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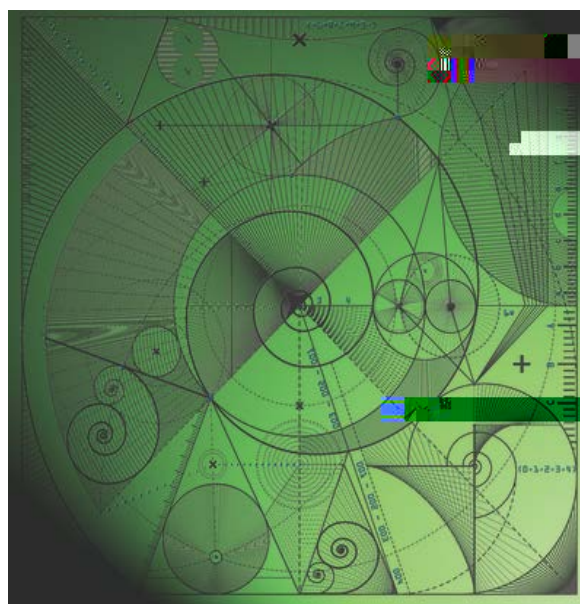
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Equivalence between weak and -solutions for symmetric hyperbolic first order PDE systems

by

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symmetric hyperbolic system (1.1) plays an important role in many contexts for both theory and applications.

In this paper we consider the question of equivalence between two completely different notions of generalised solutions to (1.1) and by assuming a commutativity

nonlinear degenerate elliptic systems

$$\begin{cases} F(x; D^2u) = f; & \text{in } \Omega; \\ u = 0; & \text{on } \partial\Omega; \end{cases}$$

and we proved existence and uniqueness of D -solution to the problem (which in general may not be even once weakly differentiable).

2. Theory of D -solutions for fully nonlinear systems

2.1. Preliminaries. We begin with some notation and some basic facts which will be used throughout the paper, perhaps without explicitly quoting this subsection.

Basics. Our measure theoretic and function space notation is either standard, e.g. as in [E2, EG] or self-explanatory. The norms $\| \cdot \|_j$ appearing will always be the Euclidean, while the Euclidean inner products will be denoted by either $\langle \cdot, \cdot \rangle$ on $\mathbb{R}^m; \mathbb{R}^M$ or by $\langle \cdot, \cdot \rangle$ on matrix spaces, e.g. on $\mathbb{R}^{M \times m}$ we have

$$\|X\|_j^2 = \sum_{i=1}^M \sum_{j=1}^m X_{ij}^2$$

Further, let $u : \mathbb{R}^m \rightarrow \mathbb{R}^M$ be any measurable map which we understand to be extended by zero on $\mathbb{R}^m \setminus \Omega$. For any $a \in \mathbb{R}^n$ with $|a| = 1$ and $h \in \mathbb{R} \setminus \{0\}$, the difference quotients of u along the direction a at x will be denoted by

$$(2.3) \quad D_a^{1;h} u(x) := \frac{1}{h} (u(x+ha) - u(x))$$

Definition 1 (Young measures) The space of Young measures is the set of all probability-valued mappings $E \in \mathbb{R}^m \mapsto P \in \overline{\mathbb{R}^{Mm}}$ which are weakly* measurable. Hence, the set of Young measures can be identified with a subset of the unit sphere of L_w^1 .

Definition 6 (D-solutions of 1st order systems) Let $\Omega \subset \mathbb{R}^m$ be open,

$$F : \mathbb{R}^M \times \mathbb{R}^{Mm} \rightarrow \mathbb{R}^d$$

a Carathéodory map and $u : \Omega \rightarrow \mathbb{R}^M$ a map in $W_{loc}^{1,1}(\Omega; \mathbb{R}^M)$. Suppose we have fixed some reference frames as in (2.1) and consider the PDE system

$$(2.5) \quad F(x; u(x); Du(x)) = 0; \quad \text{on } \Omega$$

We say that u is a D-solution of (2.5) when for any diffuse gradient of u arising from any infinitesimal sequence along subsequences (Definition 4)

$$D^{1,h_k} u \rightharpoonup^* Du \quad \text{in } Y; \bar{\mathbb{R}}^{Mm}; \quad \text{ask } 1$$

and for any $\varphi \in C_c^0(\mathbb{R}^{Mm})$, we have

$$\int_{\bar{\mathbb{R}}^{Mm}} (\varphi(X) F(x; u(x); X) - d[Du(x)](X)) = 0; \quad \text{a.e. } x \in \Omega$$

The following result asserts the fairly obvious fact that D-solutions are compatible with strong solutions.

Proposition 7 (Compatibility of D-solutions with strong solutions). Let F be a Carathéodory map as above, $u : \Omega \rightarrow \mathbb{R}^M$ a mapping in $W_{loc}^{1,1}(\Omega; \mathbb{R}^M)$ and consider the PDE system (2.5). Then, u is a D-solution on Ω if and only if u is a strong a.e. solution on Ω .

The proof of Proposition 7 is an immediate consequence of Lemma 5 and of the motivation of the notions.

Remark 8 (Nonlinearity of diffuse derivatives and relation to distributions). We summarise here some of the discussions of [K8]. In the context of the usual notions of solution (smooth, strong, weak, distributional), it is standard that the generalised derivative is a linear operation. However, without extra assumptions this may be false for diffuse derivatives; D-solutions are a genuinely nonlinear approach even when we apply them to linear PDE. More precisely, let $T_a : \mathbb{R}^{Mm} \rightarrow \mathbb{R}^{Mm}$ denote the translation by a . Given a Young measure $\mu \in \mathcal{Y}(\Omega; \bar{\mathbb{R}}^{Mm})$, we define $\mu \# T_a \in \mathcal{Y}(\Omega; \bar{\mathbb{R}}^{Mm})$;

3. Equivalence between weak and D -solutions for hyperbolic systems

3.1. Fibre spaces and the main result.

and for any $\underline{X} \in C_c^0(\mathbb{R}^{N+Nn})$, we have

$$\int_{\mathbb{R}^{N+Nn}} (\underline{X})_0 + A : \underline{X} \quad f(t; x) \, d \underline{Du}(t; x) (\underline{X}) = 0;$$

for a.e. $(t; x) \in (0; T) \times \mathbb{R}^n$.

Moreover, any D-solution u to (3.7) (and hence any weak solution) has the following regularity: the projection of the space-time gradient $(D_t u; Du)$ on the subspace $\mathbb{R}^{N+Nn}$ associated to A exists in L^2 , is spanned by rank-one matrices and for any such direction $\underline{a} \in \mathbb{R}^{N+Nn}$, we have $D_{\underline{a}} u \in L^2((0; T) \times \mathbb{R}^n)$.

The commutativity hypothesis is always satisfied if either $n = 1$ (one spatial dimension) or $N = 1$ (scalar case). By standard results on hyperbolic systems, we readily have the following consequence of Theorem 9:

Corollary 10 (Well-posedness of the Cauchy problem for D-solutions). In the setting of Theorem 9 and under the same assumptions, the Cauchy problem

$$\begin{cases} D_t u + A : Du = f; & \text{in } (0; T) \times \mathbb{R}^n; \\ u(0; \cdot) = u_0; & \text{on } \{0\} \times \mathbb{R}^n; \end{cases}$$

has a unique D-solution in the Brezis space $W^{1,2}_{\text{loc}}(0; T) \times \mathbb{R}^n; \mathbb{R}^N$, for any given $f \in L^2(0; T) \times \mathbb{R}^n; \mathbb{R}^N$, $u_0 \in L^2(\mathbb{R}^n; \mathbb{R}^N)$ and $T > 0$.

Proof of Theorem 9. The proof consists of three lemmas. In the first one below we show that the commutativity hypothesis on the matrices $A_1, \dots, A_n$ implies that the vector space of (3.3) has an orthonormal basis of rank-one directions which can be completed to an orthonormal basis of rank-one directions spanning $\mathbb{R}^{N+Nn}$.

Lemma 11. Given the consequence of

Proof of Lemma 11. We begin by observing that directly from the definitions of $\mathcal{N}(\underline{A})$ and of $\mathcal{N}(\underline{A})$, we have

$$(3.13) \quad \begin{aligned} \mathcal{N}(\underline{A}) &= \{ \underline{X} \in \mathbb{R}^{Nn} : \underline{A} : \underline{X} = 0 \} \\ \mathcal{N}(\underline{A}) &= \{ \underline{Y} = [Y_0]Y \in \mathbb{R}^{N+Nn} : Y_0(\underline{A} : \underline{X}) + Y : \underline{X} = 0, \underline{X} \in \mathbb{R}^{Nn} \} \end{aligned}$$

Next, by standard linear algebra results ([L]) we obtain that the commutativity hypothesis of the (symmetric) matrices $A_1, \dots, A_n \in \mathbb{R}^{N \times N}$ is equivalent to the requirement that there exists an orthonormal basis $e^1, \dots, e^N \in \mathbb{R}^N$ which diagonalises all the matrices $A_1, \dots, A_n$ simultaneously, namely there is a common set of eigenvectors for perhaps different eigenvalues $c^{(i)1}, \dots, c^{(i)N}$ of A_i . Thus, for any $i = 1, \dots, n$, we have

$$A_i e^j = c^{(i)j} e^j; \quad j = 1, \dots, N;$$

or, in index form (see (1.4))

$$\sum_{j=1}^N A_{ij} e^j = c^{(i)j} e^j; \quad j = 1, \dots, N;$$

whereas $(A_i) = f(c^{(i)1}, \dots, c^{(i)N})g$. We rewrite the above as

$$\sum_{j=1}^N A_{ij} e^j + c^{(i)} = 0; \quad j = 1, \dots, N; \quad i = 1, \dots, n;$$

and in view of (1.4) we may write it as

$$(3.14) \quad A : e^j + c^{(i)} = 0 \quad j = 1, \dots, N; \quad i = 1, \dots, n;$$

We now define

$$(3.15) \quad N^i := \frac{c^{(i)}}{e^j} = \begin{pmatrix} c^{(i)} & 0 & \dots & 0 & 1 & 0 & \dots & 0 \end{pmatrix}^{\text{position } (1+i)};$$

for $j = 1, \dots, N$, $i = 1, \dots, n$, and also

$$(3.16) \quad N^0 := \frac{1}{c} = \begin{pmatrix} 1 & c^{(1)} & \dots & c^{(n)} \end{pmatrix};$$

where

$$c := c^{(1)}; \dots; c^{(n)}$$

is the i -th eigenvalue vector of the matrices $A_1, \dots, A_n$. The definition of N^i and (3.14) with (3.1) immediately give that

$$\underline{A} : N^i = 0; \quad i = 1, \dots, n;$$

and hence $N^i \in \mathcal{N}(\underline{A})$. Moreover, by (3.13) and the fact that the (Nn) -many matrices $f(e^j); j = 1, \dots, N$ are an orthonormal basis of $\mathbb{R}^{Nn}$, we have that

$$\begin{aligned} \underline{Y} \in \mathcal{N}(\underline{A}) &\iff Y_0(\underline{A} : \underline{X}) + Y : \underline{X} = 0; \quad \underline{X} \in \mathbb{R}^{Nn}; \\ &\iff Y_0 : (\sum_{j=1}^N e^j) + Y : (\sum_{j=1}^N e^j) = 0; \quad j = 1, \dots, N; \quad i = 1, \dots, n; \\ (3.14) \quad &\iff Y_0 : c^{(i)} + Y : (e^j) = 0; \quad j = 1, \dots, N; \quad i = 1, \dots, n; \\ (3.15) \quad &\iff [Y_0]Y : N^i = 0; \quad i = 1, \dots, n; \end{aligned}$$

Hence, $\underline{Y} \in N(\underline{A})$ if and only if $\underline{Y} \in N^i$ for all $i = 1, \dots, N$ and $i = 1, \dots, n$. Since $N(\underline{A}) = \{ \underline{Y} \}$, this proves that

$$(3.17) \quad N(\underline{A}) = \text{span} \left[\left\{ N^i \mid i = 1, \dots, N; i = 1, \dots, n \right\} \right];$$

Moreover, the matrices N^i spanning $N(\underline{A})$ are linearly independent and hence exactly Nn -many. Indeed, for $i, j = 1, \dots, N$ and $i, j = 1, \dots, n$, by (3.15) we have

$$\begin{aligned} N^i : N^j &= \frac{c^{(i)}}{e^j} : \frac{c^{(j)}}{e^i} \\ &= c^{(i)} c^{(j)} + \delta_{ij} : \end{aligned}$$

It follows that for any $i \neq j$, N^i is orthogonal to N^j . Moreover, for any $i = 1, \dots, N$ and $i \neq j$ in $\{1, \dots, n\}$, by (3.15) we have

$$\frac{N^i}{\|N^i\|} : \frac{N^j}{\|N^j\|} = \frac{c^{(i)} c^{(j)}}{1 + (c^{(i)})^2} \frac{c^{(j)}}{1 + (c^{(j)})^2} = 2(c^{(i)} + 1)$$

and hence for each i the set of matrices $\{N^i\}$ is linearly independent. Further, by (3.15), (3.16) we have that

$$\begin{aligned} N^0 : N^i &= \frac{1}{c} : \frac{c^{(i)}}{e^i} \\ &= \frac{1}{c} : \frac{c^{(i)}}{1 + c^{(i)}} \\ &= \frac{1}{c} : \frac{c^{(i)}}{1 + c^{(i)}} \\ &= 0; \end{aligned}$$

for all $i = 1, \dots, N$ and $i = 1, \dots, n$. Moreover, by (3.16) we have

$$\begin{aligned} N^0 : N^0 &= \frac{1}{c} : \frac{1}{c} \\ &= 1 + c : c \end{aligned}$$

and as a consequence the matrices $\{N^0\}$ form an orthogonal set of N -many elements which is orthogonal to $N(\underline{A})$. Since the dimension of the space is $n + Nn$, all the above together with (3.13), (3.15), (3.16) prove that

$$(3.18) \quad N(\underline{A}) = \text{span} \left[\left\{ N^i \mid i = 1, \dots, N; i = 1, \dots, n \right\} \right];$$

We now show that the frame $\{N^i\}$ can be modified in order to be made an orthonormal basis and still consisting of rank-one matrices. First note that the matrices spanning $N(\underline{A})$ are orthogonal and we only need to fix their length. Further, note that $\mathbb{R}^n$ can be decomposed as the following direct sum of mutually orthogonal subspaces

$$\mathbb{R}^n = \bigoplus_{i=1}^N \text{span} \left[\left\{ N^i \mid i = 1, \dots, n \right\} \right] =: \bigoplus_{i=1}^N W_i$$

Since

$$W_i = \text{span} \left[\left\{ \frac{c^{(i)}}{e^j} : j = 1, \dots, n \right\} \right];$$

by the Gram-Schmidt method, we can find an orthonormal basis of W consisting of matrices of the form

$$(3.19) \quad N^i = \underline{a}^{(i)}; \quad \underline{a}^{(i)} \cdot \underline{a}^{(j)} = \delta_{ij}.$$

Finally, we define

$$E^0 := \frac{N^0}{\|N^0\|} = \rho \frac{1}{1 + |c|^2} \frac{1}{c} \in \mathbb{R}^{N+Nn};$$

$$E^i := N^i = \underline{a}^{(i)} \in \mathbb{R}^{N+Nn};$$

and also

$$E := \begin{pmatrix} E^0 \\ E^i \end{pmatrix} \in \mathbb{R}^N;$$

$$E^{(0)} := \rho \frac{1}{1 + |c|^2} \frac{1}{c}; \quad E^{(i)} := \underline{a}^{(i)} \in \mathbb{R}^{1+n};$$

where $j = 1, \dots, N$ and $i = 1, \dots, n$. By the previous it follows that $\{E^i\}_{i=1}^n = \{1, \dots, N; i = 0; 1, \dots, n\}$ is an orthonormal basis of $\mathbb{R}^{N+Nn}$ consisting of rank-one directions such that $\{E^0\}_{j=1}^N = \{1, \dots, N\}$ span the subspace and $\{E^i\}_{i=1}^n = \{1, \dots, N; i = 1, \dots, n\}$ span its complement. Moreover, $E^i = E - E^{(i)}$. We conclude the proof of the lemma by noting that (3.11), (3.12) follow by the definition of $\tilde{u}$ and standard linear algebra results.

Next, we employ the orthonormal frames constructed in Lemma 11 and the properties (3.11), (3.12) of A in order to characterise weak solutions to (3.7) as mappings in the Brezis space (3.6) which solve the equation in a pointwise "strong Brezis-wise" sense: the equation is satisfied a.e. on $(0, T) \times \mathbb{R}^n$ if we substitute the distributional gradient $(D_t u; Du)$

that $u \in W^{1,2}(0;T) \times \mathbb{R}^n; \mathbb{R}^N$ and in addition $\underline{D}u \rightharpoonup \underline{G}(u)$ in L^2 . Thus, by passing to the limit in (3.20) as $\epsilon \rightarrow 0$ and as $\epsilon \rightarrow 0$, we obtain that

$$(3.21) \quad \underline{A} : \underline{G}(u) = f; \quad \text{a.e. on } (0;T) \times \mathbb{R}^n;$$

as desired. Conversely, suppose that (3.21) holds. Then, by the definition (3.6) of the Bre space there are approximating sequences $u_\epsilon \rightarrow u$ and $\underline{D}u_\epsilon \rightarrow \underline{G}(u)$, both in L^2 as $\epsilon \rightarrow 0$. Hence, we have

$$\begin{aligned} \underline{A} : \underline{D}u_\epsilon - f &= \underline{A} : \underline{D}u_\epsilon - \underline{G}(u) \\ &= o(1); \end{aligned}$$

as $\epsilon \rightarrow 0$, in L^2 . By the above, (3.12) and (3.1)-(3.5), we have

$$\begin{aligned} D_t u_\epsilon + \underline{A} : \underline{D}u_\epsilon - f &= \underline{A} : \underline{D}u_\epsilon - f \\ &= \underline{A} : \underline{D}u_\epsilon - f \\ &= o(1); \end{aligned}$$

as $\epsilon \rightarrow 0$, in L^2_Z . Hence, for any $\varphi \in C_c^1(0;T) \times \mathbb{R}^n$, we have

$$(0;T) \times \mathbb{R}^n \quad \underline{D}u$$

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Further, for any fixed measurable set $E \subset (0; T) \times \mathbb{R}^n$ with finite measure and any $f \in C_c^0(\mathbb{R}^{N+Nn})$, by using our hypothesis $\underline{A} : \underline{G}(u) = f$, we have the estimate

$$(3.24) \quad \|\underline{D}^{1;h}u - \underline{A} : \underline{D}^{1;h}u - f\|_{L^1(E)} \leq \|\underline{D}^{1;h}u - \underline{A} : \underline{D}^{1;h}u - \underline{A} : \underline{G}(u)\|_{L^2((0;T) \times \mathbb{R}^n)} \leq \|\underline{D}^{1;h}u - \underline{A} : \underline{G}(u)\|_{L^2((0;T) \times \mathbb{R}^n)}.$$

Hence, (3.23) and (3.24) imply

$$(3.25) \quad \|\underline{D}^{1;h}u - \underline{A} : \underline{D}^{1;h}u - f\|_{L^1(E)} \rightarrow 0 \text{ in } L^1(E; \mathbb{R}^N) \text{ as } h \rightarrow 0.$$

Moreover, the Carathéodory function

$$(3.26) \quad \underline{f}(\underline{x}; \underline{X}) := \underline{X} : \underline{A} : \underline{X} - f(\underline{x}) \in \mathbb{R}$$

is an element of the space

$$L^1(0; T) \times \mathbb{R}^n; C^0(\overline{\mathbb{R}^{N+Nn}})$$

because

$$\|\underline{f}\|_{L^1(0; T) \times \mathbb{R}^n; C^0(\overline{\mathbb{R}^{N+Nn}})} \leq \|\underline{f}\|_{L^1(0; T) \times \mathbb{R}^n; C^0(\overline{\mathbb{R}^{N+Nn}})} + \|\underline{f}\|_{L^1(0; T) \times \mathbb{R}^n; C^0(\overline{\mathbb{R}^{N+Nn}})} \leq \|\underline{f}\|_{L^1(0; T) \times \mathbb{R}^n; C^0(\overline{\mathbb{R}^{N+Nn}})} + \|\underline{f}\|_{L^1(0; T) \times \mathbb{R}^n; C^0(\overline{\mathbb{R}^{N+Nn}})}.$$

Let now $(h_k)_{k=1}^\infty \subset \mathbb{R}^n$ be any infinitesimal sequence. Then, there is a subsequence $h_{k_j} \rightarrow 0$ such that

$$\|\underline{D}^{1;h_{k_j}}u - \underline{A} : \underline{D}^{1;h_{k_j}}u - f\|_{L^1(E)} \rightarrow 0 \text{ in } Y(0; T) \times \mathbb{R}^n; \overline{\mathbb{R}^{N+Nn}}; \text{ as } k_j \rightarrow \infty.$$

By strong continuity of the duality pairing between

$$\mathbb{R}^n; C^0$$

Conversely, suppose that

Such selections exist for large enough $R > 0$ by Aumann's measurable selection theorem (see e.g. [FL]), but in this specific case they can also be constructed explicitly because of the simple structure of the multi-valued mapping. By using (3.34), (3.33) implies that

$$\lim_{k \rightarrow 1} \int_{E^R} \underline{A} : T^R \underline{x}; \underline{D}^{1,h}_k u(\underline{x}) \quad f(\underline{x}) \, d\underline{x} = 0$$

and by recalling (3.12), we rewrite this as

$$(3.35) \quad \lim_{k \rightarrow 1} \int_{E^R} \underline{A} : T^R \underline{x}; \underline{D}^{1,h}_k u(\underline{x}) \quad f(\underline{x}) \, d\underline{x} = 0:$$

Hence, (3.35) implies that

$$\begin{aligned} \int_{E^R} \underline{A} : \underline{G}(u) \quad f &= \int_{E^R} \underline{A} : T^R \quad ; \quad \underline{D}^{1,h}_k u \quad f \\ &+ \int_{E^R} \underline{A} : T^R \quad ; \quad \underline{D}^{1,h}_k u \quad \underline{A} : \underline{G}(u) \\ &\alpha(1) + \int_{E^R} \underline{A} : T^R \quad ; \quad \underline{D}^{1,h}_k u \quad \underline{G}(u) \end{aligned}$$

as $k \rightarrow 1$, and as a consequence we have

$$(3.36) \quad \int_{E^R} \underline{A} : \underline{G}(u) \quad f = \int_{E^R} \underline{A} : T^R \quad ; \quad \underline{D}^{1,h}_k u \quad T^R \quad ; \underline{G}(u) \\ + \int_{E^R} \underline{A} : T^R \quad ; \underline{G}(u) \quad \underline{G}(u) + \alpha(1);$$

as $k \rightarrow 1$, for any $R > 0$. Moreover, by assumption u is in the Bre space (3.6). Hence by invoking (3.23), the Dominated convergence theorem, the fact that $|E| < 1$ and (3.34), we may pass to the limit in (3.36) as $k \rightarrow 1$ to obtain

$$\int_{E^R} \underline{A} : \underline{G}(u) \quad f = \int_{E^R} \underline{A} : T^R \quad ; \underline{G}(u) \quad \underline{G}(u);$$

for any $R > 0$. Finally, we let $R \rightarrow 1$ and recall the arbitrariness of the set $E \subset (0; T) \times \mathbb{R}^n$ and (3.34) to infer that $\underline{A} : \underline{G}(u) = f$, a.e. on $(Q; T) \times \mathbb{R}^n$. The lemma has been established.

The proof of Theorem 9 is now complete.

Remark 14 (Functional representation of the diuse gradients). In a sense, Lemma 13 says that all the diuse gradients of the D-solution u when restricted on the subspace of non-degeneracies have a certain "functional" representation inside the coefficients, given by $\underline{G}(u)$. Namely, if we decompose $\mathbb{R}^{N+Nn} = \mathbb{R}^N \oplus \mathbb{R}^{Nn}$, the restriction of any diuse space-time gradient $\underline{D}u \in Y \times \mathbb{R}^{N+Nn}$ on $\mathbb{R}^n$ is given by the Bre space-time gradient:

$$\underline{D}u(t; x) X = \underline{G}(u)(t; x); \quad \text{a.e. } (t; x) \in (0; T) \times \mathbb{R}^n:$$

This is a statement of "partial regularity type" for D-solutions: although not all of the diuse gradient is a Dirac mass, certain restrictions of it on subspaces are concentration measures.

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References

- A1. G. Aronsson, Minimization problems for the functional $\sup_x F(x; f(x); f^0(x))$, Arkiv för Mat. 6 (1965), 33 - 53.
- A2. G. Aronsson, Minimization problems for the functional $\sup_x F(x; f(x); f^0(x))$ II, Arkiv för Mat. 6 (1966), 409 - 431.
- A3. G. Aronsson, Extension of functions satisfying Lipschitz conditions, Arkiv för Mat. 6 (1967), 551 - 561.
- A4. G. Aronsson, On the partial differential equation $u_x^2 u_{xx} + 2 u_x u_y u_{xy} + u_y^2 u_{yy} = 0$, Arkiv för Mat. 7 (1968), 395 - 425.
- A5. G. Aronsson, Minimization problems for the functional $\sup_x F(x; f(x); f^0(x))$ III, Arkiv för Mat. (1969), 509 - 512.
- CFV. C. Castaing, P. R. de Fitte, M. Valadier, Young Measures on Topological spaces with Applications in Control Theory and Probability Theory, Mathematics and Its Applications, Kluwer Academic Publishers, 2004.
- Co. J.F. Colombeau, New Generalized Functions and Multiplication of distributions, North Holland, 1983.
- D. B. Dacorogna, Direct Methods in the Calculus of Variations, 2nd Edition, Volume 78, Applied Mathematical Sciences, Springer, 2008.
- DM. B. Dacorogna, P. Marcellini, Implicit Partial Differential Equations, Progress in Nonlinear Differential Equations and Their Applications, Birkhäuser, 1999.
- DPM. R.J. DiPerna, A.J. Majda, Oscillations and concentrations in weak solutions of the incompressible fluid equations, Commun. Math. Phys. 108, 667 - 689 (1987).
- Ed. R.E. Edwards, Functional Analysis: Theory and Applications, Dover Books on Mathematics, 2003.
- E. L.C. Evans, Weak convergence methods for nonlinear partial differential equations, Regional conference series in mathematics 74, AMS, 1990.
- E2. L.C. Evans, Partial Differential Equations, AMS, Graduate Studies in Mathematics Vol. 19, 1998.
- EG. L.C. Evans, R. Gariepy, Measure theory and the properties of functions, Studies in advanced mathematics, CRC press, 1992.
- FG. L.C. Florescu, C. Godet-Thobie, Young measures and compactness in metric spaces, De Gruyter, 2012.
- FL. I. Fonseca, G. Leoni, Modern methods in the Calculus of Variations: L^p spaces Springer Monographs in Mathematics, 2007.
- K. N. Katzourakis, An Introduction to viscosity Solutions for Fully Nonlinear PDE with Applications to Calculus of Variations in L^1 , Springer Briefs in Mathematics, 2015, DOI 10.1007/978-3-319-12829-0.
- K1. N. Katzourakis, L^1 -Variational Problems for Maps and the Aronsson PDE system, J. Differential Equations, Volume 253, Issue 7 (2012), 2123 - 2139.
- K2. N. Katzourakis, Explicit 2D L^1 -Harmonic Maps whose Interfaces have Junctions and Corners, Comptes Rendus Acad. Sci. Paris, Ser.I, 351 (2013) 677 - 680.
- K3. N. Katzourakis, On the Structure of L^1 -Harmonic Maps, Communications in PDE, Volume 39, Issue 11 (2014), 2091 - 2124.
- K4. N. Katzourakis, L^1 -Minimal Submanifolds, Proceedings of the Amer. Math. Soc., 142 (2014) 2797-2811.
- K5. N. Katzourakis, Nonuniqueness in Vector-valued Calculus of Variations in L^1 and some Linear Elliptic Systems, Communications on Pure and Applied Analysis, Vol. 14, 1, 313 - 327 (2015).
- K6. N. Katzourakis, Optimal L^1 -Quasiconformal Immersions, ESAIM Control, Opt. and Calc. Var., to appear (2015) DOI: <http://dx.doi.org/10.1051/cocv/2014038>.
- K7. N. Katzourakis, On Linear Degenerate Elliptic PDE Systems with Constant Coefficients, Adv. in Calculus of Variations, DOI: 10.1515/acv-2015-0004, published online June 2015.
- K8. N. Katzourakis, Generalised solutions for fully nonlinear PDE systems and existence-uniqueness theorems, ArXiv preprint, <http://arxiv.org/pdf/1501.06164.pdf>.
- K9. N. Katzourakis, Existence of generalised solutions to the equations of vectorial Calculus of Variations in L^1 , ArXiv preprint, <http://arxiv.org/pdf/1502.01179.pdf>.
- L. P. D. Lax, Linear Algebra and Its Applications, Wiley-Interscience, 2nd edition, 2007.

- KR. J. Kristensen, F. Rindler, Characterization of generalized gradient Young measures generated by sequences in $W^{1,1}$ and BV, Arch. Rational Mech. Anal. 197, 539 - 598 (2010) and erratum Arch. Rational Mech. Anal. 203, 693 - 700 (2012).
- M. S. Müller, Variational models for microstructure and phase transitions, Lecture Notes in Mathematics 1783, Springer, 85-210, 1999.
- P. P. Pedregal, Parametrized Measures and Variational Principles, Birkhäuser, 1997.
- V. M. Valadier, Young measures, in "Methods of nonconvex analysis", Lecture Notes in Mathematics 1446, 152-188 (1990).
- Y. L.C. Young, Generalized curves and the existence of an attained absolute minimum in the calculus of variations, Comptes Rendus de la Societe des Sciences et des Lettres de Varsovie, Classe III 30, 212 - 234 (1937).

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