Inverse potential problems in divergence form: some uniqueness, separation and recovery issues

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Quasilinear equations, inverse problems and their applications, Dolgoprundy september 11-15, 2016 based in part on joint work with: S. Chevillard, J. Leblond (INRIA), D. Pei, Q. Tao (Macau).

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 For a complementation of the L1(Γ) form the Coupler integral.
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In particular, we have that

$$f(\xi) = \mathcal{C}^+ - \mathcal{C}^-.$$

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- In other words, if $\varphi : \mathbb{D} \to D^+$ is a conformal map from the unit disk and if we set $\Gamma_r = \varphi(|z| = r)$, then $f \in H^p(D^+)$ iff

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• Also, the Cauchy theorem holds:

$$0=\frac{1}{2i\pi}\int_{\Gamma}\frac{f(\xi)}{\xi-z}d\xi, \qquad f\in H^p(D^{\pm}), \qquad z\in D^{\mp}.$$

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Theorem

For Γ a smooth Jordan curve and 1 there holds a topological sum:

$$L^{p}(\Gamma) = H^{p}(D^{+}) \oplus H^{p}(D^{-}).$$

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Corollary

A \mathbb{R}^2 -valued vector field of L^p -class on Γ is uniquely the sum of the trace of the gradient of a harmonic function in D^+ and the trace of the gradient of a harmonic function in D^- , where both gradients have nontangential maximal function in L^p .

• Consider a 2-D potential in divergence form supported on Γ :

$$P_{\operatorname{div} V}(X) = -\frac{1}{2\pi} \int_{\Gamma} (\operatorname{div} V)(X') \log \frac{1}{|X - X'|} \, d|X'|, \quad X \notin \operatorname{supp} V,$$

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div *m* has null potential in D^- if, and only if $(m_1 + im_2)/v \in H^p(D^+)$.
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- The generalization of the previous Hardy decomposition stems from Hodge theory for currents supported on a surface in the ambient space, but it is conveniently framed in terms of Clifford analysis that we use here as a tool.

• For $n \ge 3$, we consider harmonic potentials in divergence form:

$$P_{\operatorname{div} V}(x) = \int \frac{x - y}{|x - y|^{n-2}} \operatorname{div} V(y)$$

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- They occur frequently when modeling electro-magnetic phenomena in the quasi-static approximation to Maxwell's equations.







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- Today, inverse magnetization problems are a hot topic in Earth and Planetary Sciences.

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- This geometry is in fact realistic in scanning microscopy of rocks which are typically sanded down to thin slabs.

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$$R_j(f)(Y) := \lim_{\epsilon \to 0} \frac{1}{2\pi} \iint_{\mathbb{R}^2 \setminus B(Y,\epsilon)} f(X') \frac{(y_j - x'_j)}{|Y - X'|^3} dX', \qquad j = 1, 2,$$

are the Riesz transforms.

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• Likewise *M* is silent from below iff

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 $R_1m_1 + R_2m_2 + m_3 = 0.$

• Likewise *M* is silent from below iff

 $R_1m_1 + R_2m_2 - m_3 = 0.$

• *M* is silent (from both sides) iff $R_1m_1 + R_2m_2 = 0$ and $m_3 = 0$.



What do these quantities mean?



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To approach it, we introduce some classical function spaces.

$$\sup_{x_3>0}\int_{\mathbb{R}^2}|\nabla u(X',x_3)|^p dX'<\infty.$$

• Let \mathfrak{H}_{+}^{p} consist of ∇u , u harmonic in $\{x_3 > 0\}$, such that

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- We put \mathcal{D}^{p} for divergence-free vector fields in $L^{p}(\mathbb{R}^{2}, \mathbb{R}^{2})$.

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For 1 one has the direct sum:

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• Easily checked using $R_1^2 + R_2^2 = -\text{Id.}$

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- *M* is silent iff it is tangent and divergence-free.
- Transparent if we observe the orthogonality:

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- If M ∈ (L²(ℝⁿ))³ then P_{55²} M yields the magnetization of least (L²(ℝⁿ))³-norm which is equivalent to M from above.

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- Define Sobolev spaces W^{1,p}(M) as usual, M inherits from Rⁿ a uniform Riemaniann structure ⟨.,.⟩_M, therefore one can define tangential gradient vector fields G^p, where L^p is understood with respect to the volume form σ.

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- Note the above nontangential limits are not tangent to \mathcal{M} .

Theorem

Let \mathcal{M} be a compact simply connected Lipschitz hypersurface in \mathbb{R}^n and $p_0(\mathcal{M}) . Then, there is a direct sum$

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- The result extends with obvious modifications to the case where *M* is not connected, and also to Lipschitz graphs.
- When \mathcal{M} is smooth the decomposition holds in more general spaces of functions or distributional currents.

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- Thus, we are left to decompose V − D − ∇u which is a normal vector field on M. For this we need preliminaries in Clifford analysis. We restrict to n = 3 for simplicity.
● C is the skew unital algebra generated over ℝ by {e₁, e₂, e₃} with:

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• A typical element of \mathfrak{C} is of the form

 $z = x_0 + x_1 e_1 + x_2 e_2 + x_3 e_3 + x_{1,2} e_1 e_2 + x_{2,3} e_2 e_3 + x_{1,3} e_1 e_3 + x_{123} e_1 e_2 e_3$

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- The conjugate of z is

 $\bar{z} = x_0 - x_1 e_1 - x_2 e_2 - x_3 e_3 + x_{1,2} e_1 e_2 + x_{2,3} e_2 e_3 + x_{3,1} e_3 e_1 - x_{123} e_1 e_2 e_3.$

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$$e_j^2 = -1, \quad e_i e_j = -e_j e_i.$$

 \bullet A typical element of ${\mathfrak C}$ is of the form

 $z = x_0 + x_1 e_1 + x_2 e_2 + x_3 e_3 + x_{1,2} e_1 e_2 + x_{2,3} e_2 e_3 + x_{1,3} e_1 e_3 + x_{123} e_1 e_2 e_3$

where the x_i , the $x_{k,\ell}$ and x_{123} are real numbers.

- x_0 is the scalar part of z, denoted by Sc z;
- x₁e₁ + x₂e₂ + x₃e₃, is the vector part of z denoted as vec z; Clifford vectors get identified with Euclidean vectors in ℝ³.
- The conjugate of z is

 $\bar{z} = x_0 - x_1 e_1 - x_2 e_2 - x_3 e_3 + x_{1,2} e_1 e_2 + x_{2,3} e_2 e_3 + x_{3,1} e_3 e_1 - x_{123} e_1 e_2 e_3.$

• The norm of z is $|z| = (\sum_{0 \le k \le 3} x_k^2 + \sum_{i < j} x_{i,j}^2 + x_{123}^2)^{1/2}$.

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Lemma

A vector-valued function is left monogenic if and only if it is monogenic, if and only if it is the gradient of a harmonic function.

• If f is left monogenic in Ω^+ and its nontangential maximal function lies in $L^p(\mathcal{M})$, then f has a nontangential limit $f^+ \in L^p(\mathcal{M})$ a.e. on \mathcal{M} (Verchota), and by the Green formula (see *e.g.* "Clifford Algebras and Dirac Operators in Analysis" by Gilbert and Murray):

$$f(z) = \mathcal{C}f^+(z) := rac{1}{4\pi}\int_{\mathcal{M}}rac{\overline{y-z}}{|y-z|^3}n(y)f^+(y)d\sigma(y), \qquad z\in \Omega^+.$$

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- If $z \in \Omega^-$, then the above right hand side is zero.
- A similar result holds if *f* is left monogenic in Ω⁻; the nontangential limit on *M* from Ω⁻ is denoted by *f⁻* ∈ *L^p*(*M*).

For a C-valued h ∈ L^p(M), Ch is left monogenic on ℝ³ \ M and its nontangential maximal function lies in L^p(M) [Coifman-McIntosh-Meyer, 1982].

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$$\mathcal{C}^{\pm}h(y) = \pm \frac{h(y)}{2} + \mathcal{SC}h(y), \qquad y \in \mathcal{M},$$

where *SCh* is the *singular Cauchy integral operator:*

$$\mathcal{SCh}(y) = rac{1}{4\pi} \lim_{arepsilon o 0} \int_{\mathcal{M} \setminus \mathcal{B}(y,arepsilon)} rac{\overline{\xi-y}}{|\xi-y|^3} \mathit{n}(\xi) \mathit{h}(\xi) \mathit{d}\sigma(\xi), \qquad y \in \mathcal{M}.$$

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• This gives us an analog of the Plemelj formula:

$$\mathcal{C}^+h(y)-\mathcal{C}^-h(y)=h(y).$$

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- Write h = C⁺h − C⁻h by Plemelj formula. Since h is normal, C[±]h ∈ H^p_±. Indeed, if h(y) is normal to M at y, the C-product n(y)h(y) is scalar-valued, so the integrand in the definition of Ch is vector valued.

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- Uniqueness follows from uniqueness of the Hodge decomposition and the Liouville theorem for harmonic functions.

• Assume $V = m \otimes \delta_{\mathcal{M}}$ where $m = (m_1, m_2, m_3)^t$ is a vector field in $L^p(\mathcal{M})$.

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Theorem

The distribution is silent from outside if and only if

$$2\pi\psi(y) = -\lim_{\varepsilon \to 0} \int_{\mathcal{M} \setminus B(y,\varepsilon)} \frac{\xi - y}{|\xi - y|^3} . (\psi n + R(D))(\xi) d\sigma(\xi) \quad y \in \mathcal{M}.$$

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This is a more complicated singular integral equation involving the curvature.

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• If $\operatorname{supp} m$ is a strict subset of \mathcal{M} , things get simple:

Corollary

If supp $m \neq M$, then m is silent from outside iff it is silent from inside, which is iff

 $m \in \mathcal{D}^p(\mathcal{M}).$

Theorem

Let $m \in (L^2(\mathcal{M})^3$ with supp $m \subset \Gamma_0 \neq \mathcal{M}$, and assume we know the field $B = \nabla P_{\operatorname{div} m}$ on a surface patch Σ disjoint from \mathcal{M} .

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The proof uses that ∇ψ ∈ (KerA)[⊥], where A maps m ∈ (L²(Γ₀))³ to B_{|Σ} ⊂ (L²(Σ))³.

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• Can be used to estimate moments or spherical harmonics expansions.

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- In Rⁿ, let rational approximation mean approximation by gradients of discrete harmonic potentials with finitely many masses. The Hardy-Hodge decomposition implies:

Theorem

Let *S* be a Lipschitz regular surface patch on a compact connected smooth hypersurface $\mathcal{M} \subset \mathbb{R}^n$. Let v be \mathbb{R}^n -valued in $L^p(S)$, $1 . Then, v can be approximated arbitrarily close by rationals in <math>L^p(S)$ iff the tangential component of v is a gradient.

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- A fortiori then, a compact set containing no such arc is a grad-set for L^p (any field is approximable y a gradient). The Hardy-Hodge decomposition now implies:

Theorem

Let K be a closed set in a compact connected smooth hypersurface $\mathcal{M} \subset \mathbb{R}^n$, and assume that K contains no simple rectifiable arc of positive length. Then, each \mathbb{R}^n -valued v in $L^p(K)$ can be approximated arbitrary close by rationals in $L^p(K)$, 1 .