Basics of harmonic polynomials and spherical functions

I present a beautiful, elegant, and rather simple theory of spherical harmonics, which I heard on one of Dynkin's seminars I was attending in my third year at the university.

By \mathbb{R}^d we denote the Euclidean space of points $x = (x^1, ..., x^d)$. When it makes sense, for real-valued u(x) on \mathbb{R}^d we denote

$$u_{x^{i}} = \frac{\partial u}{\partial x^{i}}, \quad u_{x^{i}x^{j}} = \frac{\partial}{\partial x^{j}}u_{x^{i}}, \quad \Delta u = u_{x^{1}x^{1}} + \ldots + u_{x^{d}x^{d}}.$$

1. Spherical harmonics and the Laplace-Beltrami operator

Denote by \mathcal{P}_n the set of polynomials of degree (at most) n defined in \mathbb{R}^d , $d \geq 2$. Denote

$$B_R = \{ x \in \mathbb{R}^d : |x| < R \}, \quad S_R = \{ x \in \mathbb{R}^d : |x| = R \}.$$

Lemma 1.1 (maximum principle). If $u \in C^2_{loc}(B_R) \cap C(\overline{B}_R)$ and $\Delta u \ge 0$ in B, then in B

$$u \leq \max_{S_{\mathcal{P}}} u.$$

In particular, if $\Delta u = 0$ in B and u = 0 on S_R , then u = 0 in B.

Proof. Let u at some points in B_R be strictly bigger than $\max_{S_R} u$. Then, by continuity, for $\varepsilon > 0$ small enough $u - \varepsilon (R^2 - |x|^2)$ will also be at some points in B_R strictly bigger than its maximum over S_R . It follows that the maximum point, say x_{ε} of $u - \varepsilon (R^2 - |x|^2)$ over \overline{B}_R lies in B_R . At x_{ε} the second order derivatives and, hence, the Laplacian is nonpositive:

$$0 \ge \Delta (u - \varepsilon (R^2 - |x|^2))|_{x = x_{\varepsilon}} = \Delta u(x_{\varepsilon}) + 2\varepsilon d.$$

But this is impossible, since $\Delta u(x_{\varepsilon}) \geq 0$, which proves the lemma. \Box

Lemma 1.2. Let f and g be polynomials. Then there exists a unique polynomial h, such that $\Delta h = f$ in B_1 and h = g on S_1 .

Proof. By considering h - g we reduce the problem to the one with $g \equiv 0$. Let $f \in \mathcal{P}_n$.

Observe that the operator $T: v \to Tu = \Delta[(1 - |x|^2)v]$ maps \mathcal{P}_n into \mathcal{P}_n . Furthermore, by the maximum principle, if $v \in \mathcal{P}_n$ and Tv = 0, then $(1 - |x|^2)v \equiv 0$, meaning that T is a one-to-one mapping. Since \mathcal{P}_n is a finite-dimensional linear space, the equation Tv = f has a unique solution in \mathcal{P}_n , which proves the lemma in light of the fact that $1 - |x|^2$ is a polynomial.

It is interesting to discuss the following without referring to the above lemma.

Remark 1.3. If $n \ge 2$, $u \in \mathcal{P}_n$ and u = 0 on S_1 , then $u = (1 - |x|^2)v$, where $v \in \mathcal{P}_{n-2}$.

Indeed, then $f := \Delta u \in \mathcal{P}_{n-2}$, there is a unique polynomial v such that Tv = f, and this $v \in \mathcal{P}_{n-2}$.

Remark 1.4. If $f \equiv 0$ and g is a polynomial of degree n, then h is also a polynomial of degree (at most) n.

Indeed, if n = 0, 1, then h = g, and, if $n \ge 2$, then $\Delta(g - h) = p$, where p is a polynomial of degree n - 2 and it follows from the above proof that $g - h = (1 - |x|^2)v$, where v is a polynomial of degree n - 2, so that h is a polynomial of degree (at most) n. In addition, if $n \ge 2$,

$$g = h + (1 - |x|^2)v = h + (1 - |x|^2)h_1 + (1 - |x|^2)^2v_1 = \dots$$

Example 1.5. If $g = |x|^2$, $h \equiv 1$.

Remark 1.6. Every harmonic polynomial $h \in \mathcal{P}_n$ can be uniquely represented as the sum $h_n + \ldots + h_0$, where h_k are homogeneous of degree k harmonic polynomials.

Definition 1.7. The set of harmonic polynomials is denoted by \mathcal{H} . By \mathcal{H}_n we denote the set of homogeneous polynomials of order n which are harmonic. Any element of \mathcal{H}_n restricted to S_1 is called a *spherical* harmonic of degree n. The set of those is denoted by $\mathcal{H}_n(S)$

Corollary 1.8. The set

$$\bigcup_{n=0}^{\infty} \{h_0 + \dots + h_n : h_i \in \mathcal{H}_i(S), i = 0, \dots, n\}$$

is dense in $L_2(S_1)$ since the set of polynomials is dense there.

Definition 1.9. The Laplace-Beltrami operator Δ_S on S_1 is introduced on smooth functions ϕ given on S_1 by the formula

$$\Delta_S \phi(x) = \Delta \left(\phi \left(\frac{x}{|x|} \right) \right) \Big|_{|x|=1}$$

Example 1.10. Let us compute $\Delta_S h$ for $h \in \mathcal{H}_n$. Observe that by Euler $nh(x) = x^i h_{x^i}(x)$ and on functions $\phi(x) = \phi(r)$, r = |x|, we have

$$\Delta \phi = \phi'' + \frac{d-1}{r}\phi'.$$

Also

$$\Delta(\psi\eta) = \psi\Delta\eta + \eta\Delta\psi + 2\psi_{x^i}\eta_{x^i}, \quad \left(\frac{1}{|x|^n}\right)_{x^i} = -\frac{nx^i}{|x|^{n+2}}$$

Then for |x| = 1 we have

$$\Delta\left(h\left(\frac{x}{|x|}\right)\right) = \Delta\left(\frac{1}{|x|^n}h(x)\right) = \left[n(n+1) - (d-1)n\right]h - 2nx^i h_{x^i},$$
$$\Delta_S h = -n(n+d-2)h,$$

so that h is an eigenfunction of Δ_S with eigenvalue

$$\lambda_n = -n(n+d-2).$$

Theorem 1.11. The operator Δ_S is formally self-adjoint. Moreover, for any smooth ϕ, ψ given on \mathbb{R}^d

$$I := \int_{S_1} \phi \Delta_S \psi \, dS = -\int_{S_1} (D^t \phi, D^t \psi) \, dS, \qquad (1.1)$$

where $D^t \phi(x)$, $x \in S_1$, is the projection of the gradient ϕ_x of ϕ at point x on the tangent plane to S_1 at x, that is

$$(D^t \phi(x))^i = \left[\phi\left(\frac{x}{|x|}\right)\right]_{x^i} = \phi_{x^j}(x)(\delta^{ij} - x^i x^j).$$

Proof. By using polar coordinates we write that for r > 1/2

$$\int_{B_r \setminus B_{1/2}} \phi\left(\frac{x}{|x|}\right) \Delta\left[\psi\left(\frac{x}{|x|}\right)\right] dx = \int_{1/2}^r \int_{S_\rho} \phi\left(\frac{x}{|x|}\right) \Delta\left[\psi\left(\frac{x}{|x|}\right)\right] dS_\rho \, d\rho.$$
It follows that for $r = 1$

It follows that for r = 1

$$I = \frac{d}{dr} \int_{B_r \setminus B_{1/2}} \phi\left(\frac{x}{|x|}\right) \Delta\left[\psi\left(\frac{x}{|x|}\right)\right] dx.$$

We use Green's formula and observe that the boundary terms disappear because x/|x| does not change along the normals to the boundary of $B_r \setminus B_{1/2}$. Hence, I equals

$$-\frac{d}{dr} \int_{B_r \setminus B_{1/2}} \left(\left(\phi\left(\frac{x}{|x|}\right) \right)_x, \left(\phi\left(\frac{x}{|x|}\right) \right)_x \right) dx$$
$$= -\int_{S_1} \left(\left(\phi\left(\frac{x}{|x|}\right) \right)_x, \left(\phi\left(\frac{x}{|x|}\right) \right)_x \right) dS,$$

which yields (1.1). The theorem is proved.

Corollary 1.12. If $h \in \mathcal{H}_n$, $g \in \mathcal{H}_m$ and $n \neq m$, then $h \perp g$, that is $\int_{S_1} hg \, dS = 0.$

Indeed,

$$\lambda_n \int_{S_1} hg \, dS = \int_{S_1} \Delta_S hg \, dS = \lambda_m \int_{S_1} hg \, dS,$$

which implies the result since $\lambda_n \neq \lambda_m$.

Theorem 1.13. Any function $g \in L_2(S_1)$ has a unique representation

$$g = \sum_{n=0}^{\infty} h_n, \tag{1.2}$$

where $h_n \in \mathcal{H}_n$ and the series converges in $L_2(S_1)$ -sense.

To prove this theorem it suffices to refer to Corollaries 1.12 and 1.8. $\hfill \Box$

Corollary 1.14. If ϕ is a smooth function on S_1 such that $\Delta_S \phi = \lambda \phi$, where λ is a constant, then either $\phi = \text{const}$ and $\lambda \phi = 0$, or $\phi \neq \text{const}$ and there exists an m = 1, 2, ... such that $\phi \in \mathcal{H}_m(S_1)$ and $\lambda = \lambda_m$.

Indeed, for $m \ge 1$ and the projection ϕ_m of ϕ on $\mathcal{H}_m(S_1)$ we have $\|\phi_m\|_{L_2(S_1)}^2 = (\phi, \phi_m)_{L_2(S_1)} = \lambda_m^{-1}(\phi, \Delta_S \phi_m)_{L_2(S_1)} = \lambda_m^{-1}(\Delta_S \phi, \phi_m)_{L_2(S_1)}$ $= \lambda \lambda_m^{-1} \|\phi_m\|_{L_2(S_1)}^2.$

Hence, $\phi_m = 0$ if $\lambda \neq \lambda_m$. Therefore, if $\lambda \notin \{\lambda_m : m \geq 1\}$, then $\phi \in \mathcal{H}_0(S_1), \phi = \text{const}, \Delta_S \phi = 0 = \lambda \phi$. However, if $\lambda = \lambda_{m_0}$ for some $m_0 \geq 1$, then $\phi \perp \mathcal{H}_m(S_1)$ for $m \neq m_0$, since, for $h \in \mathcal{H}_m$,

$$\lambda_{m_0} \int_{S_1} h\phi \, dS = \int_{S_1} \Delta_S h\phi \, dS = \lambda_m \int_{S_1} h\phi \, dS$$

Hence, $\phi \in \mathcal{H}_{m_0}(S_1)$.

For $\rho > 0$ denote by

$$\int_{S_{\rho}} u(x) \, dS_{\rho}$$

the average value of u on S_{ρ} , that is its integral over S_{ρ} divided by the surface of S_{ρ} . Similarly introduce

$$\int_{B_{\rho}} u(x) \, dx = \frac{1}{\operatorname{Vol}\left(B_{\rho}\right)} \int_{B_{\rho}} u(x) \, dx.$$

Theorem 1.15 (mean value theorem). If $h \in \mathcal{H}$, $a \in \mathbb{R}^d$, and $\rho > 0$, then

$$\int_{S_{\rho}} h(x+a) \, dS_{\rho} = \int_{B_{\rho}} h(x+a) \, dx = h(a). \tag{1.3}$$

Indeed, it suffices to consider the case where a = 0 and h(0) = 0. In that case, the equality

$$\int_{S_{\rho}} h(x) \, dS_{\rho} = 0 \tag{1.4}$$

for all $\rho > 0$ implies

$$\int_{S_{\rho}} h(x) \, dS_{\rho} = 0, \quad \int_{B_{\rho}} h(x) \, dx = 0, \quad \oint_{B_{\rho}} h(x) \, dx = 0.$$

Therefore, we only need to prove (1.4). Scalings imply that it suffices to concentrate on $\rho = 1$. For $h \in \mathcal{H}$ there exists $n \geq 0$ such that $h = h_n + \ldots + h_0$, where $h_i \in \mathcal{H}_i$. Since h(0) = 0 and, obviously, $h_i(0) = 0, i \geq 1$, it holds that $h_0 = 0$ and it suffices to prove that (1.4) holds for $\rho = 1, n \geq 1$, and $h \in \mathcal{H}_n$. In that case (1.4) follows from Corollary 1.12 since $1 \in \mathcal{H}_0$.

2. Dirichlet problem

Any *d*-tuple $\beta = (\beta_1, ..., \beta_d)$ consisting of $\beta_i \in \{0, 1, ...\}$ is called a multi-index. For any multi-index β we set

$$|\beta| = \sum_{i=1}^{d} \beta_i, \quad D^{\beta} = D_1^{\beta_1} \cdot \ldots \cdot D_d^{\beta_d}, \quad D_i = \frac{\partial}{\partial x^i} \quad \beta! = \beta_1! \cdot \ldots \cdot \beta_d!.$$

Lemma 2.1. Let $R \in (0, \infty)$, g, u be polynomials such that $\Delta u = 0$ and u = g on ∂B_R . Then there exists a constant N = N(d) such that in B_R for any multi-index β

$$\left|D^{\beta}u(x)\right| \le \left(\frac{N|\beta|}{R-|x|}\right)^{|\beta|} \sup_{\partial B_{R}} |g|.$$
(2.1)

Proof. We employ Bernstein's method which uses only the maximum principle. First we note that scalings show that it suffices to concentrate on R = 1. In that case take any $\zeta \in C_0^{\infty}(\mathbb{R}^d)$ with support in B_1 , assume that $\zeta(0) = 1$, and consider the function

$$w := \zeta^2 |u_x|^2 + \lambda |u|^2,$$

where $\lambda > 0$ is a constant to be chosen later. We have $\Delta u = 0$, $\Delta u_{x^i} = 0$, and

$$\Delta w = |u_x|^2 \Delta(\zeta^2) + \zeta^2 \left[2u_{x^i} \Delta u_{x^i} + 2\sum_{i,k} |u_{x^i x^k}|^2 \right] + 8\zeta \zeta_{x^i} u_{x^k} u_{x^i x^k} + 2\lambda |u_x|^2 + 2\lambda u \Delta u = \left[2\lambda |u_x|^2 + |u_x|^2 \Delta(\zeta^2) \right] + 2\zeta^2 \sum_{i,k} |u_{x^i x^k}|^2 + 8[\zeta_{x^i} u_{x^k}] [\zeta u_{x^i x^k}] \geq |u_x|^2 \left[2\lambda + \Delta(\zeta^2) - 8|\zeta_x|^2 \right]$$

(we have used $2a^2 + 8ab \ge -8b^2$). We see how to take $\lambda = \lambda(d)$ so that $\Delta w \ge 0$. Fix such a λ . Then by the maximum principle

$$\left|u_x(0)\right|^2 \le \max_{\bar{B}_1} w \le \max_{\partial B_1} w = \lambda \max_{\partial B_1} |g|^2$$

This yields (2.1) for $|\beta| = 1$, R = 1, and x = 0. Then (2.1) holds for $|\beta| = 1$ and any $R \in (0, \infty)$ if x = 0. Moving the origin and observing that $\Delta u = 0$ in $B_{R-|x_0|}(x_0)$ for any $x_0 \in B_R$ we get that (2.1) holds for $|\beta| = 1$ and any $x \in B_R$.

For higher values of $|\beta|$, one obtains (2.1) by splitting $B_R \setminus B_{|x|}$ into $|\beta|$ rings of width $(R - |x|)/|\beta|$ and estimating the derivatives $D_{j_1} \cdot \ldots \cdot D_{j_k} u$ inside the k-th ring by using the above result obtained for $|\beta| = 1$ and the fact that $\Delta D_{j_1} \cdot \ldots \cdot D_{j_{k-1}} u = 0$ in B_R . The lemma is proved. \Box

Corollary 2.2. For any $g \in C(\bar{B}_R)$ there exists a unique $u \in C^2_{loc}(B_R) \cap C(\bar{B}_R)$, that satisfies $\Delta u = 0$ in B_R and is equal to g on ∂B_R . Furthermore, any such function is in $C^{\infty}_{loc}(B_R)$ and the assertions of Lemma 2.1 hold true. Also, for any $a \in B_R$ and $\rho > 0$ such that $B_{\rho}(a) \subset B_R$ we have (mean value theorem)

$$\int_{S_{\rho}} u(x+a) \, dS_{\rho} = \int_{B_{\rho}} u(x+a) \, dx = u(a)$$

Uniqueness is a consequence of the maximum principle, which also assures that, if the polynomials g_n converge uniformly to g on ∂B_R and u_n are the corresponding polynomials from Lemma 2.1, then

$$\sup_{\bar{B}_R} |u_n - u_m| \le \sup_{\partial B_R} |g_n - g_m| \to 0$$

as $n, m \to \infty$. Hence, u_n converge in \overline{B}_R uniformly to, say u, which is continuous and equal g on ∂B_R . Applying estimates (2.1) to $u_n - u_m$, we see that also all derivatives of u_n converge uniformly on any compact subset of B_R . Of course, this implies that $u \in C^{\infty}_{\text{loc}}(B_R)$ and u is the desired solution. After that estimates (2.1) for u and the last assertions are immediate.

Corollary 2.3. If $u \in C^2_{loc}(\mathbb{R}^d)$, $\Delta u = 0$ in \mathbb{R}^d , $\alpha \in [0,1)$, $n \in \{0,1,2,...\}$ and

$$\lim_{|x|\to\infty}\frac{|u(x)|}{|x|^{n+\alpha}} < \infty,$$

then $u \in \mathcal{H}_n$.

Indeed, as $R \to \infty$, (2.1) implies that all derivatives of order n + 1 vanish.

Corollary 2.4. If $u \in C^2_{loc}(B_R)$ is harmonic in B_R , then it is real analytic.

This follows from (2.1) (applied in a smaller ball) and the fact that, for any constant N and $|x| \leq (eN)^{-1}$,

$$\frac{N^k |x|^k k^k}{k!} \to 0$$

as $k \to \infty$.

Corollary 2.5. Let $u_n(x)$, n = 1, 2, ..., be a sequence of harmonic $functions in <math>B_1$ such that they are uniformly bounded in B_R for any R < 1 and at any point of B_1 they converge as $n \to \infty$ to a function u(x). Then u(x) is infinitely differentiable, harmonic in B_1 , and any derivative of u_n converges to the corresponding derivative of u locally uniformly in B_1 .

Theorem 2.6 (Harnack's inequality). Let $u \in C^2_{loc}(B_R)$ be a nonnegative harmonic in B_R . Then $u(x) \leq 5^d u(y)$ if $|x|, |y| \leq R/5$.

Proof. We have $B_{R/2}(y) \supset B_{R/10}(x)$ and

$$u(y) = \int_{B_{R/2}(y)} u(z) \, dz \ge M \, \int_{B_{R/10}(x)} u(z) \, dz = M u(x),$$

where

$$M = \frac{\text{Vol}\,B_{R/10}}{\text{Vol}\,B_{R/2}} = 5^{-d}.$$

This proves the theorem.

Corollary 2.7 (one-sided Liouville's theorem). If $u \in C^2_{\text{loc}}(\mathbb{R}^d)$ is harmonic in \mathbb{R}^d which is bounded from below, then u = const.

Indeed, in that case $v := u - \inf_{\mathbb{R}^d} u$ is also a harmonic in \mathbb{R}^d and there is a sequence y_n such that $v(y_n) \to 0$ as $n \to \infty$. By Theorem 2.6 we have $v(x) \leq Nv(y_n)$, which implies that $(0 \leq)v \leq 0, v \equiv 0$, and $u \equiv \inf_{\mathbb{R}^d} u.$

A different argument. One can assume that u > 0. Then Harnack's inequality implies that u is bounded, and, by Corollary 2.3, $u \in \mathcal{H}_0$, that is u = const.

Next, we are interested in solving the Dirichlet problem with more general boundary data.

Theorem 2.8. For any function $g \in L_2(S_1)$ there exists a unique harmonic function u in B_1 such that $u(tx) \to g(x)$ in $L_2(S_1)$ -sense as $t \uparrow 1.$

Proof. Existence. We take the right-hand side of (1.2) and consider it in B_1 . Obviously, $h_n \perp h_m$ in $L_2(B_1)$ if $n \neq m$. Furthermore,

$$\int_{B_1} |h_n|^2 dx = c \int_0^1 r^{d-1} \left(\int_{S_1} |h_n(rx)|^2 dS \right) dr$$
$$= c \int_0^1 r^{d+2n-1} dr \int_{S_1} |h_n(x)|^2 dS = c(d+2n)^{-1} ||h_n||_{L_2(S_1)}^2$$

It follows that the series in (1.2) converges in $L_2(B_1)$.

Finite sums Σ_n of this series are well defined harmonic functions satisfying

$$\Sigma_n(x) = \int_{B_r(x)} \Sigma_n(y) \, dy.$$

By Corollary 2.5, Σ_n converge uniformly in any B_r , r < 1, to a harmonic function. Call it u. Then, for $x \in S_1$ and $t \in [0, 1)$ we have

$$u(tx) = \sum_{n=0}^{\infty} t^n h_n(x).$$

Hence

$$||u(t\cdot) - g||_{L_2(S_1)}^2 = \sum_{n=0}^{\infty} (1 - t^n) ||h_n||_{L_2(S_1)}^2,$$

which goes to zero as $t \uparrow 1$ indeed.

Uniqueness. Observe that by Green's formula for t < 1

$$\frac{d}{dt} \int_{S_1} |u(tx)|^2 \, dS = 2 \int_{S_1} u(tx)(x, u_x(tx)) \, dS = 2t \int_{B_1} |u_x(tx)|^2 \, dx \ge 0$$

and this yields the result. The theorem is proved. \Box

and this yields the result. The theorem is proved.

3. Case d = 2

Let s denote the point on S_1 which is at the distance s along the circle in the counterclockwise direction from the point (1,0). As easily follows from (1.1), the Laplace-Beltrami operator on S_1 is just the second-order derivative with respect to s. Therefore, $h \in \mathcal{H}_n(S_1)$ are (the only, see Corollary 1.14) solutions of

$$h'' = \lambda_n h = -n^2 h, \quad h(s) = c_1 \sin ns + c_2 \cos ns,$$

and (1.2) is a usual representation of a function $f \in L_2(0, 2\pi)$ by its Fourier series. Amazing, at the first sight, the theory of Fourier series has nothing to do with the Laplacian in d = 2.

Also observe that since, generally, \mathcal{H}_n are obviously invariant under orthogonal transformation in this case for d = 2 and n = 1 for any twe have

$$\sin(t+s) = b(t)\cos s + c(t)\sin s$$

and one can find b(t) and c(t) by using substitutions. Say, for s = 0 we get $\sin t = b(t)$. Knowing that and interchanging s and t we get $c(t) = \cos t$.

One more remarkable thing is that $\sin ns$, which is in $\mathcal{H}_n(S)$ by Corollary 1.14, as a solution of the appropriate equation, and is therefore the trace on S_1 of a homogeneous *n*th-order harmonic *polynomial*. For n = 2 this means that $\sin 2s$ and $\cos 2s$ are linear combinations of the traces of $x^2 - y^2$ and xy, that is of $\cos^2 s - \sin^2 s$ and $\cos s \sin s$.

4. Basis in \mathcal{H}_n and its dimension

Let $\mathcal{P}_n^{\text{hom}}$ be the set of *homogeneous* polynomials of degree *n*.

Lemma 4.1. For $n \geq 2$ every $p \in \mathcal{P}_n^{\text{hom}}$ has a unique representation

$$p = h + |x|^2 r,$$

where $h \in \mathcal{H}_n$ and $r \in \mathcal{P}_{n-2}^{\text{hom}}$.

Proof. Existence. By Lemma 1.2 there is a harmonic polynomial u of degree n such that p = u on S_1 . By Remark 1.3 we have $p - u = (1 - |x|^2)q$, where $q \in \mathcal{P}_{n-2}$. Let h be the homogeneous part of u of order n and r be the homogeneous part of q of order n - 2. Then

$$p = h - |x|^{2}r + [(u - h) + r + (1 - |x|^{2})(q - r)]$$
$$= h - |x|^{2}r + p_{n-1},$$

where $p_{n-1} \in \mathcal{P}_{n-1}$. By homogeneity $p_{n-1} = 0$ and we get the desired representation.

Uniqueness. We need to prove that if $h \in \mathcal{H}_n$, $r \in \mathcal{P}_{n-2}^{\text{hom}}$ and $h + |x|^2 r \equiv 0$, then $h = r \equiv 0$.

More generally, let $k \ge 1$ be an integer such that h admits the representation $h = -|x|^{2k}r$, where r is a polynomial. We want to show that h = 0. Obviously $n - 2k \ge 0$ and

$$0 = \Delta(|x|^{2k}r) = |x|^{2k}\Delta r + 2k(2n - 2k + d - 2)|x|^{2k-2}r,$$

Here 2k(2n-2k+d-2) > 0, so that $r_1 := c\Delta r \ (c^{-1} = -2k(2n-2k+d-2))$ is a polynomial, and $h = -|x|^{2(k+1)}r_1$. By induction we see that for any $k \ge 1$ there exists polynomials r_k such that $h = |x|^{2k}r_k$, which is only possible if $h \equiv 0$ and this proves the lemma. \Box

As a simple corollary of Lemma 4.1 obtained by iteration we come to the following.

Theorem 4.2. Let $p \in \mathcal{P}_n$. Then there exist unique $h_{n-2i} \in \mathcal{H}_{n-2i}$, i = 0, 1, ..., k, where $k = \lfloor n/2 \rfloor$ such that

$$p = h_n + |x|^2 h_{n-2} + \dots + |x|^{2k} h_{n-2k}.$$

Example 4.3. Take $p = x^1 x^2 x^3$. It has a representation $p = h + |x|^2 h_1$, where $h \in \mathcal{H}_3$ and h_1 is an affine homogeneous function of $x = (x^1, x^2, x^3)$. By symmetry it follows that $h_1 = c(x^1 + x^2 + x^3)$ and $x^1 x^2 x^3 - c|x|^2 (x^1 + x^2 + x^3)$ is harmonic.

Theorem 4.4. If $n \ge 2$, then

$$\dim \mathcal{H}_n = \begin{pmatrix} d+n-1\\ d-1 \end{pmatrix} - \begin{pmatrix} d+n-3\\ d-1 \end{pmatrix}.$$

Proof. First we find dim $\mathcal{P}_n^{\text{hom}}$, which is the number of different monomials

$$x^{\alpha} = (x^1)^{\alpha_1} \cdot \ldots \cdot (x^d)^{\alpha_d}$$

such that $\alpha_1 + \ldots + \alpha_d = n$. We consider a row of d+n-1 seats, numbered from 1 to d+n-1, choose arbitrarily d-1 of them $i_1 < i_2 < \ldots < i_{d-1}$ and then define α_k as the number of seats strictly between i_k and i_{k+1} : $\alpha_1 = i_1 - 1$, $\alpha_2 = i_2 - i_1 - 1$, ..., $\alpha_d = d + n - 1 - i_{d-1}$. The number of such arrangement will be exactly the number of different monomials x^{α} such that $\alpha_1 + \ldots + \alpha_d = n$. This number is

$$\binom{d+n-1}{d-1}.$$

By Lemma 4.1 to obtain dim \mathcal{H}_n it suffices to subtract from this number dim $\mathcal{P}_{n-2}^{\text{hom}}$. This yields the result.

Remark 4.5. If d = 2 and $n \ge 1$, dim $\mathcal{H}_n = 2$.

Now we want to find a basis in \mathcal{H}_n .

Lemma 4.6. For d > 2 and any multi-index α the function

$$h := |x|^{2n+d-2} D^{\alpha} \frac{1}{|x|^{d-2}}$$

belongs to \mathcal{H}_n , where $n = |\alpha|$.

Proof. That $h \in \mathcal{P}_n^{\text{hom}}$ is checked by induction. Then, since $D^{\alpha} \frac{1}{|x|^{d-2}}$ is harmonic in $\mathbb{R}^d \setminus \{0\}$ and homogeneous of degree -(n+d-2), we have

$$\begin{aligned} |x|^{-n-d+4}\Delta h &= [(2n+d-2)(2n+d-3) + (d-1)(2n+d-2)]D^{\alpha}\frac{1}{|x|^{d-2}} \\ &+ 2(2n+d-2)x^{j}D_{j}D^{\alpha}\frac{1}{|x|^{d-2}} \\ &= [(2n+d-2)(2n+d-3) + (d-1)(2n+d-2) \\ &- 2(2n+d-2)(n+d-2)]D^{\alpha}\frac{1}{|x|^{d-2}} = 0. \end{aligned}$$

Lemma 4.7. Let $d \geq 3$ and let $u \in C^2(\mathbb{R}^d)$ have compact support. Then for any $x \in \mathbb{R}^d$

$$u(x) = -c_d \int_{\mathbb{R}^d} \frac{1}{|y - x|^{d-2}} \Delta u(y) \, dy, \tag{4.1}$$

(Newton's potential of $-\Delta u$) where c_d is a constant depending only on d.

Proof. First note that it suffices to prove (4.1) for x = 0. In that case take $\zeta \in C_0^{\infty}(\mathbb{R}^2)$ such that $\zeta = 1$ in B_1 and for $\varepsilon > 0$ set $\zeta_{\varepsilon}(y) = \zeta(y/\varepsilon)$. Observe that if in (4.1) we take $(1 - \zeta_{\varepsilon})u$ in place of u, then we can integrate by parts (x = 0) and, using the fact that $\Delta |y|^{2-d} = 0$, conclude that the integral in (4.1) is zero. Hence, the integral in (4.1) in its original form equals

$$I := \lim_{\varepsilon \downarrow 0} \int_{\mathbb{R}^d} \frac{1}{|y|^{d-2}} \Delta(\zeta_{\varepsilon} u)(y) \, dy.$$

We use that

$$\varepsilon^{-1} \int_{\mathbb{R}^d} \frac{1}{|y|^{d-2}} \zeta_{x^i}(y/\varepsilon) u_{x^i}(y) \, dy = \varepsilon \int_{\mathbb{R}^d} \frac{1}{|y|^{d-2}} \zeta_{x^i}(y) u_{x^i}(\varepsilon y) \, dy \to 0,$$
$$\int_{\mathbb{R}^d} \frac{1}{|y|^{d-2}} \zeta(y/\varepsilon) \Delta u(y) \, dy \to 0,$$

and conclude that

$$I := \lim_{\varepsilon \downarrow 0} \varepsilon^{-2} \int_{\mathbb{R}^d} \frac{1}{|y|^{d-2}} u(y) (\Delta \zeta) (y/\varepsilon) \, dy.$$

Observe that for N being the Lipschitz constant of u

$$\lim_{\varepsilon \downarrow 0} \varepsilon^{-2} \int_{\mathbb{R}^d} \frac{1}{|y|^{d-2}} |u(y) - u(0)| |\Delta \zeta| (y/\varepsilon) \, dy$$
$$\leq N \lim_{\varepsilon \downarrow 0} \varepsilon \int_{\mathbb{R}^d} \frac{1}{|y|^{d-2}} |y| |\Delta \zeta| (y) \, dy = 0.$$

Hence,

$$I = u(0) \lim_{\varepsilon \downarrow 0} \varepsilon^{-2} \int_{\mathbb{R}^d} \frac{1}{|y|^{d-2}} (\Delta \zeta) (y/\varepsilon) \, dy = u(0) \int_{\mathbb{R}^d} \frac{1}{|y|^{d-2}} \Delta \zeta(y) \, dy.$$

This proves the lemma.

Recall that we denote by $\mathcal{H}_n(S_1)$ the set of restrictions of the functions in \mathcal{H}_n to S_1 .

Lemma 4.8. If d > 2, the linear span of

$$\{D^{\alpha}|x|^{2-d}: |\alpha|=n\}$$

is $\mathcal{H}_n(S_1)$.

Proof. If n = 0, the assertion is obvious. For $n \ge 1$ assume the contrary. Then there exists $h \in \mathcal{H}_n(S_1)$ such that

$$\int_{S_1} h D^{\alpha} |x|^{2-d} \, dS = 0 \tag{4.2}$$

for all α such that $|\alpha| = n$. By Lemma 4.6 and Corollary 1.12 equation (4.2) also holds for $|\alpha| \neq n$.

Next, the function

$$F(x) := \int_{S_1} h(y) |y - x|^{2-d} \, dS$$

is harmonic outside S_1 and as such is real analytic there. Equation (4.2) says that all derivatives of F vanish at 0. It follows that F = 0 in B_1 . Furthermore, F is a bounded continuous function on \mathbb{R}^d since bounded continuous

$$\int_{S_1} h(y) |y - x|^{2-d} I_{|y - x| \ge \varepsilon} \, dS$$

converge uniformly on \mathbb{R}^d to F as $\varepsilon \downarrow 0$. In addition, $F \to 0$ as $|x| \to \infty$, hence, by the maximum principle F = 0 in B_1^c , since it is harmonic in \overline{B}_1^c . Thus, $F \equiv 0$.

Now take $u \in C_0^{\infty}(\mathbb{R}^d)$ and integrate the equality

$$0 = \int_{S_1} h(y) |y - x|^{2-d} \, dS \Delta u(x)$$

12

over \mathbb{R}^d . By using Lemma 4.7 we get

$$\int_{S_1} h(y)u(y)\,dS = 0$$

and the arbitrariness of u implies that $h \equiv 0$.

Theorem 4.9. For $d \ge 3$ the set

$$\{D^{\alpha}|x|^{2-d}: |\alpha| = n, \alpha_1 \le 1\}$$
(4.3)

is a basis in $\mathcal{H}_n(S_1)$.

Proof. The number of elements in (4.3) is not greater than the number of multi-indices α such that $|\alpha| = n$ and $\alpha_1 \leq 1$. The latter number is easily shown to be equal to dim $\mathcal{H}_n(S_1)$ (use the same interpretation as in the proof of Theorem 4.4). Therefore, due to Lemma 4.8 we only need to show that $D^{\alpha}|x|^{2-d}$ belongs to the span of (4.3) for any α_1 if $|\alpha| = n$. This is trivial because $|x|^{2-d}$ is harmonic and for any even α_1 we have

$$D_1^{\alpha_1} |x|^{2-d} = (-D_2^2 - \dots - D_d^2)^{\alpha_1/2} |x|^{2-d},$$

whereas if $\alpha_1 = 2k + 1 \ (\leq n)$,

$$D_1^{\alpha_1} |x|^{2-d} = D_1 (-D_2^2 - \dots - D_d^2)^k |x|^{2-d}.$$

5. An unexpected formula for the scalar product of spherical harmonics

Recall that for any multi-index $\alpha = (\alpha_1, ..., \alpha_d)$ and $x \in \mathbb{R}^d$ we set $x^{\alpha} = (x^1)^{\alpha_1} \cdot ... \cdot (x^d)^{\alpha_d}, \quad |\alpha| = \alpha_1 + ... + \alpha_d, \quad \alpha! = \alpha_1! \cdot ... \cdot \alpha_d!.$ Suppose

$$p = \sum_{|\alpha|=n} b_{\alpha} x^{\alpha}, \quad q = \sum_{|\alpha|=n} c_{\alpha} x^{\alpha}$$

are in \mathcal{H}_n . We want to find their scalar product in $L_2(S_1)$. It looks like the best we can do is to write

$$(p,q)_{L_2(S_1)} = \sum_{\alpha,\beta} b_\alpha c_\beta \int_{S_1} x^{\alpha+\beta} \, dS.$$

The integral over S_1 of the monomial $x^{\alpha+\beta}$ was explicitly calculated by Hermann Weyl in Section 3 of [2] (1939). Using that result would complete the formula for $(p, q)_{L_2(S_1)}$. There is however a shorter way, which I take from [1].

Lemma 5.1. If n > 0 and $p, q \in \mathcal{H}_n$, then

$$I := \int_{S_1} pq \, dS = \frac{1}{n(d+2n-2)} \sum_{j=1}^d \int_{S_1} D_j p D_j q \, dS. \tag{5.1}$$

Proof. Let Du denote the gradient of u. By Theorem 1.11 and Example 1.10

$$-\lambda_n I = \int_{S_1} \left(D\left(\frac{1}{|x|^n} p(x)\right), D\left(\frac{1}{|x|^n} q(x)\right) \right) dS$$
$$= \int_{S_1} \left(-nxp(x) + Dp(x), -nxq(x) + Dq(x) \right) dS.$$

By using that (x, Dp(x)) = np(x) we see that

$$-\lambda_n I = -n^2 I + \int_{S_1} (Dp, Dq) \, dS,$$

which is equivalent to (5.1).

Theorem 5.2. Let $p = \sum_{\alpha} b_{\alpha} x^{\alpha}$ and $q = \sum_{\alpha} c_{\alpha} x^{\alpha}$ be harmonic polynomials. Then

$$(p,q)_{L_2(S_1)} = \sum_{\alpha} b_{\alpha} c_{\alpha} w_{\alpha},$$

where

$$w_{\alpha} = \frac{\alpha!}{d(d+2) \cdot \ldots \cdot (d+2|\alpha|-2)} \quad \alpha \neq 0, \quad w_0 = 1.$$

Proof. Owing to Corollary 1.12 it suffices to prove the theorem under the assumption that $p, q \in \mathcal{H}_n$.

If n = 0, then p, q are constant and the desired result obviously holds.

If n > 0 observe that $D_j p$ and $D_j q$ are harmonic polynomials in \mathcal{H}_{n-1} , so that Lemma 5.1 is applicable. By induction we get for e_j being the *j*th basis vector in \mathbb{R}^d that

$$\sum_{j=1}^{d} (D_j p, D_j q)_{L_2(S_1)} = \sum_{j=1}^{d} \int_{S_1} \left(\sum_{\alpha} b_{\alpha} \alpha_j x^{\alpha - e_j} \right) \left(\sum_{\alpha} c_{\alpha} \alpha_j x^{\alpha - e_j} \right)$$
$$= \sum_{j=1}^{d} \sum_{\alpha} b_{\alpha} c_{\alpha} \alpha_j^2 \frac{(\alpha - e_j)!}{d(d+2) \cdot \dots \cdot (2+2n-4)}$$
$$= \sum_{\alpha} b_{\alpha} c_{\alpha} \sum_{j=1}^{d} \alpha_j \frac{\alpha!}{d(d+2) \cdot \dots \cdot (2+2n-4)}$$

$$=\sum_{\alpha}b_{\alpha}c_{\alpha}\frac{\alpha!n}{d(d+2)\cdot\ldots\cdot(2+2n-4)}.$$

By combining this with Lemma 5.1, we get the result.

References

- [1] Sheldon Axler, Paul Bourdon, and Wade Ramey, "Harmonic function theory", Springer, 2001.
- [2] Hermann Weyl, On the volume of tubes, American Journal of Mathematics, 61 (1939), 461-472.