Th. 16]. Since we are in a special situation let us show

¹³Here and below inequalities between vectors are understood in the componentwise sense.

a simple argument. Using (H_1) and the density arguments we see that if $v \in H_0^1(\Omega)$ and $v \geq 0$ a.e. on Ω , then for any $i = 1, \dots, N$,

$$\int_{\Omega} \nabla v \cdot \nabla \alpha_i dx \leq 0 \quad \text{and} \quad \int_{\Omega} \nabla v \cdot \nabla \beta_i \geq 0.$$

Let $f \in K$, i.e. $\alpha \leq f \leq \beta$, and $u = J_h(f)$, where $h > 0$. Then for any $v \in H_0^1(\Omega)$, $v \geq 0$,

$$\int_{\Omega} (u_i - \beta_i)v dx + h \int_{\Omega} (\nabla u - \nabla \beta_i) \cdot \nabla v dx = \int_{\Omega} (f_i - \beta_i)v dx - h \int_{\Omega} \nabla v \cdot \nabla \beta_i dx \leq 0.$$

Taking $v := (u_i - \beta_i)_+$, we see that $v \in H_0^1(\Omega)$ since $\beta_i \geq 0$. Hence

$$\int_{\Omega} |v|^2 dx \leq -h \int_{\{u_i > \beta_i\}} |\nabla(u_i - \beta_i)|^2 dx \leq 0$$

since $\nabla v = \chi_{\{u_i > \beta_i\}} \nabla(u_i - \beta_i)$. Thus $u_i \leq \beta_i$. Analogously we prove that $u_i \geq \alpha$ on Ω for all $i = 1, \dots, N$. This means that $u \in K \cap D(A) \subset K_0$. Applying Theorem 3.4 we end the proof. \square

Remark 4.13. The existence of solutions to (4.14) may be established by a direct use of arguments employed in the proof of Theorem 3.4 since in this situation we can use some particular issues present in the problem. Given $i = 1, \dots, N$ and $u \in W_0^{1,p}(\Omega, \mathbb{R}^N)$ let

$$\pi_i(u) = \begin{cases} \beta_i & \text{if } u_i > \beta_i, \\ u_i & \text{if } \alpha_i \leq u_i \leq \beta_i, \\ \alpha_i & \text{if } u_i < \alpha_i. \end{cases}$$

One can show that $\pi_i : \mathcal{E} := W_0^{1,p}(\Omega, \mathbb{R}^N) \rightarrow W_0^{1,p}(\Omega)$ is well-defined and continuous. The map $\pi := (\pi_1, \dots, \pi_N) : \mathcal{E} \rightarrow \mathcal{E}$ is a retraction of \mathcal{E} onto $K \cap \mathcal{E}$. Note that $E_0 \hookrightarrow \mathcal{E}$ and let

$$\Psi(u) := \Phi(\pi(u)), \quad u \in E_0.$$

Taking into account that $0 \in \rho(\Delta)$ we may consider the composition

$$\xi : E_0 \xrightarrow{\Psi} E \xrightarrow{(-\Delta)^{-1}} D(A) \rightarrow E_0.$$

This composition is a compact (at large: i.e. the range of ξ is relatively compact) upper-semicontinuous map with compact convex values. By the Glicksberg-Fan theorem (the set-valued version of the Schauder fixed point principle) we gather that ξ has a fixed point $u \in E_0$. Using (H_3) and (H_4) and the maximum principle one shows that u is located in K and, therefore, is a solution to (4.14).

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