## Hardy-Hodge decomposition of vector fields

L. Baratchart (INRIA) 25-th Summer Analysis meeting in Mathematical Analysis 25–30 June, 2016 honoring the memory of V.P. Havin

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

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

for some vector distribution  $V = (v_1, v_2, \dots, v_n)^t$  on  $\mathbb{R}^n$ .

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- They solve  $\Delta u = \operatorname{div} V$  on  $\mathbb{R}^n$  with "minimal growth" at infinity.
- They occur frequently when modeling electro-magnetic phenomena in the quasi-static approximation to Maxwell's equations.

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with  $J_p$  the so-called primary current.

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

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To approach it, we introduce some classical function spaces.

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

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

$$\sup_{\mathsf{x}_3>0}\int_{\mathbb{R}^2}|\nabla u(\mathsf{X}',\mathsf{x}_3)|^p d\mathsf{X}'<\infty.$$

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

#### Theorem (L.B., D. Hardin, E. Lima, E.B. Saff, B. Weiss)

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For 1 one has the orthogonal sum:

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where  $\mathcal{G}^p$  is the space of distributional gradients in  $L^p(\mathbb{R}^2, \mathbb{R}^2)$ .

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

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

$$\mathfrak{H}^p_+ \perp \mathfrak{H}^q_- \quad \text{and} \quad \mathcal{D}^p \times \{0\} \perp \mathfrak{H}^q_\pm, \quad 1/p + 1/q = 1.$$

# Easy extensions

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- It extends to any class of functions or of distributions invariant under Riesz transforms, e.g.  $\mathfrak{h}^1$ , BMO,  $W^{-\infty,p}$  (i.e. finite sums of derivatives of any order of  $L^p$ -functions, 1 ). The latter contains all distributions with compact support.
- If  $M \in (L^2(\mathbb{R}^n))^3$  then  $P_{\mathfrak{H}^2_-}M$  yields the magnetization of least  $(L^2(\mathbb{R}^n))^3$ -norm which is equivalent to M from above.

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- One can then define  $\mathcal{D}^p = (\mathcal{G}^q)^{\perp}$  for the pairing

$$(G,D):=\int_{\mathcal{M}}\langle G,D\rangle_{\mathcal{M}}d\sigma,\quad 1/p+1/q=1.$$

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#### Theorem

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

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- Thus, we are left to decompose  $V-D-\nabla u$  which is a normal vector field on  $\mathcal{M}$ . For this we need preliminaries in Clifford analysis. We restrict to n=3 for simplicity.

# Some Clifford analysis

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

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

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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  $\Omega^+$  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):

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

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$$\mathcal{SC}h(y) = \frac{1}{4\pi} \lim_{\varepsilon \to 0} \int_{\mathcal{M} \setminus B(y,\varepsilon)} \frac{\overline{\xi - y}}{|\xi - y|^3} n(\xi) h(\xi) d\sigma(\xi), \qquad y \in \mathcal{M}.$$

• This gives us an analog of the Plemelj formula:

$$C^+h(y)-C^-h(y)=h(y).$$

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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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#### Proposition

The distribution is silent from outside if and only if

$$2\pi\psi(y) = -\lim_{\varepsilon \to 0} \int_{\mathcal{M} \setminus B(y,\varepsilon)} \frac{\xi - y}{|\xi - y|^3} \cdot (\psi n + R^{-1}(D))(\xi) d\sigma(\xi) \quad y \in \mathcal{M}.$$

More complicated equation involving the curvature.

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More complicated equation involving the curvature. When the latter is constant one gets the previous characterization. Still silence from both sides means tangent and divergence-free.

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- In  $\mathbb{R}^n$ , 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, <math>v can be approximated arbitrarily close by rationals in  $L^p(S)$  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<sup>p</sup> (any field is approximable y a gradient).
  The Hardy-Hodge decomposition now implies:

#### $\mathsf{Theorem}$

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 .