

## Anisotropic Franklin bases on polygonal domains

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Franklin systems induced by Courant elements over multilevel nested triangulations of polygonal domains in  $\mathbb{R}^2$  are explored. Mild conditions are imposed on the triangulations which prevent them from deterioration and at the same time allow for a lot of flexibility and, in particular, arbitrarily sharp angles. It is shown that such anisotropic Franklin systems are Schauder bases for  $C$  and  $L_1$ , and unconditional bases for  $L_p$  ( $1 < p < \infty$ ) and the corresponding Hardy spaces  $H_1$ . It is also proved that the anisotropic  $H_1$  is exactly the space of all functions in  $L_1$  for which the corresponding Franklin system expansions converge unconditionally in  $L_1$ . Finally, it is shown that the Franklin bases characterize the corresponding anisotropic BMO spaces.

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### 1 Introduction

The Franklin systems in the univariate case as well as in the multivariate case in regular setups are thoroughly studied and well-known. We refer the reader to [11], [1], and [8] as references for Franklin systems.

In this article, we consider Franklin systems generated by sequences of Courant elements, i.e., piecewise linear elements, induced by multilevel nested triangulations of compact polygonal domains in  $\mathbb{R}^2$ . For a given polygonal domain  $E$  in  $\mathbb{R}^2$  we consider a sequence of nested triangulations  $\mathcal{T}_0, \mathcal{T}_1, \dots$  of a general nature. Mild conditions are imposed on the triangulations which prevent them from deterioration. At the same time these conditions allow for a great deal of flexibility and, in particular, arbitrarily sharp angles.

We show that the Franklin systems obtained by applying the Gram–Schmidt orthogonalization process to the corresponding Courant elements are Schauder bases for  $C$  and  $L_1$ , and unconditional bases for  $L_p$  ( $1 < p < \infty$ ) and the corresponding Hardy space  $H_1$ . Further, we prove that  $H_1$  is exactly the space of all functions in  $L_1$  for which the corresponding Franklin system expansions converge unconditionally in  $L_1$ . Finally, we show that the anisotropic Franklin systems characterize the corresponding BMO spaces. Thus we show that the basic and well-known results on Franklin bases in the regular case have analogues in the anisotropic case. We do not consider anisotropic atomic Hardy spaces  $H_p$  with  $0 < p < 1$  in this article.

The motivation for this article is two-fold. On the one hand the spaces induced by general multilevel nested triangulations are an example of spaces on homogeneous type. There are no bases available for such spaces in general and hence such bases are worthy to be studied. On the other hand the Franklin bases that we explore are the only anisotropic bases over general sequences of nested triangulations. There are no constructions of spline wavelet or prewavelet bases over such triangulations available as for now.

In [12] we show that the anisotropic Franklin systems considered in this article characterize the anisotropic  $B$ -space (generalized Besov spaces) which are naturally associated with hierarchical sequences of nested triangulations. These spaces play a fundamental role in nonlinear spline approximation (see [4, 5, 9, 10]).

The paper is organized as follows. In §2 we give all auxiliary results needed for the development of the anisotropic Franklin bases. In §3 we introduce the anisotropic Franklin systems and state and proof our main

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results. In the Appendix we give an example which shows that the anisotropic  $H_1$  spaces essentially depend on the triangulations which are used.

**Notation** Throughout this article for a set  $G \subset \mathbb{R}^2$ ,  $|G|$  denotes the Lebesgue measure of  $G$ , while  $G^\circ$  means the interior of  $G$ ;  $\mathbb{1}_G$  denotes the characteristic function of  $G$ , and  $\tilde{\mathbb{1}}_G := |G|^{-1/2} \mathbb{1}_G$ . For a finite set  $G$ ,  $\#G$  denotes the cardinality of  $G$ . Positive constants are denoted by  $c, c_1, \dots$  (if not specified, they may vary at every occurrence),  $A \approx B$  means  $c_1 A \leq B \leq c_2 B$ , and  $A := B$  or  $B =: A$  stands for “ $A$  is by definition equal to  $B$ ”. We set  $\langle f, g \rangle := \int_E fg$ .

## 2 Preliminaries

In this section we collect all prerequisites regarding triangulations, maximal operators, Hardy spaces on spaces of homogeneous type, and other results, which will be needed in the development of the Franklin bases. Most of these facts are well-known and we give only the essentials and suitable references for them.

### 2.1 Multilevel triangulations

We call  $E \subset \mathbb{R}^2$  a *bounded polygonal domain* if its interior  $E^\circ$  is connected and  $E$  is the union of a finite set  $\mathcal{T}_0$  of closed triangles with disjoint interiors:  $E = \bigcup_{\Delta \in \mathcal{T}_0} \Delta$ . Following [9] we call

$$\mathcal{T} = \bigcup_{m=0}^{\infty} \mathcal{T}_m$$

a *locally regular triangulation* of  $E$  or briefly an *LR-triangulation* with levels  $(\mathcal{T}_m)_{m \geq 0}$  if the following conditions are fulfilled:

- Every level  $\mathcal{T}_m$  is a partition of  $E$ , that is,  $E = \bigcup_{\Delta \in \mathcal{T}_m} \Delta$  and  $\mathcal{T}_m$  consists of closed triangles with disjoint interiors.
- The levels  $(\mathcal{T}_m)$  of  $\mathcal{T}$  are nested, i.e.,  $\mathcal{T}_{m+1}$  is a refinement of  $\mathcal{T}_m$ .
- Each triangle  $\Delta \in \mathcal{T}_m$  has at least two and at most  $M_0$  children (subtriangles) in  $\mathcal{T}_{m+1}$ , where  $M_0 \geq 2$  is a constant.
- The valence  $N_v$  of each vertex  $v$  of any triangle  $\Delta \in \mathcal{T}_m$  (the number of the triangles from  $\mathcal{T}_m$  which share  $v$  as a vertex) is at most  $N_0$ , where  $N_0$  is a constant.
- No hanging vertices condition*: No vertex of any triangle  $\Delta \in \mathcal{T}_m$  which belongs to the interior of  $E$  lies in the interior of an edge of another triangle from  $\mathcal{T}_m$ .
- There exist constants  $0 < r < \rho < 1$  ( $r \leq \frac{1}{2}$ ) such that for each  $\Delta \in \mathcal{T}_m$  ( $m \geq 0$ ) and any child  $\Delta' \in \mathcal{T}_{m+1}$  of  $\Delta$ ,

$$r|\Delta| \leq |\Delta'| \leq \rho|\Delta|. \quad (2.1)$$

- There exists a constant  $0 < \delta \leq 1$  such that for  $\Delta', \Delta'' \in \mathcal{T}_m$  ( $m \geq 0$ ) with a common vertex,

$$\delta \leq |\Delta'|/|\Delta| \leq \delta^{-1}. \quad (2.2)$$

The notion of a *regular triangulation* will be needed later on. We call  $\mathcal{T} = \bigcup_{m=0}^{\infty} \mathcal{T}_m$  a regular triangulation of a bounded polygonal domain  $E \subset \mathbb{R}^2$  if  $\mathcal{T}$  satisfies conditions (a)–(e) of LR-triangulations and also the *minimal angle condition*, that is,  $\min \text{angle}(\Delta) \geq \beta$  for every triangle  $\Delta \in \mathcal{T}$ , where  $\beta > 0$  is a constant. Evidently, every regular triangulation is locally regular but not the other way around. For other types of triangulations, see [9].

We denote by  $\mathcal{V}_m$  the set of all vertices of triangles from  $\mathcal{T}_m$ , where if  $v$  is on the boundary of  $E$  we include in  $\mathcal{V}_m$  as many different copies of  $v$  as is its multiplicity. The multiplicity of a vertex  $v \in \mathcal{T}_0$  on the boundary of  $E$  can be larger than one if the interior of the union of all triangles which share  $v$  as a vertex is not connected. Also, cuts in  $E$  along edges of triangles from  $\mathcal{T}_0$  are possible; such edges belong to the boundary of  $E$ .

We let  $\mathcal{E}_m$  denote the set of all edges of triangles in  $\mathcal{T}_m$ . We also set  $\mathcal{V} := \bigcup_{m \geq 0} \mathcal{V}_m$  and  $\mathcal{E} := \bigcup_{m \geq 0} \mathcal{E}_m$ .

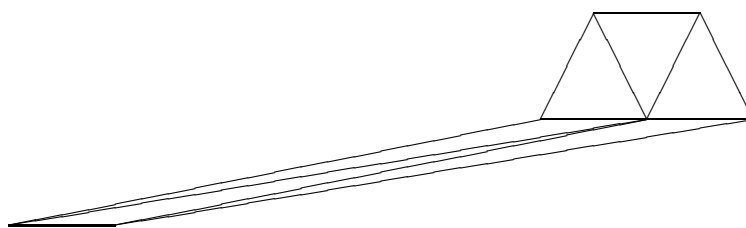


Fig. 1 A skewed cell in an LR-triangulation

We next clarify a number of issues concerning LR-triangulations, which are discussed in detail in [9] (see also [5, 10]).

The constants  $M_0, N_0, r, \rho, \delta$ , and  $\#\mathcal{T}_0$  (the cardinality of  $\mathcal{T}_0$ ) associated with an LR-triangulation  $\mathcal{T}$  are assumed fixed. We refer to them as *parameters* of  $\mathcal{T}$ .

It is an important observation that the collection of all LR-triangulations with given (fixed) parameters is invariant under affine transforms. More precisely, if  $\mathcal{T}$  is an LR-triangulation of  $E \subset \mathbb{R}^2$  and  $\mathbf{A}$  is an affine transform of  $\mathbb{R}^2$ , then  $\mathbf{A}(\mathcal{T}) := \{\mathbf{A}(\Delta) : \Delta \in \mathcal{T}\}$  is an LR-triangulation of the polygonal domain  $\mathbf{A}(E)$  with the same parameters.

The most important conditions (f)–(g) on LR-triangulations involve only areas of triangles but not angles. Consequently, if  $\mathcal{T}$  is an LR-triangulation and  $\Delta', \Delta'' \in \mathcal{T}_m$  have a common edge, then it may happen that  $\Delta'$  is an equilateral triangle (or close to an equilateral triangle) but  $\Delta''$  has a uncontrollably sharp angle (see Fig. 1).

In an LR-triangulation  $\mathcal{T}$  there can be an equilateral (or close to such) triangle  $\Delta^\diamond$  at any level  $\mathcal{T}_m$  with descendants  $\Delta_1 \supset \Delta_2 \supset \dots$  such that  $\min \text{angle}(\Delta_j) \rightarrow 0$  as  $j \rightarrow \infty$ .

It is important to know how fast the area  $|\Delta|$  of a triangle  $\Delta \in \mathcal{T}_m$  may change when  $\Delta$  moves away from a fixed triangle within the same level. Condition (f) suggests a geometric rate of change but in fact it is polynomial.

**Lemma 2.1** *If  $\Delta, \Delta' \in \mathcal{T}_m$  can be connected by  $n$  intermediate edges from  $\mathcal{E}_m$ , then*

$$c_1^{-1}(n+1)^{-s} \leq |\Delta'|/|\Delta''| \leq c_1(n+1)^s, \tag{2.3}$$

where  $s, c_1 > 0$  depend only on the parameters of  $\mathcal{T}$ .

This result follows easily by the following lemma (see [9, Lemma 2.4]):

**Lemma 2.2** *Let  $\mathcal{T}$  be an LR-triangulation of  $E \subset \mathbb{R}^2$ . Suppose that  $\Delta', \Delta'' \in \mathcal{T}_m$  (with  $m$  sufficiently large) and  $\Delta'$  and  $\Delta''$  can be connected by less than  $2^\nu$  intermediate edges from  $\mathcal{E}_m$  with (pairwise) common vertices. Then there exist  $\Delta_1, \Delta_2 \in \mathcal{T}_{m-2N_0\nu}$  with a common vertex such that  $\Delta' \subset \Delta_1$  and  $\Delta'' \subset \Delta_2$ .*

From the above discussion (see Fig. 1) it follows that for two triangles  $\Delta', \Delta'' \in \mathcal{T}_m$  which share a vertex  $|\max \text{edge}(\Delta')|/|\max \text{edge}(\Delta'')|$  can be uncontrollably large (or small). However, when going in depth the maximal edges of the triangles behave similarly as their areas.

**Lemma 2.3** *If  $\mathcal{T}$  is an LR-triangulation of  $E$ , there exist constants  $0 < r_1 < \rho_1 < 1$  depending only on the parameters of  $\mathcal{T}$  such that if  $\Delta' \subset \Delta, \Delta \in \mathcal{T}_m$  ( $m \geq 0$ ), and  $\Delta' \in \mathcal{T}_{m+3N_0\nu}$ ,  $\nu \geq 1$ , then*

$$r_1^\nu \leq \frac{|\max \text{edge}(\Delta')|}{|\max \text{edge}(\Delta)|} \leq \rho_1^\nu. \tag{2.4}$$

*Proof.* Evidently, it suffices to prove (2.4) for  $\nu = 1$  only. Let  $\Delta \in \mathcal{T}_m$  and let  $e$  be an edge of  $\Delta$ . If it is also an edge of a child of  $\Delta$ , then the valence of at least one of the two endpoints of  $e$  will increase by one at level  $m + 1$ . (Recall that there are always at least two children, so that a child and a parent cannot be the same triangle.) Therefore,  $e$  will be subdivided at least once after at most  $S := 2(N_0 - 3) + 1$  steps of refinement. By (2.1) it readily follows that any edge  $e'$  obtained by subdividing  $e$  satisfies  $|e'| \leq \rho|e|$ .

We say that an edge of a descendant of  $\Delta$  is a *cutting edge* for  $\Delta$  if one of its endpoints is a vertex of  $\Delta$  and the other lies in the interior of the opposite edge of  $\Delta$ . Since all cutting edges must emanate from the same vertex of  $\Delta$ , there are totally no more than  $M := N_0 - 3$  such edges for  $\Delta$ . Therefore, no new cutting edges for  $\Delta$  will be created at levels  $l > m + N_0 - 3$ . (It is easy to see that, as long as no new cutting edges are created

at a level  $l$ , they cannot be created at any further level.) Using this and the above observation, we conclude that there will be no cutting edges at levels  $l > m + M + S$  since they all will be subdivided. Therefore, each edge  $e'$  inside  $\Delta$  at these levels is either a proper part of an edge of  $\Delta$ , or has both of its endpoints in the interiors of two different edges of  $\Delta$ , or it has at least one endpoint in the interior of  $\Delta$ . In all cases, condition (2.1) ensures that  $|e'| \leq \rho |\max \text{edge}(\Delta)|$ . Consequently, if  $\Delta' \in \mathcal{T}_{m+3N_0}$  then  $|\max \text{edge}(\Delta')| \leq \rho |\max \text{edge}(\Delta)|$ , since  $3N_0 > M + S + 1$ . Thus the upper bound in (2.4) is established.

The argument for the proof of the lower bound in (2.4) is simpler. Suppose  $\Delta \in \mathcal{T}_m$ ,  $\Delta' \in \mathcal{T}_{m+1}$ , and  $\Delta' \subset \Delta$ . Let  $e_{\max}$  and  $e'_{\max}$  be the largest edges of  $\Delta$  and  $\Delta'$ , respectively. Denote by  $h$  the length of the height to  $e_{\max}$  in  $\Delta$  and by  $h'$  the length of the height to  $e'_{\max}$  in  $\Delta'$ . Further, let  $R$  and  $R'$  be the radii of the circles inscribed in  $\Delta$  and  $\Delta'$  respectively. A simple geometric argument shows that  $R < h < 3R$  as well as  $R' < h' < 3R'$ . Since  $\Delta' \subset \Delta$ , then  $R' \leq R$  and hence  $h' < 3h$ . We use this and (2.1) to obtain

$$(1/2)r |e_{\max}| h = r |\Delta| \leq |\Delta'| \leq (1/2) |e'_{\max}| h' \leq (3/2) |e'_{\max}| h$$

which implies  $|e'_{\max}| \geq (r/3) |e_{\max}|$ . This obviously yields the lower bound in (2.4). □

**Graph distance.** We next introduce the  $m$ th level graph distance between vertices, which will play a vital role in our further development: For any two vertices  $v', v'' \in \mathcal{T}_m$ ,  $m \geq 0$ , we define the *graph distance*  $\rho_m(v', v'')$  as the minimum number of edges from  $\mathcal{E}_m$  needed to connect  $v'$  and  $v''$ .

By the conditions on  $LR$ -triangulations, in particular condition (d), it follows that every edge in  $\mathcal{E}$  is divided at least once after  $2N_0$  steps of refinement. This immediately implies the important inequality

$$\rho_{m+2N_0\nu}(v', v'') \geq 2^\nu \rho_m(v', v''), \quad v', v'' \in \mathcal{V}_m, \quad m, \nu \geq 0. \tag{2.5}$$

The following lemma is a consequence of Lemma 2.2.

**Lemma 2.4** *There exist constants  $c > 0$  and  $t > 0$  depending only on the parameters of  $\mathcal{T}$  such that for any  $v^\diamond \in \mathcal{V}_m$*

$$\#\{v \in \mathcal{V}_m : \rho_m(v, v^\diamond) \leq n\} \leq cn^t, \quad n \geq 1. \tag{2.6}$$

Furthermore, for any  $v', v'' \in \mathcal{V}_m$  with  $\rho_m(v', v'') = n$ ,

$$\#\{v \in \mathcal{V}_m : \rho_m(v, v') + \rho_m(v, v'') = n + k\} \leq c(n + k)^t, \quad k \geq 0. \tag{2.7}$$

*Proof.* To prove (2.6) choose  $\nu \geq 1$  so that  $2^{\nu-1} \leq n < 2^\nu$ . Assume first that  $m \geq 2N_0\nu$ . Denote by  $\mathcal{T}^\diamond$  the set of all triangles  $\Delta \in \mathcal{T}_m$  which have at least one vertex in the set  $\{v \in \mathcal{V}_m : \rho_m(v, v^\diamond) \leq n\}$ . Let  $\Delta^\diamond \in \mathcal{T}^\diamond$  have  $v^\diamond$  as a vertex. Applying Lemma 2.2 with  $\Delta^\diamond$  and any other triangle in  $\mathcal{T}^\diamond$ , it readily follows that there exists a set  $\mathcal{T}^0$ , say, consisting of a triangle  $\Delta_0 \in \mathcal{T}_{m-2N_0\nu}$  and its neighbors (triangles in  $\mathcal{T}_{m-2N_0\nu}$  which share a vertex with  $\Delta_0$ ) such that  $\bigcup_{\Delta \in \mathcal{T}^\diamond} \Delta \subset \bigcup_{\Delta \in \mathcal{T}^0} \Delta$ . Evidently  $\#\mathcal{T}^0 < 3N_0$ .

Let  $\Delta_{\min}$  be a triangle in  $\mathcal{T}^\diamond$  with maximum area. By (2.1) and (2.2), we infer  $\delta r^{2N_0\nu} |\Delta_0| \leq |\Delta_{\min}|$ . We use this to obtain

$$\#\mathcal{T}^\diamond |\Delta_{\min}| \leq \sum_{\Delta \in \mathcal{T}^0} |\Delta| \leq 3N_0 \delta^{-1} |\Delta_0| \leq 3N_0 \delta^{-2} r^{-2N_0\nu} |\Delta_{\min}|$$

and hence

$$\#\{v \in \mathcal{V}_m : \rho_m(v, v^\diamond) \leq n\} \leq 3\#\mathcal{T}^\diamond \leq 9N_0 \delta^{-2} r^{-2N_0\nu} \leq cn^t$$

for some  $c, t > 0$  depending on  $N_0, r$ , and  $\delta$ .

If  $m < 2N_0\nu$  one proceeds in the same way using  $\mathcal{T}_0$  in place of  $\mathcal{T}^0$ .

Estimate (2.7) follows easily by (2.6). □

**Cells.** For any vertex  $v \in \mathcal{V}_m$  ( $m \geq 0$ ), we denote by  $\theta_v$  the union of all triangles from  $\mathcal{T}_m$  which have  $v$  as a common vertex. (Here we take into account the above observation about the multiplicities of vertices from  $\mathcal{V}_m$ .) We denote by  $\Theta_m$  the set of all such cells  $\theta_v$  with  $v \in \mathcal{V}_m$  and set  $\Theta = \bigcup_{m \geq 0} \Theta_m$ . For a given cell  $\theta \in \Theta$ , we

shall denote by  $v_\theta$  the “central” vertex of  $\theta$ . We let  $l(\theta)$  denote the *level* of  $\theta$ . Thus  $l(\theta) = m$  if  $\theta \in \Theta_m$ .

For given  $\theta', \theta'' \in \Theta_m$ , we define the *graph distance*  $\rho_m(\theta', \theta'')$  between  $\theta'$  and  $\theta''$  by  $\rho_m(\theta', \theta'') := \rho_m(v_{\theta'}, v_{\theta''})$ , where  $v_{\theta'}, v_{\theta''} \in \mathcal{V}_m$  are the “central vertices” of  $\theta'$  and  $\theta''$ . Evidently,  $\rho_m(\cdot, \cdot)$  is a true distance on  $\Theta_m$ .

Furthermore, if  $\theta', \theta'' \in \Theta_{m-1} \cup \Theta_m$  then  $v_{\theta'}, v_{\theta''} \in \mathcal{V}_m$ , and we define the *mth level graph distance*  $\rho_m(\theta', \theta'')$  between  $\theta'$  and  $\theta''$  by  $\rho_m(\theta', \theta'') := \rho_m(v_{\theta'}, v_{\theta''})$ . Evidently, if  $\Lambda \subset \Theta_{m-1} \cup \Theta_m$  consists of cells with distinct “central” points, then  $\rho_m(\cdot, \cdot)$  is a true distance on  $\Lambda$ . This will be needed in §3.2.

**Definition of  $\theta_x^m$ .** We want to associate with each  $x \in E$  a cell  $\theta_x^m \in \Theta_m$ ,  $m \geq 0$ , which contains  $x$ . Since the cells from  $\Theta_m$  overlap, this needs some care. We first associate with each triangle  $\Delta \in \mathcal{T}_m$  a cell  $\theta_\Delta^m \in \Theta_m$  such that  $\Delta \subset \theta_\Delta^m$ . Such a cell can be selected in three different ways. We choose one of them for each  $\Delta \in \mathcal{T}_m$ . Then for each  $x \in E$  such that  $x \in \Delta^\circ$  with  $\Delta \in \mathcal{T}_m$ , we define  $\theta_x^m := \theta_\Delta^m$ . If  $x$  lies on the edge of a triangle from  $\mathcal{T}_m$ , we define  $\theta_x^m$  as any cell from  $\theta_m$  such that  $x$  belongs to its interior, but if  $x = v_\theta$  for some  $\theta \in \Theta_m$ , we set  $\theta_x^m := \theta$ .

From the above definition of  $\theta_x^m$  it readily follows that for any  $\theta \in \Theta_m$  the function  $h(x) := \rho_m(\theta, \theta_x^m)$  is piecewise constant over  $\mathcal{T}_m$  and hence it is measurable, which will be needed later on.

We now introduce the *mth level graph “distance”* between any two points  $x, y \in E$  by

$$\rho_m(x, y) := \rho_m(\theta_x^m, \theta_y^m). \tag{2.8}$$

The following useful inequality is immediate from (2.5): For  $m \geq 0$ ,  $\nu \geq 1$ , and  $x, y \in E$ ,

$$\rho_{m+2N_0\nu}(x, y) \geq 2^{\nu-2} \rho_m(x, y) \quad \text{if} \quad \rho_m(x, y) \geq 3. \tag{2.9}$$

**Stars.** In order to deal with graph distances and neighborhood relations it is convenient to employ the notion of the *mth level star of a set*: For any set  $G \subset E$  and  $m \geq 0$ , we define the first *mth level star* of  $G$  by

$$\text{Star}_m(G) := \text{Star}_m^1(G) := \bigcup \{ \theta \in \Theta_m : \theta^\circ \cap G \neq \emptyset \} \tag{2.10}$$

and inductively

$$\text{Star}_m^k(G) := \text{Star}_m^1(\text{Star}_m^{k-1}(G)), \quad k > 1. \tag{2.11}$$

When  $G$  consists of a single point  $x$ , in slight abuse of notation, we shall write  $\text{Star}_m^k(x)$  instead of  $\text{Star}_m^k(\{x\})$ . For instance,  $\text{Star}_m^1(v) = \theta_v$  if  $v \in \mathcal{V}_m$ .

**Courant elements.** The no-hanging-vertices condition (e) on *LR-triangulations* guarantees the existence of *Courant elements*, that is, for every cell  $\theta \in \Theta_m$  there exists a unique continuous piecewise linear function  $\varphi_\theta$  on  $E$  which is supported on  $\theta$  and satisfies  $\varphi_\theta(v_\theta) = 1$ . We denote  $\Phi_m := \Phi_{m,\mathcal{T}} := (\varphi_\theta)_{\theta \in \Theta_m}$ .

We let  $\mathcal{S}_m$  denote the space of all continuous piecewise linear functions over  $\mathcal{T}_m$ . Clearly,  $S \in \mathcal{S}_m$  if and only if  $S = \sum_{v \in \mathcal{V}_m} S(v) \varphi_{\theta_v}$ . Evidently,  $\mathcal{S}_0 \subset \mathcal{S}_1 \subset \dots$  and by Lemma 2.3 it follows that  $\bigcup_{m \geq 0} \mathcal{S}_m = L_p(E)$ ,  $0 < p \leq \infty$ .

We shall frequently use the obvious fact that all norms of a polynomial on a triangle are equivalent, namely, if  $P$  is a polynomial of degree  $\leq k$  and  $\Delta$  is a triangle in  $\mathbb{R}^2$ , then

$$\|P\|_{L_p(\Delta)} \approx |\Delta|^{1/p-1/q} \|P\|_{L_q(\Delta)}, \quad 0 < p, q \leq \infty, \tag{2.12}$$

with constants of equivalence depending only on  $k, p$ , and  $q$ .

The  $L_p$ -stability of  $\Phi_m = (\varphi_\theta)_{\theta \in \Theta_m}$  is immediate from (2.12). In fact we shall need the following obvious modification of this fact: Let  $\Lambda \subset \Theta_{m-1} \cup \Theta_m$  consists of cells with distinct “central points”. If  $(a_\theta)_{\theta \in \Lambda}$  is an arbitrary sequence of real numbers and  $S := \sum_{\theta \in \Lambda} a_\theta \varphi_\theta$ , then

$$\|S\|_p \approx \left( \sum_{\theta \in \Lambda} \|a_\theta \varphi_\theta\|_p^p \right)^{1/p} \approx \left( \sum_{\theta \in \Lambda} |\theta| |a_\theta|^p \right)^{1/p}, \quad 0 < p \leq \infty. \tag{2.13}$$

### 2.2 Quasi-distance and maximal operators

Here we introduce a quasi-distance and maximal operators induced by  $LR$ -triangulations.

We begin by recalling the definition of a quasi-distance on a set  $X$ : The mapping  $d : X \times X \rightarrow [0, \infty)$  is called a *quasi-distance* on  $X$  if for  $x, y, z \in X$ ,

- (a)  $d(x, y) = 0 \iff x = y$ ,
- (b)  $d(y, x) = d(x, y)$ ,
- (c)  $d(x, z) \leq K(d(x, y) + d(y, z))$  with  $K \geq 1$ .

Assuming that  $\mathcal{T}$  is an  $LR$ -triangulation of a polygonal domain  $E \subset \mathbb{R}^2$ , we define the quasi-distance  $d_{\mathcal{T}} : E \times E \rightarrow [0, \infty)$  by

$$d_{\mathcal{T}}(x, y) := \min\{|\theta| : \theta \in \Theta \text{ and } x, y \in \theta\}, \tag{2.14}$$

if  $x, y$  belong to at least one cell from  $\Theta$ , and by  $d_{\mathcal{T}}(x, y) := |E|$  otherwise.

**Lemma 2.5** *The mapping  $d_{\mathcal{T}} : E \times E \rightarrow [0, \infty)$  defined above is a quasi-distance on  $E$ .*

*Proof.* Conditions (a) and (b) on quasi-distances are apparently satisfied by  $d_{\mathcal{T}}(\cdot, \cdot)$ . To prove that condition (c) holds let  $x, y, z$  be three distinct points in  $E$ . Assume that  $d(x, z) = |\theta'|$ , where  $\theta' \in \Theta_m$  is a cell containing  $x$  and  $z$ , and for some cell  $\theta'' \in \Theta_n$  containing  $y$  and  $z$  let  $d(y, z) = |\theta''|$ . Suppose  $m \leq n$ . Since  $x, z \in \theta'$ ,  $x$  and  $z$  lie in two triangles from  $\mathcal{T}_m$  with a common vertex or in the same triangle. Since  $m \leq n$ , the same is true for  $y$  and  $z$ . In other words, there exist triangles  $\Delta_1, \Delta_2 \in \mathcal{T}_m$ , which can be connected by less than  $2^2$  edges from  $\mathcal{T}_m$ , so that  $x \in \Delta_1$  and  $y \in \Delta_2$ .

Assume that  $m \geq 4N_0$ . Then by Lemma 2.2 there exists  $\theta \in \Theta_{m-4N_0}$  such that  $\Delta_1, \Delta_2 \subset \theta$  and hence  $d(x, y) \leq |\theta|$ . By (2.1) and (2.2) there exists a constant  $c$  depending on the parameters of  $\mathcal{T}$  such that  $|\theta| \leq c|\theta'|$ . Therefore,  $d(x, y) \leq c(d(x, z) + d(z, y))$ . If  $m < 4N_0$ , we use  $E$  instead of  $\theta$  with the same result.  $\square$

With the following lemma we relate the quasi-distance  $d_{\mathcal{T}}(\cdot, \cdot)$  to the  $m$ th level graph distance introduced in §2.1.

**Lemma 2.6** *There exist constants  $\beta > 0$  and  $c > 0$  such that for  $\theta \in \Theta_m$  ( $m \geq 0$ ) and  $x \in E$ ,*

$$d_{\mathcal{T}}(v_{\theta}, x) \leq c|\theta| \rho_m(\theta, \theta_x^m)^{\beta} \quad \text{if} \quad \rho_m(\theta, \theta_x^m) \geq 2. \tag{2.15}$$

*Proof.* Clearly, if we prove (2.15) when  $x = v'$  with  $v' \in \mathcal{V}_m$ , then it will hold in general with a different constant  $c$ . Let  $d_{\mathcal{T}}(v_{\theta}, v') = |\omega|$  with  $\omega \in \Theta_{\ell}$  ( $\ell \geq 0$ ), i.e., the cell  $\omega$  is of minimum area and  $v_{\theta}, v' \in \omega$ . Let  $\omega^* \in \Omega$  be of maximum level such that  $v_{\theta}, v' \in \omega^*$ . Evidently,  $|\omega| \leq |\omega^*|$ . Since  $\omega \cap \omega^* \neq \emptyset$ , it follows from (2.1) and (2.2) that

$$l(\omega) \leq l(\omega^*) \leq l(\omega) + \nu_0$$

where  $\nu_0 > 0$  is a constant depending only on the parameters of  $\mathcal{T}$ .

Our next claim is that  $\rho_m(v_{\theta}, v') \geq 2^n$ , where  $n := [(m - l(\omega^*)) / 2N_0] + 1$ . Assume to the contrary that  $\rho_m(v_{\theta}, v') < 2^n$ . Then by Lemma 2.2 it follows that there is  $\eta \in \Theta_{m-2N_0n}$  such that  $v_{\theta}, v' \in \eta$  and hence  $m - 2N_0n \geq l(\omega^*)$  or  $n \leq (m - l(\omega^*)) / 2N_0$ , which is a contradiction. Therefore,

$$\rho_m(v_{\theta}, v') \geq 2^{(m-l(\omega^*)) / 2N_0} \geq 2^{(m-l(\omega)-\nu_0) / 2N_0} \geq c 2^{(m-l(\omega)) / 2N_0}. \tag{2.16}$$

Evidently, by (2.1) and (2.2),  $|\omega| / |\theta| \leq c(1/r)^{m-l(\omega)}$ . Combining this with (2.16), there exists  $\beta > 1$  such that

$$d_{\mathcal{T}}(v_{\theta}, v') = |\omega| \leq c|\theta| \left( 2^{(m-l(\omega)) / 2N_0} \right)^{\beta} \leq c|\theta| \rho_m(v_{\theta}, v')^{\beta}.$$

Above it may happen that  $\omega = \theta_0 := E$ . The proof is the same.  $\square$

The quasi-distance  $d_{\mathcal{T}}(\cdot, \cdot)$  induces a maximal operator. Denote by  $B(y, a)$  the “ball” centered at  $y$  of radius  $a > 0$  with respect to this quasi-distance. Then for any  $s > 0$  the maximal operator  $\mathcal{M}_{d_{\mathcal{T}}}^s$  is defined by

$$(\mathcal{M}_{d_{\mathcal{T}}}^s f)(x) := \sup_{B: x \in B} \left( \frac{1}{|B|} \int_B |f(y)|^s dy \right)^{1/s}, \quad x \in E, \tag{2.17}$$

where the supremum is over all balls  $B$  containing  $x$ .

For our purposes it is more convenient to use the equivalent maximal operator  $\mathcal{M}_{\mathcal{T}}^s$  defined by

$$(\mathcal{M}_{\mathcal{T}}^s f)(x) := \sup_{\theta: x \in \theta} \left( \frac{1}{|\theta|} \int_{\theta} |f(y)|^s dy \right)^{1/s} \tag{2.18}$$

where the supremum is over all cells  $\theta \in \Theta$  containing  $x$  or  $\theta = E$ .

**Lemma 2.7** For any measurable function  $f$

$$\mathcal{M}_{\mathcal{T}}^s f(x) \approx \mathcal{M}_{d_{\mathcal{T}}}^s f(x), \quad x \in \mathbb{R}^2, \tag{2.19}$$

where the constants of equivalence depend only on  $s$  and the parameters of  $\mathcal{T}$ .

This equivalence is immediate from the following lemma which will be needed later on as well.

**Lemma 2.8** (a) Given a ball  $B = B(x, r)$ ,  $x \in E$ ,  $r > 0$ , there exist  $\theta' \in \Theta$  and  $\theta'' \in \Theta$  or  $\theta'' = E$  such that

$$\theta' \subset B \subset \theta'' \quad \text{and} \quad r \leq |\theta''| \leq c_1 |B| \leq c_2 |\theta'| < c_2 r. \tag{2.20}$$

(b) For any  $\theta \in \Theta$  there exists a ball  $B \subset E$  (with respect to the quasi-distance  $d_{\mathcal{T}}$ ) such that

$$\theta \subset B \quad \text{and} \quad |B| \leq c |\theta| \tag{2.21}$$

Here the constants depend only on the parameters of  $\mathcal{T}$ .

*Proof.* (a) Fix a ball  $B = B(x, r)$ ,  $x \in E$ ,  $r > 0$ . Let  $\theta'$  be a cell of minimal level, say  $m$ , such that  $x \in \theta' \subset B$ . Clearly,

$$B \subset \bigcup_{\theta \in \Theta_m: x \in \theta} \theta \subset \text{Star}_m^2(x)$$

and by Lemma 2.2 there exists  $\theta'' \in \Theta_{m-4N_0}$  or  $\theta'' = E$  if  $m < 4N_0$  such that

$$\theta' \subset B \subset \text{Star}_m^2(x) \subset \theta''.$$

By properties (f) and (g) of LR-triangulations (§2.1), it follows that  $|\theta''| \leq c |\theta'|$ . Evidently  $|\theta'| < r$  and  $|\theta''| \geq r$ , and (2.20) follows.

(b) Suppose that  $\theta \in \Theta_n$  ( $n \geq 0$ ) is with “central” vertex  $v$ . Let  $\delta := \max \{|\theta| : \theta \subset \text{Star}_n^2(v)\}$ . Then for sufficiently small  $\varepsilon > 0$ ,

$$\theta \subset B(v, \delta + \varepsilon) = \bigcup_{|\tilde{\theta}| < \delta + \varepsilon: v \in \tilde{\theta}} \tilde{\theta} \subset \text{Star}_n^2(v),$$

which yields  $|B(v, \delta + \varepsilon)| \leq c |\theta|$ . □

We now come to the main point in this subsection. It is well-known that the Fefferman–Stein vector valued maximal inequality holds for maximal functions generalized by quasi-distances as in our case (see [15]). This combined with Lemma 2.7 gives the needed maximal inequality:

**Proposition 2.9** Let  $\mathcal{T}$  be an LR-triangulation of  $E \subset \mathbb{R}^2$ . If  $0 < p < \infty$ ,  $0 < q \leq \infty$ , and  $0 < s < \min\{p, q\}$ , then for any sequence of functions  $(f_j)_{j=1}^{\infty}$  on  $E$ ,

$$\left\| \left( \sum_{j=1}^{\infty} |\mathcal{M}_{\mathcal{T}}^s f_j|^q \right)^{1/q} \right\|_p \leq c \left\| \left( \sum_{j=1}^{\infty} |f_j|^q \right)^{1/q} \right\|_p, \tag{2.22}$$

where  $c$  depends only on  $p, q, s$ , and the parameters of  $\mathcal{T}$ .

### 2.3 Spaces of homogeneous type on polygonal domains

Spaces of homogeneous type were first introduced in [2] as a means to extend the Calderon–Zygmund theory of singular integral operators to more general settings.

Let  $X$  be a topological space endowed with a Borel measure  $\mu$  and a quasi-distance  $d(\cdot, \cdot)$  (see §2.2). Assume that the balls  $B(x, r) := \{y \in X : d(x, y) < r\}$ ,  $x \in X$ ,  $r > 0$ , form a basis for the topology  $T$  in  $X$ , and  $\mu(B(x, r)) > 0$  if  $r > 0$ . The space  $(X, d, \mu)$  is said to be of *homogeneous type* if there exists a constant  $A$  such that for all  $x \in X$  and  $r > 0$ ,

$$0 < \mu(B(x, 2r)) \leq A\mu(B(x, r)). \tag{2.23}$$

**The space of homogeneous type  $(E, d_{\mathcal{T}}, m)$ .** Suppose that  $E$  is a bounded polygonal domain and let  $\mathcal{T}$  be a  $LR$ -triangulation on  $E$ . Also, let  $d_{\mathcal{T}}(\cdot, \cdot)$  be the quasi-distance on  $E$  defined in (2.14). Finally, denote by  $m$  the Lebesgue measure on  $E$ . It is easy to see that  $(E, d_{\mathcal{T}}, m)$  is a space of homogeneous type, so that we can utilize the machinery developed in [2]. Indeed, by Lemma 2.5,  $d_{\mathcal{T}}(\cdot, \cdot)$  is a quasi-distance on  $E$  and evidently  $m(B(x, r)) = |B(x, r)| > 0$  for  $x \in E$  and  $r > 0$ . Further, it follows by Lemma 2.8 that condition (2.23) is fulfilled as well.

**The Hardy space  $H_1(E, \mathcal{T})$ .** We next define the Hardy space  $H_1 := H_1(E, \mathcal{T})$  associated with the space  $(E, d_{\mathcal{T}}, m)$  by means of atomic representations (see [3]).

According to Coifmann and Weiss [3], a function  $a(x)$  is said to be a  $q$ -atom ( $1 < q \leq \infty$ ) if there exist  $x_0 \in E$  and  $r > 0$  such that

- (i)  $\text{supp } a \subset B(x_0, r)$ ,
- (ii)  $\|a\|_q \leq |B(x_0, r)|^{1/q-1}$ ,
- (iii)  $\int a(x) dx = 0$ .

In addition,  $|E|^{-1} \mathbb{1}_E$  is by definition a  $q$ -atom as well.

We adopt the following slightly different but equivalent definition for a  $q$ -atom which better suits our purposes.

**Definition 2.10** A function  $a(x)$  is said to be a  $q$ -atom ( $1 < q \leq \infty$ ) for  $H_1(E, \mathcal{T})$  if there is  $\theta \in \Theta$  or  $\theta = E$  such that

- (a)  $\text{supp } a \subset \theta$ ,
- (b)  $\|a\|_q \leq |\theta|^{1/q-1}$ ,
- (c)  $\int_E a(x) dx = 0$ .

We also postulate  $|E|^{-1} \mathbb{1}_E$  to be a  $q$ -atom.

The equivalence of the two definitions for a  $q$ -atom is immediate by Lemma 2.8 (a).

**Definition 2.11** The space  $H_1^q := H_1^q(E, \mathcal{T})$  ( $1 < q \leq \infty$ ) is defined to be the set of all functions  $f \in L_1(E)$  admitting an atomic decomposition

$$f = \sum_{j=0}^{\infty} \lambda_j a_j,$$

where the  $a_j$ 's are  $q$ -atoms and  $\sum_{j=0}^{\infty} |\lambda_j| < \infty$ . Moreover, the norm of  $f \in H_1^q$  is given by

$$\|f\|_{H_1^q} := \inf \left\{ \sum_{j=0}^{\infty} |\lambda_j| : f = \sum_{j=0}^{\infty} \lambda_j a_j, a_j \text{ } q\text{-atoms} \right\}.$$

A fundamental fact in the theory of Hardy spaces is that  $H_1^q = H_1^\infty$  whenever  $1 < q \leq \infty$  with equivalent norms (see [3, Theorem A]). Thus all spaces  $H_1^q$  are the same and we shall drop the index  $q$ . In the following we shall only work with the norm in  $H_1$  defined by using 2-atoms.

An important fact is that the spaces  $H_1(E, \mathcal{T})$  essentially depend on the triangulations  $\mathcal{T}$ . We call  $H_1(E, \mathcal{T}^*)$  a *regular  $H_1$ -space* if  $\mathcal{T}^*$  is a regular multilevel triangulation of  $E$  (see §2.1). It is readily seen that if  $H_1(E, \mathcal{T}^*)$  is regular, then it is the same (with equivalent norms) as the space  $H_1(E)$  defined using atoms generated by the Euclidean distance on  $E$ . Thus all regular spaces  $H_1(E, \mathcal{T})$  are the same. Consider the case when  $E := [-1, 1]^2$  and denote by  $H_1(E)$  the regular  $H_1$ -space on  $[-1, 1]^2$ . As will be shown in the appendix, there exists an  $LR$ -triangulation  $\mathcal{T}$  such that  $H_1(E, \mathcal{T}) \neq H_1(E)$ . The reason for this is that there exist  $LR$ -triangulations on



$[-1, 1]^2$  containing triangles with uncontrollably sharp angles (see §2.1). The fact that the spaces  $H_1(E, \mathcal{T})$ , where  $\mathcal{T}$  is allowed to vary, are not all the same is not a surprise since as is well-known the norm in  $H_1(\mathbb{R}^d)$  ( $d > 1$ ) is not invariant (like the  $L_1$ -norm) under linear transforms with determinant one. We do not explore in more detail the relationship between the various spaces  $H_1(E, \mathcal{T})$  in this article.

It is not hard to prove that  $H_1(E, \mathcal{T})$  is a Banach space and  $\|f\|_{L_1(E)} \leq c \|f\|_{H_1(E, \mathcal{T})}$  for  $f \in H_1(E, \mathcal{T})$ .

Another fundamental result is that the dual of  $H_1(E, \mathcal{T})$  is the space  $BMO := BMO(E, \mathcal{T})$  which can be defined in our case as the set of all functions  $f$  on  $E$  such that

$$\|f\|_{BMO} := \left| \int_E f(x) dx \right| + \sup_{\theta} \left( \frac{1}{|\theta|} \int_{\theta} |f(x) - f_{\theta}|^2 dx \right)^{1/2} < \infty, \tag{2.24}$$

where  $f_{\theta} := \frac{1}{|\theta|} \int_{\theta} f(x) dx$  and the supremum is taken over all  $\theta \in \Theta$  or  $\theta = E$ . More precisely, for  $g \in BMO(E, \mathcal{T})$  and  $f \in H_1(E, \mathcal{T})$  with an atomic decomposition  $f = \sum_{j=1}^{\infty} \lambda_j a_j$ ,

$$\lim_{m \rightarrow \infty} \sum_{j=1}^m \lambda_j \int_E g(x) a_j(x) dx \tag{2.25}$$

defines a continuous linear functional on  $H_1$  whose norm is equivalent to  $\|g\|_{BMO}$  and vice versa each continuous linear functional on  $H_1$  is of this form.

Note that an equivalent norm in  $BMO(E, \mathcal{T})$  can be defined by replacing in (2.24)  $\left( \frac{1}{|\theta|} \int_{\theta} |f(x) - f_{\theta}|^2 dx \right)^{1/2}$  by  $\frac{1}{|\theta|} \int_{\theta} |f(x) - f_{\theta}| dx$ . For more details, see [3].

Finally, we observe that since  $H_1(E, \mathcal{T}) \neq H_1(E)$  for some LR-triangulations  $\mathcal{T}$ , then by a duality argument it follows that for the same triangulations  $BMO(E, \mathcal{T}) \neq BMO(E)$ , where  $BMO(E)$  stands for the regular BMO space on  $E$ . Thus in general  $BMO(E, \mathcal{T})$  depends on the triangulation  $\mathcal{T}$ .

One of the advantages of introducing  $H_1$  via atomic decompositions is that questions related to the boundedness of Calderon–Zygmund operators (CZO) on  $H_1$  can be answered by focusing on individual atoms. Evidently, any operator  $T$  would be bounded if  $T$  maps atoms into atoms. Coifman and Weiss observed that for certain type of operators  $T$ , for every atom  $a(x)$ ,  $Ta$  is a function with similar structure, which they term a *molecule*. We shall use the following definition of a molecule.

**Definition 2.12** For a given  $\varepsilon > 0$ , we say that  $m(x)$  is an  $\varepsilon$ -molecule for  $H_1(E, \mathcal{T})$  centered at  $x_0 \in E$  if

$$\left( \int_E |m(x)|^2 dx \right) \left( \int_E |m(x)|^2 d_{\mathcal{T}}(x, x_0)^{1+\varepsilon} dx \right)^{1/\varepsilon} \leq 1 \quad \text{and} \quad \int_E m(x) dx = 0. \tag{2.26}$$

It is trivial to see that every 2-atom is an  $\varepsilon$ -molecule for any  $\varepsilon > 0$ . More importantly,  $\|m\|_{H_1} \leq c$  for each  $\varepsilon$ -molecule  $m(x)$  (see [3, Theorem C]). From this it follows that a linear operator mapping atoms into molecules has a bounded extension to  $H_1$ .

The following result [3] will play an important role in our further development:

**Proposition 2.13** Let  $T : L_2(E) \rightarrow L_2(E)$  be a bounded linear operator given by

$$(Tf)(x) = \int_E K(x, y) f(y) dy.$$

Suppose that for each 2-atom  $a \neq |E|^{-1} \mathbb{1}_E$

$$\int_E (Ta)(x) dx = 0 \tag{2.27}$$

and there is  $\varepsilon > 0$  such that for  $d(x, y_0) > c d(y, y_0)$  the kernel  $K(\cdot, \cdot)$  satisfies

$$|K(x, y) - K(x, y_0)| \leq c \left( \frac{d(y, y_0)}{d(x, y_0)} \right)^{\varepsilon} \frac{1}{d(x, y_0)}. \tag{2.28}$$

Then  $Ta$  is a constant multiple of an  $\varepsilon$ -molecule for any atom  $a \neq |E|^{-1} \mathbb{1}_E$ .

### 3 Anisotropic Franklin bases

In this section we explore the Franklin system  $\mathcal{F}_T$  generated by the Courant elements associated with an arbitrary locally regular triangulation  $T$  of a compact polygonal domain  $E$  in  $\mathbb{R}^2$ . We shall show that each such Franklin system is a Schauder basis for  $C(E)$  and  $L_1(E)$ , and it is an unconditional basis for  $H_1(E, T)$  and  $L_p(E)$  ( $1 < p < \infty$ ). We also prove some related results.

#### 3.1 Definition of the Franklin system. Main results

Throughout this section, we assume that  $T = \bigcup_{m=0}^\infty \mathcal{T}_m$  is an  $LR$ -triangulation of  $E$ . We recall that  $\mathcal{V}_m$  denotes the set of all vertices of triangles from  $\mathcal{T}_m$ . We set  $\mathcal{V}_0^* = \mathcal{V}_0$  and  $\mathcal{V}_m^* = \mathcal{V}_m \setminus \mathcal{V}_{m-1}$  for  $m \geq 1$  and write  $\mathcal{V}^* = \bigcup_{m=0}^\infty \mathcal{V}_m^*$ .

Let  $\theta_0 := E$ . Choose  $\theta_{\max} \in \Theta_0$  to be of maximum area and denote  $\Theta_0^* := \{\theta_0\} \cup \Theta_0 \setminus \{\theta_{\max}\}$ , i.e., we replace  $\theta_{\max}$  by  $\theta_0 = E$ . Moreover, we associate  $\theta_0$  with  $v_{\theta_{\max}}$  and set  $\varphi_{\theta_0} := \mathbb{1}_{\theta_0}$ . For  $m \geq 1$  denote by  $\Theta_m^*$  the set of all cells  $\theta \in \Theta_m$  with ‘‘central’’ vertices  $v_\theta \in \mathcal{V}_m^*$  and set  $\Theta^* := \bigcup_{m=0}^\infty \Theta_m^*$ .

Note that for each  $m$ , the set  $\{\varphi_\theta : \theta \in \bigcup_{i=0}^m \Theta_i^*\}$  is linearly independent. Also,

$$S_m = \text{span}\{\varphi_\theta : \theta \in \Theta_m\} = \text{span}\left\{\varphi_\theta : \theta \in \bigcup_{i=0}^m \Theta_i^*\right\}.$$

For  $\theta \in \Theta$  we denote by  $\tilde{\varphi}_\theta$  the  $L_2$ -normalized version of the Courant element  $\varphi_\theta$ , i.e.,  $\tilde{\varphi}_\theta := \|\varphi_\theta\|_2^{-1} \varphi_\theta \approx |\theta|^{-1/2} \varphi_\theta$ .

We consider an arbitrary (but fixed) linear order  $\preceq$  on  $\Theta^*$  satisfying the following conditions:

- (i) If  $\theta \in \Theta_m^*$  and  $\theta' \in \Theta_n^*$  with  $m < n$ , then  $\theta \preceq \theta'$ ,
  - (ii)  $\theta_0 \preceq \theta$ , for all  $\theta \in \Theta^*$ .
- (3.1)

We now define the *Franklin system*  $\mathcal{F}_T$  by applying the Gram–Schmidt orthogonalization process to  $\{\tilde{\varphi}_\theta\}_{\theta \in \Theta^*}$  in  $L_2(E)$  with respect to the order  $\preceq$ . We obtain an orthonormal system  $\mathcal{F}_T := \{f_\theta\}_{\theta \in \Theta^*}$  in  $L_2(E)$  consisting of continuous piecewise linear functions. Each Franklin function  $f_\theta$  is uniquely determined (up to a multiple  $\pm 1$ ) by the conditions:

- (a)  $f_\theta \in \text{span}\{\varphi_{\theta'} : \theta' \preceq \theta\}$ .
- (b)  $\langle f_\theta, \varphi_{\theta'} \rangle = 0$  for all  $\theta' \prec \theta$ ,
- (c)  $\|f_\theta\|_2 = 1$ .

Note that  $f_{\theta_0} = \pm \tilde{\mathbb{1}}_{\theta_0} := \pm |E|^{-1/2} \mathbb{1}_E$ .

We next state our main results on Franklin systems  $\mathcal{F}_T$ , where  $T$  is an arbitrary  $LR$ -triangulation of a bounded polygonal domain  $E \subset \mathbb{R}^2$ .

**Theorem 3.1** *The Franklin system  $\mathcal{F}_T := \{f_\theta\}_{\theta \in \Theta^*}$  is a Schauder basis for  $L_p(E)$ ,  $1 \leq p \leq \infty$ , with  $L_\infty(E) := C(E)$ .*

**Theorem 3.2** *The Franklin system  $\mathcal{F}_T := \{f_\theta\}_{\theta \in \Theta^*}$  is an unconditional basis for  $H_1(E, T)$  and  $L_p(E)$ ,  $1 < p < \infty$ .*

**Theorem 3.3** *The following conditions are equivalent:*

- (a)  $f \in H_1(E, T)$ ;
- (b) *The series  $\sum_{\theta \in \Theta^*} \langle f, f_\theta \rangle f_\theta$  converges unconditionally in  $L_1$ ;*
- (c)  $S_f(x) := \left(\sum_{\theta \in \Theta^*} |\langle f, f_\theta \rangle|^2 |f_\theta(x)|^2\right)^{1/2} \in L_1$ ;
- (d)  $F_f(x) := \left(\sum_{\theta \in \Theta^*} |\langle f, f_\theta \rangle|^2 |\tilde{\mathbb{1}}_\theta(x)|^2\right)^{1/2} \in L_1$ .

Furthermore, if  $f \in H_1(E, T)$ , then

$$\|f\|_{H_1} \approx \|S_f\|_{L_1} \approx \|F_f\|_{L_1}. \tag{3.2}$$

**Theorem 3.4** A function  $f \in BMO(E, \mathcal{T})$  if and only if

$$\sup_{\theta} \left( \frac{1}{|\theta|} \sum_{\eta \in \Theta^*: \eta \subset \theta} |\langle f, f_{\eta} \rangle|^2 \right)^{1/2} < \infty, \tag{3.3}$$

where the supremum is taken over all  $\theta \in \Theta$  or  $\theta = E$ . Furthermore,  $\|f\|_{BMO(E, \mathcal{T})}$  is equivalent to the quantity in (3.3).

**3.2 Representation of the Franklin functions and proof of Theorem 3.1**

The exponential decay of the Franklin functions is a central issue in the study of Franklin systems. We begin with a generalization of the well-known result of Demko [6] on the inverses of band matrices, given in [13].

**Proposition 3.5** Suppose  $K$  is a finite set of indices and let  $\rho$  be a distance on  $K$ . Let  $A = [a_{k,l}]_{k,l \in K}$  be an invertible band matrix of order  $r \geq 1$ , i.e.,  $a_{k,l} = 0$  if  $\rho(k,l) > r$ . Let  $A^{-1} = [b_{k,l}]_{k,l \in K}$  be the inverse matrix of  $A$ . Suppose that for some  $1 \leq p \leq \infty$ ,

$$\|A\|_{\ell_p(K) \rightarrow \ell_p(K)} \leq M_1 \quad \text{and} \quad \|A^{-1}\|_{\ell_p(K) \rightarrow \ell_p(K)} \leq M_2.$$

Then there exist constants  $c > 0$  and  $0 < q < 1$  depending only on  $M_1, M_2, r$ , and  $p$  such that

$$|b_{k,l}| \leq cq^{\rho(k,l)} \quad \text{for } k, l \in K.$$

For any  $\eta \in \Theta_m^*$  ( $m \geq 0$ ), denote

$$\Lambda_{\eta} := \Theta_{m-1} \cup \{\theta \in \Theta_m^* : \theta \leq \eta\}. \tag{3.4}$$

Note that the cells  $\theta \in \Lambda_{\eta}$  have distinct ‘‘central’’ points  $v_{\theta}$  and hence the set  $\{\varphi_{\theta} : \theta \in \Lambda_{\eta}\}$  is linearly independent. Let  $G_{\eta}$  be the Gram matrix given by

$$G_{\eta} = [a_{\theta\theta'}]_{\theta, \theta' \in \Lambda_{\eta}} \quad \text{with} \quad a_{\theta\theta'} := \langle \tilde{\varphi}_{\theta}, \tilde{\varphi}_{\theta'} \rangle, \tag{3.5}$$

and denote  $G_{\eta}^{-1} =: [b_{\theta\theta'}]_{\theta, \theta' \in \Lambda_{\eta}}$ .

**Lemma 3.6** There exist constants  $0 < q < 1$  and  $c > 0$  such that for any  $\eta \in \Theta_m^*$  ( $m \geq 0$ ) we have, for the entries of  $G_{\eta}^{-1}$ , the following estimate

$$|b_{\theta\theta'}| \leq cq^{\rho_m(\theta, \theta')}, \quad \theta, \theta' \in \Lambda_{\eta}, \tag{3.6}$$

where  $\rho_m(\theta, \theta')$  is the  $m$ th level graph distance between  $\theta$  and  $\theta'$ , introduced in §2.1.

*Proof.* By condition (c) on LR-triangulations, every triangle  $\Delta \in \mathcal{T}$  has at most  $M_0$  children. Then  $a_{\theta\theta'} := \langle \tilde{\varphi}_{\theta}, \tilde{\varphi}_{\theta'} \rangle = 0$  if  $\rho_m(\theta, \theta') > 2M_0$  and hence the Gram matrix  $G_{\eta}$  is  $r$ -banded with  $r := 2M_0 + 1$ . Since  $G_{\eta}$  is symmetric, then

$$\begin{aligned} \|G_{\eta}\|_{\ell_2(\Lambda_{\eta}) \rightarrow \ell_2(\Lambda_{\eta})} &= \max\{\lambda : \lambda \text{ eigenvalue of } G_{\eta}\} \quad \text{and} \\ \|G_{\eta}^{-1}\|_{\ell_2(\Lambda_{\eta}) \rightarrow \ell_2(\Lambda_{\eta})} &= \max\{1/\lambda : \lambda \text{ eigenvalue of } G_{\eta}\}. \end{aligned}$$

On the other hand, for any vector  $x := (x_{\theta})_{\theta \in \Lambda_{\eta}}$  we have  $\|\sum_{\theta \in \Lambda_{\eta}} x_{\theta} \tilde{\varphi}_{\theta}\|_{L_2(E)}^2 = \langle Gx, x \rangle$  and thus by (2.13) there exist constants  $c_1, c_2 > 0$  such that

$$c_1 \|x\|_{\ell_2(\Lambda_{\eta})}^2 \leq \langle G_{\eta}x, x \rangle \leq c_2 \|x\|_{\ell_2(\Lambda_{\eta})}^2,$$

which implies  $c_1 \leq \lambda \leq c_2$  for every eigenvalue  $\lambda$ . Hence

$$\|G_{\eta}\|_{\ell_2(\Lambda_{\eta}) \rightarrow \ell_2(\Lambda_{\eta})} \leq c_2 \quad \text{and} \quad \|G_{\eta}^{-1}\|_{\ell_2(\Lambda_{\eta}) \rightarrow \ell_2(\Lambda_{\eta})} \leq c_1^{-1}.$$

Finally, note that since the ‘‘central’’ points of the cells from  $\Lambda_{\eta}$  are distinct points in  $\mathcal{V}_m$ , the  $m$ th level graph distance  $\rho_m(\cdot, \cdot)$  is a true distance on  $\Lambda_{\eta}$ . Thus  $G_{\eta}$  satisfies the conditions of Lemma 3.5 with  $\rho(\cdot, \cdot) = \rho_m(\cdot, \cdot)$  and hence (3.6) holds.  $\square$

We are now prepared to deduce an important representation of the Franklin functions.

**Lemma 3.7** For any  $\theta \in \Theta_m^*$  ( $m \geq 0$ ) the Franklin function  $f_\theta$  has a representation of the form

$$f_\theta = \sum_{\eta \in \Theta_m} c_{\theta\eta} \tilde{\varphi}_\eta, \tag{3.7}$$

with coefficients  $c_{\theta\eta}$  satisfying

$$|c_{\theta\eta}| \leq c q^{\rho_m(\theta,\eta)}, \quad \eta \in \Theta_m, \tag{3.8}$$

where the constants  $0 < q < 1$  and  $c > 0$  depend only on the parameters of  $\mathcal{T}$ .

*Proof.* It is readily seen (and well-known) that the function  $g_\theta$  defined by

$$g_\theta := \sum_{\xi \in \Lambda_\theta} b_{\theta\xi} \tilde{\varphi}_\xi = \sum_{\xi \in \Theta_{m-1}} b_{\theta\xi} \tilde{\varphi}_\xi + \sum_{\xi \in \Theta_m^* : \xi \preceq \theta} b_{\theta\xi} \tilde{\varphi}_\xi, \tag{3.9}$$

where  $b_{\theta\xi}$  are entries of  $G_\theta^{-1}$ , has the property  $\langle g_\theta, \tilde{\varphi}_\eta \rangle = \delta_{\theta\eta}$ . Here, we set  $b_{\theta\xi} := 0$  if  $\xi \in \Theta_{-1}$ . Therefore,  $f_\theta = \pm \|g_\theta\|_2^{-1} g_\theta$ . Evidently, for  $\xi \in \Theta_{m-1}$ ,

$$\tilde{\varphi}_\xi = \sum_{\eta \in \Theta_m : \eta \subset \xi} \tilde{\varphi}_\xi(v_\eta) \|\varphi_\eta\|_2 \tilde{\varphi}_\eta.$$

Substituting this in (3.9) and switching the order of summation, we arrive at

$$f_\theta = \sum_{\eta \in \Theta_m} c_{\theta\eta} \tilde{\varphi}_\eta,$$

where

$$c_{\theta\eta} = \begin{cases} \|g_\theta\|_2^{-1} \|\varphi_\eta\|_2 \sum_{\xi \in \Theta_{m-1} : \eta \subset \xi} b_{\theta\xi} \tilde{\varphi}_\xi(v_\eta) & \text{if } \eta \in \Theta_m \setminus \Lambda_\theta, \\ \|g_\theta\|_2^{-1} \left( b_{\theta\eta} + \|\varphi_\eta\|_2 \sum_{\xi \in \Theta_{m-1} : \eta \subset \xi} b_{\theta\xi} \tilde{\varphi}_\xi(v_\eta) \right) & \text{if } \eta \in \Theta_m \cap \Lambda_\theta. \end{cases}$$

Note that

$$1 = \langle g_\theta, \tilde{\varphi}_\theta \rangle \leq \|g_\theta\|_2 \|\tilde{\varphi}_\theta\|_2 = \|g_\theta\|_2$$

and  $\rho_m(\eta, \xi) \leq M_0$ , whenever  $\eta \in \Theta_{m-1}$ ,  $\xi \in \Theta_m$ , and  $\xi \subset \eta$ . Also,  $\|\tilde{\varphi}_\xi\|_\infty = \|\varphi_\xi\|_2^{-1} \approx |\xi|^{-1/2}$ . We use these along with properties (2.1) and (2.2) of LR-triangulations, and the estimate for  $|b_{\theta\eta}|$  from Lemma 3.6 to obtain, for  $\eta \in \Theta_m \cap \Lambda_\theta$ ,

$$\begin{aligned} |c_{\theta\eta}| &\leq |b_{\theta\eta}| + c |\eta|^{1/2} \sum_{\xi \in \Theta_{m-1} : \eta \subset \xi} |b_{\theta\xi}| \tilde{\varphi}_\xi(v_\eta) \\ &\leq c q^{\rho_m(\theta,\eta)} + c (|\eta|/|\xi|)^{1/2} \sum_{\xi \in \Theta_{m-1} : \eta \subset \xi} q^{\rho_m(\theta,\xi)} \leq c q^{\rho_m(\theta,\eta)}. \end{aligned}$$

The estimate of  $|c_{\theta\eta}|$  when  $\eta \in \Theta_m \setminus \Lambda_\theta$  is the same. □

*Proof of Theorem 3.1.* As was already mentioned in §2.1,  $\overline{\bigcup_{m=0}^\infty \mathcal{S}_m} = L_p(E)$ ,  $1 \leq p \leq \infty$ , and hence  $\mathcal{F}_\mathcal{T}$  is dense in  $L_p(E)$ . It remains to prove that the orthogonal projector operator  $P_\eta f := \sum_{\theta \preceq \eta} \langle f, f_\theta \rangle f_\theta$  is bounded on  $L_p(E)$  ( $1 \leq p \leq \infty$ ), i.e., considered as an operator from  $L_p$  into  $L_p$ . It is easy to see that

$$P_\eta f(x) = \int_E \sum_{\theta, \theta' \in \Lambda_\eta} b_{\theta, \theta'} \tilde{\varphi}_\theta(x) \tilde{\varphi}_{\theta'}(y) f(y) dy, \quad \Lambda_\eta := \Theta_{m-1} \cup \{\theta \in \Theta_m^* : \theta \preceq \eta\}.$$

Denote  $K_\eta(x, y) := \sum_{\theta, \theta' \in \Lambda_\eta} b_{\theta, \theta'} \tilde{\varphi}_\theta(x) \tilde{\varphi}_{\theta'}(y)$ . Then for each  $x \in E$ ,

$$\|K_\eta(x, \cdot)\|_{L_1} \leq \sum_{\theta, \theta' \in \Lambda_\eta} |b_{\theta, \theta'}| |\tilde{\varphi}_\theta(x)| \|\tilde{\varphi}_{\theta'}\|_{L_1} \leq c \sum_{\theta, \theta' \in \Lambda_\eta} q^{\rho_m(\theta, \theta')} (|\theta|/|\theta'|)^{1/2} \mathbb{1}_\theta(x),$$

where we used (3.6) and that  $\|\tilde{\varphi}_\theta\|_2 = 1$ . Our goal is to show that  $\|K_\eta(x, \cdot)\|_{L_1} \leq C < \infty$ . Fix  $x \in E$ . Since there are at most 6 cells  $\theta \in \Lambda_\eta$  such that  $x \in \theta^\circ$ , it suffices to show that for each  $\theta \in \Lambda_\eta$ ,  $\sum_{\theta' \in \Lambda_\eta} q^{\rho_m(\theta, \theta')} (|\theta|/|\theta'|)^{1/2} \leq c < \infty$ . By (2.3) it follows that if  $\rho_m(\theta, \theta') = n$  ( $n \geq 1$ ), then  $|\theta|/|\theta'| \leq cn^s$ . On the other hand, Lemma 2.4 yields that  $\#\{\theta' \in \Lambda_\eta : \rho_m(\theta, \theta') = n\} \leq cn^t$ . Using the above, we obtain

$$\sum_{\theta' \in \Lambda_\eta} q^{\rho_m(\theta, \theta')} (|\theta|/|\theta'|)^{1/2} \leq c \sum_{n=1}^\infty q^n n^{t+s/2} \leq c < \infty$$

since  $0 < q < 1$ . Consequently,  $\|K_\eta(x, \cdot)\|_{L_1} \leq C$ . This estimate implies

$$\|P_\eta\|_{L_1 \rightarrow L_1} < \infty \quad \text{and} \quad \|P_\eta\|_{L_\infty \rightarrow L_\infty} < \infty.$$

By interpolation it follows that  $\|P_\eta\|_{L_p \rightarrow L_p} < \infty$ ,  $1 < p < \infty$ . □

### 3.3 Localization and smoothness of the Franklin functions

Here we show that the Franklin functions belong to  $\text{Lip } \varepsilon$  (for some  $\varepsilon > 0$ ) with respect to the quasi-distance  $d_{\mathcal{T}}(\cdot, \cdot)$  introduced in (2.14) and have exponential rate of decay with respect to the corresponding graph distance.

We shall systematically use the notation introduced in §2.1. We recall, in particular, that for  $x \in E$ ,  $\theta_x^m$  is a cell from  $\Theta_m$  containing  $x$ , and  $\rho_m(\cdot, \cdot)$  is the  $m$ th level graph distance.

**Theorem 3.8** *There exist constants  $\varepsilon > 0$ ,  $0 < q_1 < 1$ , and  $c > 0$  depending only on the parameters of  $\mathcal{T}$  such that for any  $\theta \in \Theta_m^*$  ( $m \geq 0$ ),*

$$|f_\theta(x) - f_\theta(y)| \leq c \frac{d_{\mathcal{T}}(x, y)^\varepsilon}{|\theta|^{1/2+\varepsilon}} \left( q_1^{\rho_m(\theta, \theta_x^m)} + q_1^{\rho_m(\theta, \theta_y^m)} \right), \quad x, y \in E. \tag{3.10}$$

Moreover,

$$|f_\theta(x)| \leq c |\theta|^{-1/2} q_1^{\rho_m(\theta, \theta_x^m)}, \quad x \in E, \tag{3.11}$$

and for any  $s > 0$  there exist a constant  $c_s$  such that

$$|f_\theta(x)| \leq c_s |\theta|^{-1/2} (\mathcal{M}_{\mathcal{T}}^s \mathbb{1}_\theta)(x), \quad x \in E, \tag{3.12}$$

where  $\mathcal{M}_{\mathcal{T}}^s$  is the maximal operator defined in (2.18).

In addition,

$$c_p^{-1} |\theta|^{1/p-1/2} \leq \|f_\theta\|_{L_p(\theta)} \leq \|f_\theta\|_p \leq c_p |\theta|^{1/p-1/2}, \quad 0 < p \leq \infty, \tag{3.13}$$

and

$$\|f_\theta\|_{H_1} \leq c |\theta|^{1/2}. \tag{3.14}$$

We first prove that each Courant element  $\varphi_\theta$  is an  $\text{Lip } \varepsilon$  function with respect to the quasi-distance  $d_{\mathcal{T}}$ .

**Lemma 3.9** *There exist constants  $\varepsilon > 0$  and  $c > 0$  depending only on the parameters of  $\mathcal{T}$  such that for any  $\theta \in \Theta$*

$$|\varphi_\theta(x) - \varphi_\theta(y)| \leq c \frac{d_{\mathcal{T}}(x, y)^\varepsilon}{|\theta|^\varepsilon}, \quad x, y \in E. \tag{3.15}$$

*Proof.* Let  $\theta \in \Theta_n$ ,  $n \geq 0$ , and assume that  $d_{\mathcal{T}}(x, y) > 0$  and  $x \in \theta$  or  $y \in \theta$ . (Otherwise the claim is trivial.) Also, let  $d_{\mathcal{T}}(x, y) = |\theta^\circ|$  with  $\theta^\circ \in \Theta_m$  ( $m \geq 0$ ) and  $\theta^\circ$  containing  $x$  and  $y$ .

If  $m \leq n$ , (3.15) is trivial because there is a constant  $c > 0$  such that  $|\theta^\circ|/|\theta| > c$  ( $\theta \cap \theta^\circ \neq \emptyset$ ).

Assume that  $n + 3N_0k < m \leq n + 3N_0(k + 1)$  for some  $k \geq 0$ .

Case 1: Let  $x, y \in \Delta$ , where  $\Delta \in \mathcal{T}_n$  is one of the triangles forming  $\theta$ . It is easy to see that estimate (3.15) is invariant under affine transforms. So, without loss of generality we may assume that  $\Delta$  is an equilateral triangle with side lengths 1. Then there exist two triangles  $\Delta', \Delta'' \in \mathcal{T}_m$  with a common vertex such that  $x \in \Delta', y \in \Delta''$  and  $\Delta', \Delta'' \subset \Delta$ . (It may happen that  $\Delta' = \Delta''$ .) By Lemma 2.3,

$$\max\{|\max \text{edge}(\Delta')|, |\max \text{edge}(\Delta'')|\} \leq \rho_1^k |\max \text{edge}(\Delta)| \leq \rho_1^k$$

with  $0 < \rho_1 < 1$ . Choose  $\varepsilon_1 > 0$  so that  $\rho_1 = r^{\varepsilon_1}$ , where  $0 < r < 1$  is from (2.1). Then

$$\begin{aligned} |\phi_\theta(x) - \phi_\theta(y)| &\leq \|\nabla \varphi_\theta\|_{L_\infty(\Delta)} |x - y| \leq 2 \max\{|\max \text{edge}(\Delta')|, |\max \text{edge}(\Delta'')|\} \\ &\leq 2\rho_1^k \leq 2(r^{3N_0k})^{\frac{\varepsilon_1}{3N_0}} \leq 2(r^{3N_0k} |\Delta|)^{\frac{\varepsilon_1}{3N_0}} \leq c |\Delta'|^\varepsilon \leq c |\theta^\circ|^\varepsilon, \end{aligned}$$

where  $\varepsilon = \varepsilon_1/3N_0$  and we used (2.1). Since  $|\theta| \approx |\Delta| \approx 1$ , estimate (3.15) follows.

Case 2: Let  $x \in \Delta_1$  and  $y \in \Delta_2$ , where  $\Delta_1, \Delta_2 \in \mathcal{T}_n$  are distinct triangles with a common vertex  $v_\theta$  and  $\Delta_1, \Delta_2 \subset \theta$ . Since  $x, y \in \theta^\circ$ , there are two smaller subtriangles  $\Delta', \Delta'' \in \mathcal{T}_m$  of  $\theta^\circ$  containing  $x$  and  $y$ , respectively, such that  $\Delta' \subset \Delta_1$  and  $\Delta'' \subset \Delta_2$ . Choose  $z \in \Delta' \cap \Delta''$ . Then using that estimate (3.15) holds in Case 1, we obtain

$$\begin{aligned} |\varphi_\theta(x) - \varphi_\theta(y)| &\leq |\varphi_\theta(x) - \varphi_\theta(z)| + |\varphi_\theta(z) - \varphi_\theta(y)| \\ &\leq c \left( \frac{d_{\mathcal{T}}(x, z)^\varepsilon}{|\theta|^\varepsilon} + \frac{d_{\mathcal{T}}(z, y)^\varepsilon}{|\theta|^\varepsilon} \right) \leq c \frac{d_{\mathcal{T}}(x, y)^\varepsilon}{|\theta|^\varepsilon}, \end{aligned}$$

where we used that  $d_{\mathcal{T}}(x, z) \leq d_{\mathcal{T}}(x, y)$  since  $x, z \in \theta^\circ$ , and similarly  $d_{\mathcal{T}}(z, y) \leq d_{\mathcal{T}}(x, y)$ .

Case 3: Let  $x \in \theta$  and  $y \notin \theta$  (or  $y \in \theta$  and  $x \notin \theta$ ). This case reduces to the first case by introducing an appropriate point  $z$  on the boundary of  $\theta$  such that  $d_{\mathcal{T}}(x, z) \leq d_{\mathcal{T}}(x, y)$  and taking into account that  $\varphi_\theta(y) = \varphi_\theta(z) = 0$ . □

**Proof of Theorem 3.8.** By Lemmas 3.7 and 3.9, we have for  $\theta \in \Theta_m^*$ ,

$$\begin{aligned} |f_\theta(x) - f_\theta(y)| &\leq \sum_{\eta \in \Theta_m} |\eta|^{-1/2} |c_{\theta\eta}| |\varphi_\eta(x) - \varphi_\eta(y)| \\ &\leq \sum_{\eta \in \Theta_m: x \in \eta^\circ \text{ or } y \in \eta^\circ} |\eta|^{-1/2-\varepsilon} q^{\rho_m(\theta, \eta)} d_{\mathcal{T}}(x, y)^\varepsilon \\ &\leq |\theta|^{-1/2-\varepsilon} d_{\mathcal{T}}(x, y)^\varepsilon \left( \sum_{\eta \in \Theta_m: x \in \eta^\circ} + \sum_{\eta \in \Theta_m: y \in \eta^\circ} \right) (|\theta|/|\eta|)^{1/2+\varepsilon} q^{\rho_m(\theta, \eta)}. \end{aligned}$$

Note first that for any  $x \in E$  there are at most 3 cells  $\eta \in \Theta_m$  such that  $x \in \eta^\circ$ . By the definition of  $\theta_x^m$ , we have  $x \in \theta_x^m$  and  $\theta_x^m \in \Theta_m$ , and hence  $|\eta| \approx |\theta_x^m|$  if  $\eta \in \Theta_m$  and  $x \in \eta$ . Also, by (2.3) it follows that  $|\theta|/|\theta_x^m| \leq c(\rho_m(\theta, \theta_x^m) + 1)^s$ . Finally, we choose the constants  $c_1 > 0$  and  $q < q_1 < 1$ , depending only on  $q, \varepsilon$ , and  $s$ , so that  $(\nu + 1)^{(1/2+\varepsilon)s} q^\nu \leq c_1 q_1^\nu$  for  $\nu \geq 1$ . We use this preparation to obtain

$$\sum_{\eta \in \Theta_m: x \in \eta^\circ} (|\theta|/|\eta|)^{1/2+\varepsilon} q^{\rho_m(\theta, \eta)} \leq c(\rho_m(\theta, \theta_x^m) + 1)^{(1/2+\varepsilon)s} q^{\rho_m(\theta, \theta_x^m)} \leq c q_1^{\rho_m(\theta, \theta_x^m)}.$$

We similarly estimate the second sum above, and (3.10) follows.

Estimate (3.11) follows in a similar way but it is easier and will be omitted.

To prove (3.12) we need estimate  $(\mathcal{M}_{\mathcal{T}}^s \mathbb{1}_\theta)(x)$  ( $\theta \in \Theta_m$ ) from below. By (2.2) it follows that

$$(\mathcal{M}_{\mathcal{T}}^s \mathbb{1}_\theta)(x) \geq c > 0 \quad \text{for } x \in \text{Star}_m^1(\theta). \tag{3.16}$$

Suppose  $x \in E \setminus \text{Star}_m^1(\theta)$ . Then  $\rho_m(\theta, \theta_x^m) \geq 2$ . Let  $d_{\mathcal{T}}(v_\theta, x) = |\theta^\circ|$  with  $\theta^\circ \in \Theta_\ell$  and  $\theta^\circ$  containing  $v_\theta$  and  $x$ . Evidently,  $\ell < m$  and  $|\theta \cap \theta^\circ| \geq c|\theta|$ . Further, by Lemma 2.6,  $d_{\mathcal{T}}(v_\theta, x) \leq c|\theta| \rho_m(\theta, \theta_x^m)^\beta$ . Using all of the

above, we obtain

$$(\mathcal{M}_T^s \mathbb{1}_\theta)(x) \geq \left( \frac{1}{|\theta^\circ|} \int_{\theta^\circ} \mathbb{1}_\theta(y) dy \right)^{1/s} = (|\theta \cap \theta^\circ|/|\theta^\circ|)^{1/s} \geq c \rho_m(\theta, \theta_x)^{-\beta/s}.$$

Combining this estimate with (3.16) and (3.11) yields (3.12).

We next prove (3.13). If  $p < \infty$  and  $\theta \in \Theta_m^*$ , by using (3.11) we have

$$\|f_\theta\|_p^p \leq \sum_{\nu=0}^\infty \sum_{\eta \in \mathcal{X}_m^\nu} \|f_\theta\|_{L_p(\eta)}^p \leq |\theta|^{-p/2} \sum_{\nu=0}^\infty \sum_{\eta \in \mathcal{X}_m^\nu} |\eta| q_1^{\rho_m(\theta, \eta)}, \tag{3.17}$$

where  $\mathcal{X}_m^\nu := \{\eta \in \theta_m : \rho_m(\theta, \eta) = \nu\}$ . As above  $|\eta|/|\theta| \leq c(\rho_m(\theta, \eta) + 1)^s = c(\nu + 1)^s$  if  $\eta \in \mathcal{X}_m^\nu$ , and by Lemma 2.4 it follows that  $\#\mathcal{X}_m^\nu \leq c(\nu + 1)^t$ . Using these in (3.17), we find

$$\|f_\theta\|_p^p \leq c |\theta|^{-p/2+1} \sum_{\nu=0}^\infty (\nu + 1)^{s+t} q_1^\nu \leq c |\theta|^{p(1/p-1/2)}. \tag{3.18}$$

We now estimate  $\|f_\theta\|_p$  from below. From the proof of Lemma 3.7 it follows that  $f_\theta = \pm \|g_\theta\|_2^{-1} g_\theta$ , where  $g_\theta = \sum_{\xi \in \Lambda_\theta} b_{\theta\xi} \tilde{\varphi}_\xi$ , and also  $\langle g_\theta, \tilde{\varphi}_\theta \rangle = 1$ . Exactly as in (3.17) and (3.18) (with  $f_\theta$  replaced by  $g_\theta$ ) we obtain  $\|g_\theta\|_2 \leq c$ . Therefore,

$$|\langle f_\theta, \tilde{\varphi}_\theta \rangle| = \|g_\theta\|_2^{-1} |\langle g_\theta, \tilde{\varphi}_\theta \rangle| = \|g_\theta\|_2^{-1} \geq c.$$

On the other hand  $|\langle f_\theta, \tilde{\varphi}_\theta \rangle| \leq |\theta|^{1/2} \|f_\theta\|_{L_\infty(\theta)}$ . Consequently,  $\|f_\theta\|_{L_\infty(\theta)} \approx |\theta|^{-1/2}$  and using (2.12), we infer  $\|f_\theta\|_{L_p(\theta)} \approx |\theta|^{1/p-1/2}$ ,  $0 < p \leq \infty$ .

It remains to prove that each Franklin function  $f_\theta$  belongs to  $H_1$  and (3.14) holds true. To this end it suffices to prove that  $g_\theta := |\theta|^{-1/2} f_\theta$  is a constant multiple of an  $\varepsilon$ -molecule centered at  $v_\theta$ . Evidently  $\int_E |g_\theta(x)|^2 dx = |\theta|^{-1}$ . We use Lemma 2.6 and (3.11) to obtain

$$\begin{aligned} \int_E |g_\theta(x)|^2 d_T(x, v_\theta)^{1+\varepsilon} dx &\leq c |\theta|^{-1+\varepsilon} \int_E q_1^{\rho_m(\theta, \theta_x)} (\rho_m(\theta, \theta_x) + 1)^{\beta(1+\varepsilon)} dx \\ &\leq c |\theta|^{-1+\varepsilon} \sum_{\eta \in \Theta_m} |\eta| q_1^{\rho_m(\theta, \eta)} (\rho_m(\theta, \eta) + 1)^{\beta(1+\varepsilon)} \\ &\leq c |\theta|^\varepsilon, \end{aligned} \tag{3.19}$$

where for the latter estimate we proceed exactly as in (3.17) and (3.18). Therefore, according to the definition of a molecule (see Definition 2.12)  $g_\theta$  is a constant multiple of an  $\varepsilon$ -molecule and hence  $f_\theta \in H_1$  and  $\|f_\theta\|_{H_1} \leq c |\theta|^{1/2}$ . □

### 3.4 Proof of Theorem 3.2

We first observe that since by Theorem 3.8 each Franklin function  $f_\theta$  belongs to  $H_1$ , then  $(f_\theta, f_\theta)_{\theta \in \Theta^*}$  is a biorthogonal system in  $H_1$ .

We next prove a technical result which will provide the main step in the proof of Theorem 3.2.

**Proposition 3.10** *For any 2-atom  $a(x)$ ,*

$$a = \sum_{\theta \in \Theta^*} \langle a, f_\theta \rangle f_\theta, \tag{3.20}$$

where the series converges unconditionally in  $H_1$ . Moreover, there exists a constant  $c > 0$  depending only on the parameters of  $\mathcal{T}$  such that for any  $\mathcal{M} \subset \Theta^*$  and any sequence  $\omega = (\omega_\theta)_{\theta \in \mathcal{M}}$  with  $\omega_\theta = \pm 1$ ,

$$\left\| \sum_{\theta \in \mathcal{M}} \omega_\theta \langle a, f_\theta \rangle f_\theta \right\|_{H_1} \leq c \tag{3.21}$$

and

$$\left\| \left( \sum_{\theta \in \Theta^*} |\langle a, f_\theta \rangle|^2 |f_\theta(x)|^2 \right)^{1/2} \right\|_{L_1} \leq c. \tag{3.22}$$

For given  $\mathcal{M} \subset \Theta^*$  and  $\omega = (\omega_\theta)_{\theta \in \mathcal{M}}$  with  $\omega_\theta = \pm 1$ , consider the linear operator

$$(T_{\mathcal{M}, \omega} f)(x) := \int_E K_{\mathcal{M}, \omega}(x, y) f(y) dy \tag{3.23}$$

with kernel

$$K_{\mathcal{M}, \omega}(x, y) := \sum_{\theta \in \mathcal{M}} \omega_\theta f_\theta(x) f_\theta(y). \tag{3.24}$$

With the next lemma we show that the kernel  $K_{\mathcal{M}, \omega}$  satisfies condition (2.28) of Proposition 2.13

**Lemma 3.11** *There exists a constant  $c > 0$  such that if  $d_{\mathcal{T}}(x, y_0) > 2Kd_{\mathcal{T}}(y, y_0)$ , where  $K$  is the constant from the definition of the quasi-distance  $d_{\mathcal{T}}(\cdot, \cdot)$ , and  $\min\{l(\theta) : \theta \in \mathcal{M}\} \geq \mu \geq 0$ , then*

$$|K_{\mathcal{M}, \omega}(x, y) - K_{\mathcal{M}, \omega}(x, y_0)| \leq \gamma_\mu \left( \frac{d_{\mathcal{T}}(y, y_0)}{d_{\mathcal{T}}(x, y_0)} \right)^\varepsilon \frac{1}{d_{\mathcal{T}}(x, y_0)}, \tag{3.25}$$

where  $\varepsilon > 0$  is from Theorem 3.8,  $0 < \gamma_\mu \leq c$ , and  $\gamma_\mu \rightarrow 0$  as  $\mu \rightarrow \infty$ .

*Proof.* Denote  $K_m(x, y) := \sum_{\theta \in \mathcal{M} \cap \Theta_m^*} \omega_\theta f_\theta(x) f_\theta(y)$ . We use (3.10) and (3.11) to obtain

$$\begin{aligned} |K_m(x, y) - K_m(x, y_0)| &\leq \sum_{\theta \in \Theta_m} |f_\theta(x)| |f_\theta(y) - f_\theta(y_0)| \\ &\leq c d_{\mathcal{T}}(y, y_0)^\varepsilon \sum_{\theta \in \Theta_m} |\theta|^{-1-\varepsilon} q_1^{\rho_m(\theta, \theta_x^m)} \left( q_1^{\rho_m(\theta, \theta_y^m)} + q_1^{\rho_m(\theta, \theta_{y_0}^m)} \right). \end{aligned} \tag{3.26}$$

We now claim that there exist constants  $c > 0$  and  $0 < q_2 < 1$  such that

$$\sum_{\theta \in \Theta_m} |\theta|^{-1-\varepsilon} q_1^{\rho_m(\theta, \theta_x^m) + \rho_m(\theta, \theta_y^m)} \leq c |\theta_y^m|^{-1-\varepsilon} q_2^{\rho_m(\theta_x^m, \theta_y^m)}. \tag{3.27}$$

To see this set  $n := \rho_m(\theta_x^m, \theta_y^m)$  and define

$$A_k := \{ \theta \in \Theta_m : \rho_m(\theta, \theta_x^m) + \rho_m(\theta, \theta_y^m) = k + n \}.$$

By Lemma 2.4,  $\#A_k \leq c(n+k)^t$ , and by Lemma 2.1,  $|\theta_y^m|/|\theta| \leq c\rho_m(\theta, \theta_y^m)^s$  if  $\theta \in \Theta_m$ . We use the above to obtain

$$\begin{aligned} \sum_{\theta \in \Theta_m} \frac{q_1^{\rho_m(\theta, \theta_x^m) + \rho_m(\theta, \theta_y^m)}}{|\theta|^{1+\varepsilon}} &= |\theta_y^m|^{-1-\varepsilon} \sum_{k \geq 0} \sum_{\theta \in A_k} (|\theta_y^m|/|\theta|)^{1+\varepsilon} q_1^{n+k} \\ &\leq c |\theta_y^m|^{-1-\varepsilon} \sum_{k \geq 0} \sum_{\theta \in A_k} \rho_m(\theta, \theta_y^m)^{s(1+\varepsilon)} q_1^{n+k} \\ &\leq c |\theta_y^m|^{-1-\varepsilon} \sum_{k \geq 0} (n+k)^{t+s(1+\varepsilon)} q_1^{n+k} \\ &\leq c |\theta_y^m|^{-1-\varepsilon} q_2^n \\ &= c |\theta_y^m|^{-1-\varepsilon} q_2^{\rho_m(\theta_x^m, \theta_y^m)}, \end{aligned}$$

where  $q_1 < q_2 < 1$ . Thus (3.27) is established.

Applying (3.27) in (3.26) we get

$$|K_m(x, y) - K_m(x, y_0)| \leq c d_{\mathcal{T}}(y, y_0)^\varepsilon \left( \frac{q_2^{\rho_m(\theta_x^m, \theta_y^m)}}{|\theta_y^m|^{1+\varepsilon}} + \frac{q_2^{\rho_m(\theta_x^m, \theta_{y_0}^m)}}{|\theta_{y_0}^m|^{1+\varepsilon}} \right)$$



and hence

$$\begin{aligned}
 |K_{\mathcal{M},\omega}(x, y) - K_{\mathcal{M},\omega}(x, y_0)| &\leq \sum_{m=0}^{\infty} |K_m(x, y) - K_m(x, y_0)| \\
 &\leq c d_{\mathcal{T}}(y, y_0)^{\varepsilon} \sum_{m=\mu}^{\infty} \left( \frac{q_2^{\rho_m(\theta_x^m, \theta_y^m)}}{|\theta_y^m|^{1+\varepsilon}} + \frac{q_2^{\rho_m(\theta_x^m, \theta_{y_0}^m)}}{|\theta_{y_0}^m|^{1+\varepsilon}} \right).
 \end{aligned}
 \tag{3.28}$$

To estimate the latter quantity above let us assume that  $d_{\mathcal{T}}(x, y) = |\theta^*|$  with  $\theta^* \in \Theta_{m^*}$ . Then

$$\sum_{m=\mu}^{\infty} \frac{q_2^{\rho_m(\theta_x^m, \theta_y^m)}}{|\theta_y^m|^{1+\varepsilon}} = \left( \sum_{m=\mu}^{m^*} + \sum_{m=m^*+1}^{\infty} \right) \frac{q_2^{\rho_m(\theta_x^m, \theta_y^m)}}{|\theta_y^m|^{1+\varepsilon}} =: \sigma_1 + \sigma_2.
 \tag{3.29}$$

For  $\sigma_1$ , using (2.1) and (2.2), we have

$$\begin{aligned}
 \sigma_1 &\leq \frac{c}{|\theta^*|^{1+\varepsilon}} \sum_{m=\mu}^{m^*} (|\theta^*|/|\theta_y^m|)^{1+\varepsilon} q_2^{\rho_m(\theta_x^m, \theta_y^m)} \\
 &\leq \frac{c}{|\theta^*|^{1+\varepsilon}} \sum_{m=\mu}^{m^*} \rho^{(m^*-m)(1+\varepsilon)} \\
 &\leq \frac{c}{d_{\mathcal{T}}(x, y)^{1+\varepsilon}}.
 \end{aligned}
 \tag{3.30}$$

We now estimate  $\sigma_2$ . Denote  $\mathcal{I}_k := [m^* + 2N_0k, m^* + 2N_0(k + 1))$ . Note that by (2.9)

$$\rho_{\ell}(\theta_x^{\ell}, \theta_y^{\ell}) \geq 2^{k-2} \rho_{m^*+2N_0}(\theta_x^{m^*+2N_0}, \theta_y^{m^*+2N_0}) \geq 2^{k-2},$$

whenever  $\ell \geq m^* + 2N_0k$ . We use this, (2.1) and (2.2) to obtain

$$\begin{aligned}
 \sigma_2 &\leq \frac{c}{|\theta^*|^{1+\varepsilon}} \sum_{m=m^*+1}^{\infty} (|\theta^*|/|\theta_y^m|)^{1+\varepsilon} q_2^{\rho_m(\theta_x^m, \theta_y^m)} \\
 &\leq \frac{c}{d_{\mathcal{T}}(x, y)^{1+\varepsilon}} \sum_{m=m^*+1}^{\infty} r^{(m^*-m)(1+\varepsilon)} q_2^{\rho_m(\theta_x^m, \theta_y^m)} \\
 &\leq \frac{c}{d_{\mathcal{T}}(x, y)^{1+\varepsilon}} \sum_{k=0}^{\infty} \sum_{m \in \mathcal{I}_k} r^{(m^*-m)(1+\varepsilon)} q_2^{2^k} \\
 &\leq \frac{c}{d_{\mathcal{T}}(x, y)^{1+\varepsilon}} \sum_{k=0}^{\infty} r^{-2N_0(k+1)(1+\varepsilon)} q_2^{2^k} \\
 &\leq \frac{c}{d_{\mathcal{T}}(x, y)^{1+\varepsilon}}.
 \end{aligned}
 \tag{3.31}$$

Estimates (3.30) and (3.31) yield

$$\sum_{m=0}^{\infty} \frac{q^{\rho_m(\theta_x^m, \theta_y^m)}}{|\theta_y^m|^{1+\varepsilon}} \leq \frac{c}{d_{\mathcal{T}}(x, y)^{1+\varepsilon}}.$$

Using this inequality twice in (3.28) (for  $y$  and  $y_0$ ) we obtain

$$|K_{\mathcal{M},\omega}(x, y) - K_{\mathcal{M},\omega}(x, y_0)| \leq c d_{\mathcal{T}}(y, y_0)^{\varepsilon} \left( \frac{1}{d_{\mathcal{T}}(x, y)^{1+\varepsilon}} + \frac{1}{d_{\mathcal{T}}(x, y_0)^{1+\varepsilon}} \right).
 \tag{3.32}$$

Finally, since  $d_{\mathcal{T}}(x, y_0) > 2K d_{\mathcal{T}}(y, y_0)$ , it follows that  $d(x, y) \approx d(x, y_0)$  and (3.32) implies (3.25) with  $\gamma_{\mu} = c$ .

If  $m^* + 2N_0k^{\natural} \leq \mu < m^* + 2N_0(k^{\natural} + 1)$  for some  $k^{\natural} \geq 1$ , then  $\sigma_1 = 0$  and as above

$$\sigma_2 \leq \frac{c}{d_{\mathcal{T}}(x, y)^{1+\varepsilon}} \sum_{k=k^{\natural}}^{\infty} r^{-2N_0(k+1)(1+\varepsilon)} q_2^{2k} \leq \frac{cq_3^{2(\mu-m^*)/2N_0}}{d_{\mathcal{T}}(x, y)^{1+\varepsilon}}$$

for some  $0 < q_3 < 1$ . This estimate yields (3.25) with  $\gamma_{\mu} := cq_3^{2(\mu-m^*)/2N_0}$ . Evidently,  $\gamma_{\mu} \rightarrow 0$  as  $\mu \rightarrow \infty$ . The proof of Lemma 3.11 is complete. □

**Proof of Proposition 3.10.** Suppose first that  $a(x) = \pm|E|^{-1}\mathbb{1}_E$ . If  $\theta_0 \in \mathcal{M}$ , then using the orthogonality of the Franklin functions and that  $f_{\theta_0} = \pm|E|^{-1/2}\mathbb{1}_E$  it follows that  $T_{\mathcal{M},\omega}a = \pm a$ , while if  $\theta_0 \notin \mathcal{M}$  then  $T_{\mathcal{M},\omega}a = 0$ . In both cases the result obviously holds.

Let now  $a(x)$  be a 2-atom and  $a \neq \pm|E|^{-1}\mathbb{1}_E$ . Then  $\int_E a = 0$ . As above, we see that  $\int_E T_{\mathcal{M},\omega}a = 0$ . Then by Proposition 2.13 and Lemma 3.11 with  $\mathcal{M} := \{\theta \in \Theta^* : l(\theta) \geq \mu\}$  and an arbitrary sequence  $\omega = (\omega_{\theta})_{\theta \in \Theta^*}$  with  $\omega_{\theta} = \pm 1$ , it follows that the function

$$h_{\mu} := \sum_{\theta \in \Theta^*, l(\theta) \geq \mu} \omega_{\theta} \langle a, f_{\theta} \rangle f_{\theta}$$

belongs to  $H_1$  and  $\|h_{\mu}\|_{H_1} \leq \gamma_{\mu}$ , where  $\gamma_{\mu} \rightarrow 0$  as  $\mu \rightarrow \infty$ . Therefore, the series in (3.20) converges unconditionally in  $H_1$ . We also know that (3.20) holds in  $L_2$ . Since both spaces  $H_1$  and  $L_2$  are continuously embedded in  $L_1$ , it follows that (3.20) holds in  $H_1$ .

Estimate (3.21) follows by Proposition 2.13 and Lemma 3.11.

It remains to prove (3.22). From above for an arbitrary sequence  $\omega = (\omega_{\theta})_{\theta \in \Theta^*}$  with  $\omega_{\theta} = \pm 1$ , we have

$$\left\| \sum_{\theta \in \Theta^*} \omega_{\theta} \langle a, f_{\theta} \rangle f_{\theta} \right\|_{L_1} \leq c \left\| \sum_{\theta \in \Theta^*} \omega_{\theta} \langle a, f_{\theta} \rangle f_{\theta} \right\|_{H_1} \leq c.$$

Applying now the usual trick with Khintchine’s inequality, we infer

$$\left\| \sum_{\theta \in \Theta^*} |\langle a, f_{\theta} \rangle|^2 |f_{\theta}(x)|^2 \right\|_{L_1} \leq c \left\| \sum_{\theta \in \Theta^*} \omega_{\theta} \langle a, f_{\theta} \rangle f_{\theta} \right\|_{L_1} \leq c. \quad \square$$

**Proof of Theorem 3.2.** We begin by proving that the Franklin system  $\mathcal{F}_{\mathcal{T}}$  is a unconditional basis for  $H_1 := H_1(E, \mathcal{T})$ . As we have already mentioned,  $f_{\theta} \in H_1$  for all  $\theta \in \Theta^*$ , and hence  $(f_{\theta}, f_{\theta})_{\theta \in \Theta^*}$  is a biorthogonal system in  $H_1$ .

By Proposition 3.10 for any atom  $a(x)$ ,  $a = \sum_{\theta \in \Theta^*} \langle a, f_{\theta} \rangle f_{\theta}$  in  $H_1$ . This along with the definition of  $H_1$  yields that  $\mathcal{F}_{\mathcal{T}}$  is dense in  $H_1$ .

By estimate (3.21) from Proposition 3.10, it readily follows that for any  $\mathcal{M} \subset \Theta^*$ ,

$$\left\| \sum_{\theta \in \mathcal{M}} \langle f, f_{\theta} \rangle f_{\theta} \right\|_{H_1} \leq c \|f\|_{H_1} \quad \text{for } f \in H_1. \tag{3.33}$$

Therefore,  $\mathcal{F}_{\mathcal{T}}$  is a unconditional basis for  $H_1$ .

We now turn to  $L_p$  ( $1 < p < \infty$ ). Taking into account Lemma 2.3, it is obvious that  $\mathcal{F}_{\mathcal{T}}$  is dense in  $L_p(E)$ . For a given  $\mathcal{M} \subset \Theta^*$ , consider the operator  $T_{\mathcal{M}} := T_{\mathcal{M},\omega}$  with  $\omega = (1)_{\theta \in \mathcal{M}}$ , where  $T_{\mathcal{M},\omega}$  is defined in (3.23). By (3.33)  $T_{\mathcal{M}}$  is bounded on  $H_1$  and, since  $\mathcal{F}_{\mathcal{T}}$  is an orthogonal basis for  $L_2$ ,  $T_{\mathcal{M}}$  is bounded on  $L_2$  as well. Then by interpolation it follows that  $T_{\mathcal{M}}$  is bounded on  $L_p$  for  $1 < p \leq 2$ . Finally, by a standard duality argument, it easily follows that  $T_{\mathcal{M}}$  is bounded on  $L_p$ ,  $2 < p < \infty$ , as well. Consequently,  $\mathcal{F}_{\mathcal{T}}$  is a unconditional basis for  $L_p(E)$ ,  $1 < p < \infty$ . □

### 3.5 Proof of Theorem 3.3

We first note that the implication (a) → (b) is immediate from the fact that  $\mathcal{F}_{\mathcal{T}}$  is a unconditional basis for  $H_1$  (see Theorem 3.2), since  $H_1$  is embedded in  $L_1$ .

One applies Khintchine’s inequality as usual to show that (b) ↔ (c) (see e.g. [11, 17]).

We now show that (c) ↔ (d). We know from Theorem 3.8 that for any  $s > 0$ ,  $|f_{\theta}(x)| \leq c_s(\mathcal{M}_{\mathcal{T}}^s \tilde{\mathbb{I}}_{\theta})(x)$ ,  $x \in E$ . Then choosing  $0 < s < 1$  and applying Proposition 2.9 we obtain  $\|S_f\|_{L_1} \leq c\|F_f\|_{L_1}$ ; thus (d) → (c).

For the other direction, we first note that by Theorem 3.8,  $\|f_{\theta}\|_{L_{\infty}(\theta)} \geq c|\theta|^{-1/2}$ . Then since  $f_{\theta}$  is linear on each triangle of  $\theta$ , there exist a set  $G_{\theta} \subset \theta$  with  $|G_{\theta}| \geq \alpha|\theta|$  such that  $|f_{\theta}(x)| \geq c|G_{\theta}|^{-1/2}$  for  $x \in G_{\theta}$ , where  $c > 0$  and  $0 < \alpha < 1$  are constants depending only on the parameters of  $\mathcal{T}$ . Therefore,  $\tilde{\mathbb{I}}_{G_{\theta}}(x) \leq c|f_{\theta}(x)|$  for  $x \in E$ . Denote

$$\Gamma_f(x) := \left( \sum_{\theta \in \Theta^*} |\langle f, f_{\theta} \rangle|^2 |\tilde{\mathbb{I}}_{G_{\theta}}(x)|^2 \right)^{1/2}.$$

Then by the above estimate  $\|\Gamma_f\|_{L_1} \leq c\|S_f\|_{L_1}$ . On the other hand,  $\tilde{\mathbb{I}}_{\theta}(x) \leq c_s(\mathcal{M}_{\mathcal{T}}^s \tilde{\mathbb{I}}_{G_{\theta}})(x)$ , since  $|G_{\theta}| \geq \alpha|\theta|$ . Applying again Proposition 2.9 we infer  $\|F_f\|_{L_1} \leq c\|\Gamma_f\|_{L_1} \leq c\|S_f\|_{L_1}$ . Consequently, (c) → (d).

It remains to show that (d) → (a) which is the main step in the proof of Theorem 3.3. We give it in the following proposition.

**Proposition 3.12** *Suppose that for a collection of numbers  $(a_{\theta})_{\theta \in \Theta^*}$  we have*

$$F(x) := \left( \sum_{\theta \in \Theta^*} |a_{\theta}|^2 |\tilde{\mathbb{I}}_{\theta}(x)|^2 \right)^{1/2} \in L_1 \tag{3.34}$$

which is equivalent to

$$S(x) := \left( \sum_{\theta \in \Theta^*} |a_{\theta}|^2 |f_{\theta}(x)|^2 \right)^{1/2} \in L_1. \tag{3.35}$$

Then  $f := \sum_{\theta \in \Theta^*} a_{\theta} f_{\theta}$  belongs to  $H_1$  and  $\|f\|_{H_1} \leq c\|F\|_{L_1}$ .

*Proof.* We shall use the idea of the proof in the wavelet case (see e.g. [14, 17]). Note first that the equivalence  $\|F\|_{L_1} \approx \|S\|_{L_1}$  follows by the argument which we used above to show that (c) ↔ (d).

Let us denote

$$G_k := \{x \in E : F(x) > 2^k\}, \quad k \in \mathbb{Z}.$$

It is easy to see that

$$\sum_{k \in \mathbb{Z}} 2^k |G_k| \leq 2\|F\|_{L_1(E)} \tag{3.36}$$

(see e.g. [17, Proposition 8.15]).

We introduce the collections of cells

$$\mathcal{C}_k := \{\theta \in \Theta^* : |\theta \cap G_k| > |\theta|/2\}.$$

Since  $G_{k+1} \subset G_k$ , then  $\mathcal{C}_{k+1} \subset \mathcal{C}_k$ . It is easy to see that

$$\Theta^* = \bigcup_{j \in \mathbb{Z}} \mathcal{C}_j.$$

Indeed, if  $a_{\theta} \neq 0$ , then  $|a_{\theta}||\theta|^{-1/2} > 2^j$  for some  $j \in \mathbb{Z}$ , and hence  $F(x) \geq |a_{\theta}|\tilde{\mathbb{I}}(x) > 2^j$  on  $\theta$ . Therefore,  $\theta \in \mathcal{C}_j$ .

Denote now

$$G_k^* := \bigcup_{\theta \in \mathcal{C}_k} \theta.$$

It is not hard to see that the Lebesgue differentiation theorem holds with the cells from  $\Theta$ . (For its proof one can use the maximal operator  $\mathcal{M}_{\mathcal{T}}^1$  introduced in §2.2.) Consequently,  $G_k \subset G_k^*$  modulo a set of measure zero.

By the coloring lemma in [10] (Lemma 3.2),  $\Theta$  can be represented as a finite disjoint union of subsets  $(\Theta^\nu)_{\nu=1}^K$  with  $K = K(N_0, M_0)$  such that each  $\Theta^\nu$  has a tree structure with respect to the inclusion relation, i.e., if  $\theta', \theta'' \in \Theta^\nu$ , then  $(\theta')^\circ \cap (\theta'')^\circ = \emptyset$  or  $\theta' \subset \theta''$  or  $\theta'' \subset \theta'$ .

Further, denote by  $\mathcal{M}_{k\nu}$  the set of all maximal cells in  $\mathcal{C}_k \cap \Theta^\nu$ , i.e.,  $\theta \in \mathcal{M}_{k\nu}$  if  $\theta \in \mathcal{C}_k \cap \Theta^\nu$  and  $\theta$  is not contained in any other cell from  $\mathcal{C}_k \cap \Theta^\nu$ . Clearly,  $G_k^* = \bigcup_{\nu=1}^K \bigcup_{\theta \in \mathcal{M}_{k\nu}} \theta$  and

$$|G_k^*| \leq \sum_{\nu=1}^K \sum_{\theta \in \mathcal{M}_{k\nu}} |\theta| \leq 2 \sum_{\nu=1}^K \sum_{\theta \in \mathcal{M}_{k\nu}} |G_k \cap \theta| \leq 2 \sum_{\nu=1}^K \left| G_k \cap \bigcup_{\theta \in \mathcal{M}_{k\nu}} \theta \right| \leq 2K |G_k|. \quad (3.37)$$

Denote  $\mathcal{D}_k := \mathcal{C}_k \setminus \mathcal{C}_{k+1}$ . Now for any  $\theta \in \mathcal{M}_{k\nu}$  we define  $\mathcal{D}_\theta := \{\eta \in \mathcal{D}_k \cap \Theta^\nu : \eta \subset \theta\}$ . Clearly, the sets  $\mathcal{D}_\theta$  are disjoint and

$$\Theta^* = \bigcup_{k \in \mathbb{Z}} \bigcup_{\nu=1}^K \bigcup_{\theta \in \mathcal{M}_{k\nu}} \mathcal{D}_\theta.$$

We also define

$$\mathcal{A}_\theta := \sum_{\eta \in \mathcal{D}_\theta} a_\eta f_\theta, \quad \theta \in \mathcal{M}_{k\nu}.$$

We next show that  $m_\theta := 2^{-k} |\theta|^{-1} \mathcal{A}_\theta$  is a constant multiple of an  $\varepsilon$ -molecule, which will imply  $\|m_\theta\|_{H_1} \leq c$  and as a consequence  $\|\mathcal{A}_\theta\|_{H_1} \leq 2^k |\theta|$ . Note first that, for  $\theta \in \mathcal{M}_{k\nu}$ ,

$$\int_{\theta \setminus G_{k+1}} F^2(x) dx \geq \sum_{\eta \in \mathcal{D}_\theta} |a_\eta|^2 \int_{\theta \setminus G_{k+1}} \tilde{\mathbb{I}}_\eta^2(x) dx \geq \sum_{\eta \in \mathcal{D}_\theta} |a_\eta|^2 |\eta|^{-1} |\eta \setminus G_{k+1}|.$$

Since  $\eta \notin \mathcal{C}_{k+1}$ ,  $|G_{k+1} \cap \eta| \leq |\eta|/2$  and hence  $|\eta \setminus G_{k+1}| \geq |\eta|/2$ . Therefore,

$$\int_{\theta \setminus G_{k+1}} F^2(x) dx \geq \frac{1}{2} \sum_{\eta \in \mathcal{D}_\theta} |a_\eta|^2, \quad \theta \in \mathcal{M}_{k\nu}. \quad (3.38)$$

On the other hand,

$$\int_{\theta \setminus G_{k+1}} F^2(x) dx \leq 2^{2(k+1)} |\theta \setminus G_{k+1}| \leq 4 \cdot 2^{2k} |\theta|.$$

Combining this with (3.38) we arrive at

$$\|\mathcal{A}_\theta\|_2^2 = \sum_{\eta \in \mathcal{D}_\theta} |a_\eta|^2 \leq c 2^{2k} |\theta|. \quad (3.39)$$

To prove that  $m_\theta$  is a constant multiple of an  $\varepsilon$ -molecule it suffices to show that  $m_\theta$  satisfies (2.26) with the 1 in the right-hand side of (2.26) replaced by a constant  $c > 0$  (for some  $\varepsilon > 0$ ). This is apparently equivalent to

$$\left( \int_E |\mathcal{A}_\theta(x)|^2 dx \right) \left( \int_E |\mathcal{A}_\theta(x)|^2 d_{\mathcal{T}}(x, v_\theta)^{1+\varepsilon} dx \right)^{1/\varepsilon} \leq c (2^{2k} |\theta|^2)^{1+1/\varepsilon}.$$

We choose an arbitrary  $\varepsilon > 0$ , e.g.  $\varepsilon = 1$ , and fix it. Taking into account (3.39) it suffices to show that

$$\int_E |\mathcal{A}_\theta(x)|^2 d_{\mathcal{T}}(x, v_\theta)^{1+\varepsilon} dx \leq c 2^{2k} |\theta|^{2+\varepsilon}. \quad (3.40)$$

To prove this estimate we need the following lemma:

**Lemma 3.13** For  $\theta \in \Theta_m$  ( $m \geq 0$ ) and  $x \in E$  with  $\rho_m(\theta, \theta_x^m) \geq 3$ , we have

$$R(x) := \sum_{\eta \in \Theta^*, \eta \subset \theta} |f_\eta(x)| \leq c |\theta|^{-1/2} q_*^{\rho_m(\theta, \theta_x^m)}, \tag{3.41}$$

where the constants  $c > 0$  and  $0 < q_* < 1$  depend only on the parameters of  $\mathcal{T}$ .

*Proof.* Write briefly  $\mathcal{N} := \rho_m(\theta, \theta_x^m)$  and  $\mathcal{I}_k := [2N_0k, 2N_0(k+1))$ , where  $N_0$  is from condition (d) on LR-triangulations (§2.1). If  $\nu \in \mathcal{I}_k$  ( $k \geq 0$ ), then using (2.9) we obtain  $\rho_{m+\nu}(v_\theta, x) \geq 2^{k-2} \rho_m(v_\theta, x) = 2^{k-2} \mathcal{N}$ . Then by (3.11) we have for  $\nu \in \mathcal{I}_k$ ,

$$\begin{aligned} \sum_{\eta \in \Theta_{m+\nu}^*, \eta \subset \theta} |f_\eta(x)| &\leq c \sum_{\eta \in \Theta_{m+\nu}^*, \eta \subset \theta} |\eta|^{-1/2} q_1^{\rho_{m+\nu}(\eta, \theta_x^{m+\nu})} \\ &\leq c \sum_{\mu \geq 2^{k-2} \mathcal{N}} \sum_{\eta \in \mathcal{X}_{m+\nu}^\mu} |\eta|^{-1/2} q_1^\mu, \end{aligned} \tag{3.42}$$

where

$$\mathcal{X}_{m+\nu}^\mu := \{ \eta \in \Theta_{m+\nu} : \rho_{m+\nu}(\eta, \theta_x^{m+\nu}) = \mu \text{ and } \eta \subset \theta \}.$$

By Lemma 2.4,  $\#\mathcal{X}_{m+\nu}^\mu \leq c\mu^t$ , and by (2.1) and (2.2) we have for  $\eta \in \mathcal{X}_{m+\nu}^\mu$ ,  $|\eta| \geq cr^\nu |\theta| \geq cr^{2N_0k} |\theta|$ , since  $\eta \subset \theta$ . Using these in (3.42), we obtain

$$\begin{aligned} \sum_{\eta \in \Theta_{m+\nu}^*, \eta \subset \theta} |f_\eta(x)| &\leq c |\theta|^{-1/2} r^{-N_0k} \sum_{\mu \geq 2^{k-1} \mathcal{N}} q_1^\mu \#\mathcal{X}_{m+\nu}^\mu \\ &\leq c |\theta|^{-1/2} r^{-N_0k} \sum_{\mu \geq 2^{k-1} \mathcal{N}} q_1^\mu \mu^t \\ &\leq c |\theta|^{-1/2} q_*^{2^k \mathcal{N}}, \end{aligned}$$

where  $0 < q_* < 1$  and we used that  $0 < q_1 < 1$ . Summing up the above estimates we get

$$R(x) \leq c |\theta|^{-1/2} \sum_{k=0}^\infty 2N_0 q_*^{2^k \mathcal{N}} \leq c |\theta|^{-1/2} q_*^{\mathcal{N}} = c |\theta|^{-1/2} q_*^{\rho_m(\theta, \theta_x^m)}.$$

□

We are now prepared to prove (3.40). Assuming that  $\theta \in \Theta_m$ ,  $m \geq 0$ , we can write

$$\int_E |\mathcal{A}_\theta(x)|^2 d_{\mathcal{T}}(x, v_\theta)^{1+\varepsilon} dx = \int_{\text{Star}_m^3(v_\theta)} + \int_{E \setminus \text{Star}_m^3(v_\theta)} =: J_0 + J_1.$$

To estimate the first integral we note that  $d_{\mathcal{T}}(v_\theta, x) \leq c|\theta|$  if  $x \in \text{Star}_m^3(v_\theta)$  and using (3.39) we obtain

$$J_0 \leq c |\theta|^{1+\varepsilon} \|\mathcal{A}_\theta\|_2^2 \leq c 2^{2k} |\theta|^{2+\varepsilon}. \tag{3.43}$$

To estimate  $J_1$  we first observe that by (3.39) it follows that  $|a_\eta| \leq c 2^k |\theta|^{1/2}$  for  $\eta \in \mathcal{D}_\theta$ . Using this and Lemma 2.6, we obtain

$$\begin{aligned} J_1 &= \int_{E \setminus \text{Star}_m^3(v_\theta)} |\mathcal{A}_\theta(x)|^2 d_{\mathcal{T}}(x, v_\theta)^{1+\varepsilon} dx \\ &\leq \int_{E \setminus \text{Star}_m^3(v_\theta)} \left( \sum_{\eta \in \mathcal{D}_\theta} |a_\eta| |f_\eta(x)| \right)^2 d_{\mathcal{T}}(x, v_\theta)^{1+\varepsilon} dx \\ &\leq c 2^{2k} |\theta|^{2+\varepsilon} \int_{E \setminus \text{Star}_m^3(v_\theta)} \left( \sum_{\eta \in \Theta^*, \eta \subset \theta} |f_\eta(x)| \right)^2 \rho_m(\theta, \theta_x^m)^{\beta(1+\varepsilon)} dx. \end{aligned}$$

Applying Lemma 3.13 we obtain

$$J_1 \leq c 2^{2k} |\theta|^{1+\varepsilon} \int_{E \setminus \text{Star}_m^3(v_\theta)} q_*^{\rho_m(\theta, \theta_x^m)} \rho_m(\theta, \theta_x^m)^{\beta(1+\varepsilon)} dx$$

and exactly as in (3.19) (see also (3.17) and (3.18)) we obtain  $J_1 \leq c 2^{2k} |\theta|^{2+\varepsilon}$ . This and (3.43) yield (3.40). Consequently,  $\|\mathcal{A}_\theta\|_{H_1} \leq c 2^k |\theta|$ , which implies

$$\|f\|_{H_1} \leq \sum_{k \in \mathbb{Z}} \sum_{\nu=1}^K \sum_{\theta \in \mathcal{D}_{k\nu}} \|\mathcal{A}_\theta\|_{H_1} \leq c \sum_{k \in \mathbb{Z}} 2^k \sum_{\nu=1}^K \sum_{\theta \in \mathcal{M}_{k\nu}} |\theta| \leq c \sum_{k \in \mathbb{Z}} 2^k |G_k| \leq c \|F\|_{L_1},$$

where we used (3.36). This completes the proof of Proposition 3.12. □

It remains to prove equivalences (3.2). The estimate  $\|f\|_{H_1} \leq c \|S_f\|_{L_1}$  and the equivalence  $\|S_f\|_{H_1} \approx \|F_f\|_{L_1}$  are immediate from Proposition 3.12. The estimate  $\|S_f\|_{L_1} \leq c \|f\|_{H_1}$  holds since by Proposition 3.10 it is true for each individual atom. The proof of Theorem 3.3 is complete. □

### 3.6 Proof of Theorem 3.4

We shall follow the scheme of the proof of Wojtaszczyk [16] combined with our techniques of this article. We begin with one technical lemma.

**Lemma 3.14** *For any  $\theta \in \Theta_m$  ( $m \geq 0$ ) and  $1 \leq q < \infty$ , we have*

$$\left\| \sum_{\eta \in \Theta^*, \eta \subset \theta} |\eta|^{1/2} |f_\eta| \right\|_{L_q(E \setminus \theta)} \leq c |\theta|^{1/q} \tag{3.44}$$

and

$$\left\| \sum_{\eta \in \Theta^*, \eta \not\subset \theta, l(\eta) > m} |\eta|^{1/2} |f_\eta| \right\|_{L_q(\theta)} \leq c |\theta|^{1/q}. \tag{3.45}$$

*Proof.* We shall prove only (3.44) since the proof of (3.45) is the same. We first observe that by (3.11) and Lemma 2.4,

$$\sum_{\eta \in \Theta_n^*} |\eta|^{1/2} |f_\eta(x)| \leq c, \quad x \in E, \quad n \geq 0. \tag{3.46}$$

For  $x \in E \setminus \theta$  and  $\nu \geq 1$ , we define  $\rho_{m+\nu}(\theta, \theta_x^{m+\nu}) := \inf_{y \in \theta} \rho_{m+\nu}(\theta_y^{m+\nu}, \theta_x^{m+\nu})$ . Exactly as in the proof of Lemma 3.13 one shows that there exist constants  $0 < q_* < 1$  and  $c > 0$  such that

$$\sum_{\eta \in \Theta^*, \eta \subset \theta, l(\eta) \geq m+\nu} |\eta|^{1/2} |f_\eta(x)| \leq c q_*^{\rho_{m+\nu}(\theta, \theta_x^{m+\nu})}, \quad x \in E \setminus \text{Star}_{m+\nu}^2(\theta), \tag{3.47}$$

and, in particular,

$$\sum_{\eta \in \Theta^*, \eta \subset \theta} |\eta|^{1/2} |f_\eta(x)| \leq c q_*^{\rho_m(\theta, \theta_x^m)}, \quad x \in E \setminus \text{Star}_m^2(\theta). \tag{3.48}$$

Taking into account (2.5), estimates (3.46) and (3.47) yield

$$\sum_{\eta \in \Theta^*, \eta \subset \theta} |\eta|^{1/2} |f_\eta(x)| \leq c(\nu + 1), \quad x \in E \setminus \text{Star}_{m+\nu}^2(\theta), \quad \nu \geq 0. \tag{3.49}$$

For the proof of (3.44) we also need the estimate

$$|\text{Star}_{m+\nu}^2(\theta) \setminus \theta| \leq c \rho_2^\nu |\theta|, \quad \nu \geq 0, \tag{3.50}$$

where  $0 < \rho_2 < 1$  and  $c > 0$  are constants depending only on the parameters of  $\mathcal{T}$ . Estimate (3.50) is an immediate consequence of estimate (2.4) from Lemma 2.3 and properties (2.1) and (2.2) of  $LR$ -triangulations.

Denote briefly  $F(x) := \sum_{\eta \in \Theta^*, \eta \subset \theta} |\eta|^{1/2} |f_\eta(x)|$ . We have

$$\int_{E \setminus \theta} |F(x)|^q dx = \int_{E \setminus \text{Star}_m^2(\theta)} |F(x)|^q dx + \sum_{\nu=0}^{\infty} \int_{\text{Star}_{m+\nu}^2(\theta) \setminus \text{Star}_{m+\nu+1}^2(\theta)} |F(x)|^q dx. \tag{3.51}$$

Using (3.48) we obtain

$$\int_{E \setminus \text{Star}_m^2(\theta)} |F(x)|^q dx \leq c \sum_{\omega \in \Theta_m: \omega \cap \text{Star}_m^2(\theta) = \emptyset} |\omega| q_*^{q\rho_m(\theta, \omega)} \leq c|\theta|, \tag{3.52}$$

where for the latter estimate we proceed exactly as in (3.17) and (3.18).

By (3.49) and (3.50) we get

$$\sum_{\nu=0}^{\infty} \int_{\text{Star}_{m+\nu}^2(\theta) \setminus \text{Star}_{m+\nu+1}^2(\theta)} |F(x)|^q dx \leq c \sum_{\nu=0}^{\infty} (\nu + 1) \rho_2^\nu |\theta| \leq c|\theta|.$$

This and (3.52) yield (3.44). □

We are now in a position to prove Theorem 3.4.

(a) Assume that (3.3) holds and denote by  $A$  the quantity in (3.3). We shall prove that  $f \in BMO(E, \mathcal{T})$  and  $\|f\|_{BMO} \leq cA$ .

Let  $\theta \in \Theta_m$  ( $m \geq 0$ ) (the case  $\theta = E$  is trivial). We write

$$f = \left( \sum_{\eta \in \Theta^*, \eta \subset \theta} + \sum_{\eta \in \Theta^*, \eta \not\subset \theta, l(\eta) > m} + \sum_{\eta \in \Theta^*, l(\eta) \leq m} \right) \langle f, f_\eta \rangle f_\eta =: F_1 + F_2 + F_3.$$

Using (3.3) we have

$$\|F_1\|_{L_2(E)}^2 = \sum_{\eta \subset \theta} |\langle f, f_\eta \rangle|^2 \leq A^2 |\theta|. \tag{3.53}$$

From (3.3),  $|\langle f, f_\eta \rangle| \leq A |\eta|^{1/2}$  and using (3.45) with  $q = 2$ , we obtain

$$\|F_2\|_{L_2(\theta)} \leq cA |\theta|^{1/2}. \tag{3.54}$$

We use (3.10) and that  $|\langle f, f_\eta \rangle| \leq A |\eta|^{1/2}$  to obtain, for  $x \in \theta$ ,

$$\begin{aligned} |F_3(x) - F_3(v_\theta)| &\leq cA \sum_{\eta \in \Theta^*, l(\eta) \leq m} |\eta|^{1/2} |f_\eta(x) - f_\eta(v_\theta)| \\ &\leq cA \sum_{\nu=0}^m \sum_{\eta \in \Theta_\nu^*} d_{\mathcal{T}}(v_\theta, x)^\varepsilon |\eta|^{-\varepsilon} \left( q_1^{\rho_\nu(\eta, \theta_x^\nu)} + q_1^{\rho_\nu(\eta, \theta_{v_\theta}^\nu)} \right) \\ &\leq cA |\theta|^\varepsilon \sum_{\nu=0}^m \sum_{\eta \in \Theta_\nu^*} |\eta|^{-\varepsilon} q_1^{\rho_\nu(\eta, \theta_x^\nu)}. \end{aligned} \tag{3.55}$$

Now, exactly as in the proof of Lemma 3.11, we have

$$\sum_{\eta \in \Theta_\nu^*} |\eta|^{-\varepsilon} q_1^{\rho_\nu(\eta, \theta_x^\nu)} \leq c |\theta|^{-\varepsilon} \rho^{(m-\nu)\varepsilon} \sum_{\eta \in \Theta_\nu^*} \rho_\nu(\eta, \theta_x^\nu)^{\varepsilon} q_1^{\rho_\nu(\eta, \theta_x^\nu)} \leq c |\theta|^{-\varepsilon} \rho^{(m-\nu)\varepsilon}.$$

We use this in (3.55) to obtain

$$|F_3(x) - F_3(v_\theta)| \leq cA, \quad x \in \theta. \tag{3.56}$$

By (3.53)–(3.56) it readily follows that

$$\left( \frac{1}{|\theta|} \int_{\theta} |f(x) - f(v_{\theta})|^2 dx \right)^{1/2} \leq cA.$$

Consequently,  $\|f\|_{BMO} \leq cA$ .

(b) Assume that  $f \in BMO$  and  $\|f\|_{BMO} = B$ . Fix  $\theta \in \Theta_m$  ( $m \geq 0$ ) (the case  $\theta = E$  is easier). We shall prove that

$$\sum_{\eta \in \Theta^*, \eta \subset \theta} |\langle f, f_{\eta} \rangle|^2 \leq cB^2 |\theta|, \tag{3.57}$$

which implies (3.3).

We write

$$f = \left( \sum_{\eta \in \Theta^*, \eta \subset \theta} + \sum_{\eta \in \Theta^*, \eta \not\subset \theta, l(\eta) > m} + \sum_{\eta \in \Theta^*, l(\eta) \leq m} \right) \langle f, f_{\eta} \rangle f_{\eta} =: F_1 + F_2 + F_3.$$

As was shown in the proof of Theorem 3.2  $\|f_{\eta}\|_{H_1} \leq c|\eta|^{1/2}$ . Then by (2.25) it follows that

$$|\langle f, f_{\eta} \rangle| \leq \|f\|_{BMO} \|f_{\eta}\|_{H_1} \leq cB |\eta|^{1/2}. \tag{3.58}$$

We use this and (3.44) to obtain

$$\begin{aligned} \left| \int_{\theta} F_1(x) dx \right| &\leq \sum_{\eta \in \Theta^*, \eta \subset \theta} |\langle f, f_{\eta} \rangle| \left| \int_{\theta} f_{\eta}(x) dx \right| \\ &\leq cB \sum_{\eta \in \Theta^*, \eta \subset \theta} |\eta|^{1/2} \left| \int_{E \setminus \theta} f_{\eta}(x) dx \right| \\ &\leq cB \left\| \sum_{\eta \in \Theta^*, \eta \subset \theta} |\eta|^{1/2} |f_{\eta}| \right\|_{L_1(E \setminus \theta)} \\ &\leq cB |\theta|^{1/2}. \end{aligned} \tag{3.59}$$

Using (3.45) we have

$$\|F_2\|_{L_2(\theta)} \leq cB \left\| \sum_{\eta \in \Theta^*, \eta \not\subset \theta, l(\eta) > m} |\eta|^{1/2} |f_{\eta}| \right\|_{L_2(\theta)} \leq cB |\theta|^{1/2}. \tag{3.60}$$

Finally, exactly as in (3.56) (using (3.58) instead of  $|\langle f, f_{\eta} \rangle| \leq cA |\eta|^{1/2}$ ), we obtain

$$|F_3(x) - F_3(v_{\theta})| \leq cB, \quad x \in \theta. \tag{3.61}$$

Now, (3.59)–(3.61) readily imply

$$|F_3(v_{\theta}) - f_{\theta}| = \left| F_3(v_{\theta}) - \frac{1}{|\theta|} \int_{\theta} f(x) dx \right| \leq cB.$$

This combined with the definition of BMO yields

$$\left( \frac{1}{|\theta|} \int_{\theta} |f(x) - F_3(v_{\theta})|^2 dx \right)^{1/2} \leq cB.$$



This in turn along with (3.60) and (3.61) implies  $\|F_1\|_{L_2(\theta)} \leq cB |\theta|^{1/2}$ . On the other hand using (3.44) we have

$$\|F_1\|_{L_2(E \setminus \theta)} \leq cB \left\| \sum_{\eta \in \Theta^*, \eta \subset \theta} |\eta|^{1/2} |f_\eta(x)| \right\|_{L_2(E \setminus \theta)} \leq cB |\theta|^{1/2}.$$

Therefore,  $\|F_1\|_{L_2(E)} \leq cB |\theta|^{1/2}$ . Consequently,

$$\sum_{\eta \in \Theta^*, \eta \subset \theta} |\langle f, f_\eta \rangle|^2 = \|F_1\|_{L_2(E)}^2 \leq cB^2 |\theta|,$$

which is (3.57). The proof of Theorem 3.4 is complete. □

### A Appendix: Example of a space $H_1(E, \mathcal{T}) \neq H_1(E)$

Here  $H_1(E)$  denotes the regular  $H_1$ -space on  $E$ . We consider the case  $E := [-1, 1]^2$ .

Denote

$$g(x) = g(x_1, x_2) := \begin{cases} |\ln |x_1||^{1/2} |\ln |x_2||^{1/2} & \text{if } x \in [-1, 1], \\ 0 & \text{if } x \in \mathbb{R}^2 \setminus [-1, 1]. \end{cases} \tag{A.1}$$

In the following we denote by  $BMO(E)$  the regular  $BMO$  space on  $E$ .

**Lemma A.1** *The above defined function  $g(x)$  belongs to  $BMO(\mathbb{R}^2)$ .*

*Proof.* Let  $I := [a, b] \times [c, d]$  with  $b - a = d - c = h > 0$ . We shall prove that there is a constant  $C$  such that

$$\frac{1}{|I|} \int_I |g(x) - C| dx \leq c < \infty, \tag{A.2}$$

where  $c > 0$  is an absolute constant.

We first consider the important situation when  $I \subset [-\frac{1}{4}, \frac{1}{4}] \times [-1, 1] \cup [-1, 1] \times [-\frac{1}{4}, \frac{1}{4}]$  and hence  $0 < h \leq 1/2$ . Note that  $g(-x_1, x_2) = g(x_1, -x_2) = g(x_1, x_2)$ . Therefore, all possibilities for  $I$  are covered by considering the following three cases:

*Case 1:*  $a < 0 < b$  and  $c < 0 < d$ . Because of the symmetry of  $g(x)$  we may assume that  $|a| \leq b$  and  $|c| \leq d$ . By integration by parts we get

$$\int_0^u |\ln t|^{1/2} dt = u |\ln u|^{1/2} + \frac{1}{2} \int_0^u |\ln t|^{-1/2} dt, \quad 0 < u < 1/2,$$

and hence

$$\frac{1}{u} \int_0^u |\ln t|^{1/2} dt = |\ln u|^{1/2} + R_u, \quad \text{where } 0 < R_u \leq \frac{1}{2|\ln u|^{1/2}}. \tag{A.3}$$

Denote  $I_1 := [0, b] \times [0, d]$ ,  $I_2 := [a, 0] \times [0, d]$ ,  $I_3 := [a, 0] \times [c, 0]$ , and  $I_4 := [0, b] \times [c, 0]$ . From the definition of  $g(x)$  it follows that  $0 < g(b, d) = \min_{x \in I} g(x)$  and hence

$$\frac{1}{|I|} \int_I |g(x) - g(b, d)| dx = \frac{1}{|I|} \int_I (g(x) - g(b, d)) dx = \sum_{j=1}^4 \frac{1}{|I|} \int_{I_j} (g(x) - g(b, d)) dx.$$

Using the assumptions and the symmetry of  $g(x)$ , each integral  $\int_{I_j} g(x) dx$  can be written in the form

$$\int_{I_j} g(x) dx = \int_0^u \int_0^v |\ln x_1|^{1/2} |\ln x_2|^{1/2} dx_1 dx_2$$

for some  $u, v$  satisfying  $0 \leq u \leq b$  and  $0 \leq v \leq d$ . Then using (A.3), we have

$$\begin{aligned} & \frac{1}{|I|} \int_{I_j} (g(x) - g(b, d)) dx \\ &= uv |I|^{-1} \left[ (|\ln u|^{1/2} + R_u)(|\ln v|^{1/2} + R_v) - |\ln b|^{1/2} |\ln d|^{1/2} \right] \\ &\leq uv (bd)^{-1} \left[ |\ln u|^{1/2} |\ln v|^{1/2} - |\ln b|^{1/2} |\ln d|^{1/2} \right] \\ &\quad + uv (bd)^{-1} |\ln u|^{1/2} R_v + uv (bd)^{-1} |\ln v|^{1/2} R_u + uv (bd)^{-1} R_u R_v \\ &=: A_1 + A_2 + A_3 + A_4. \end{aligned}$$

To estimate  $A_1$  we substitute  $u = sb$  and  $v = td$ ,  $0 < s, t \leq 1$ , and obtain

$$\begin{aligned} A_1 &= st \left[ (|\ln s| + |\ln b|)^{1/2} (|\ln t| + |\ln d|)^{1/2} - |\ln b|^{1/2} |\ln d|^{1/2} \right] \\ &= \frac{st \left[ (|\ln s| + |\ln b|)(|\ln t| + |\ln d|) - |\ln b| |\ln d| \right]}{(|\ln s| + |\ln b|)^{1/2} (|\ln t| + |\ln d|)^{1/2} + |\ln b|^{1/2} |\ln d|^{1/2}} \\ &\leq \frac{st (|\ln s| |\ln t| + B |\ln t| + D |\ln s|)}{(BD)^{1/2}}, \end{aligned}$$

where  $B := |\ln b|$  and  $D := |\ln d|$ . Since  $b - a = d - c = h > 0$ ,  $|a| \leq b$  and  $|c| \leq d$ , then  $\ln 2 \leq |\ln h| \leq B$  and  $D \leq |\ln(h/2)| \leq 2 |\ln h|$ . We use this above to obtain

$$A_1 \leq (\ln 2)^{-1} s |\ln s| \cdot t |\ln t| + 2t |\ln t| + 2s |\ln s| \leq c < \infty.$$

To estimate  $A_2$  we use (A.3) and again replace  $u$  and  $v$  by  $u = sb$  and  $v = td$  ( $0 < s, t \leq 1$ ), respectively. We have

$$A_2 \leq \frac{uv |\ln u|^{1/2}}{bd |\ln v|^{1/2}} = st \left( \frac{|\ln s| + |\ln b|}{|\ln t| + |\ln d|} \right)^{1/2} \leq 2st \left( \frac{|\ln s| + |\ln h|}{|\ln t| + |\ln h|} \right)^{1/2}.$$

Since  $1 < |\ln h| < \infty$  ( $0 < h \leq 1/2$ ), it is readily seen that

$$\frac{|\ln s| + |\ln h|}{|\ln t| + |\ln h|} \leq \frac{|\ln s| + 1}{|\ln t| + 1} + 1 \leq |\ln s| + 2.$$

Consequently,

$$A_2 \leq 2s (|\ln s| + 2)^{1/2} \leq c < \infty.$$

Exactly in the same way we get  $A_3 \leq c < \infty$ . Also, by (A.3),

$$A_4 \leq \frac{uv}{4bd (|\ln u| |\ln v|)^{1/2}} \leq \frac{1}{4 \ln 2}.$$

The above estimates for  $A_1, A_2, A_3$  and  $A_4$  imply (A.2) with  $C := g(b, d)$  and  $c > 0$  an absolute constant.

*Case 2:*  $0 \leq a < b$  and  $0 \leq c < d$ . In this case we shall make use of the following simple identity

$$\int_a^b |\ln t|^{1/2} dt = (b - a) |\ln b|^{1/2} + \frac{1}{2} \int_a^b \frac{t - a}{t |\ln t|^{1/2}} dt, \tag{A.4}$$

which can be verified by differentiating both sides with respect to  $b$ . We next use (A.4) to prove the following ( $h := b - a$ ):

$$\frac{1}{h} \int_a^b |\ln t|^{1/2} dt = |\ln b|^{1/2} + R_{ab}, \quad \text{where } 0 < R_{ab} \leq \frac{2}{|\ln h|^{1/2}}. \tag{A.5}$$

Note first that from our assumptions it follows that  $h \leq 1/4$ . If  $0 \leq a \leq h$ , then by (A.4)

$$R_{ab} \leq \frac{1}{2h} \int_a^b \frac{1}{|\ln t|^{1/2}} dt \leq \frac{1}{2|\ln b|^{1/2}} \leq \frac{1}{2|\ln 2h|^{1/2}} \leq \frac{1}{|\ln h|^{1/2}}.$$

If  $h < a \leq 1/4$ , then (A.4) implies

$$R_{ab} \leq \frac{1}{2} \int_a^b \frac{1}{t|\ln t|^{1/2}} dt \leq \frac{h}{2h|\ln h|^{1/2}} = \frac{1}{2|\ln h|^{1/2}},$$

where we used that the function  $t|\ln t|^{1/2}$  is increasing on  $(0, 1/2]$  and  $b = a + h \leq 1/2$ . Finally, (A.5) is trivial if  $a > 1/4$ . Thus (A.5) holds true.

Evidently,  $0 < g(b, d) = \min_{x \in I} g(x)$  and hence

$$\begin{aligned} \frac{1}{|I|} \int_I |g(x) - g(b, d)| dx &= \frac{1}{|I|} \int_I (g(x) - g(b, d)) dx \\ &= \frac{1}{h^2} \int_a^b |\ln x_1|^{1/2} dx_1 \cdot \int_c^d |\ln x_2|^{1/2} dx_2 - |\ln b|^{1/2} |\ln d|^{1/2}. \end{aligned}$$

Now, employing (A.5), we obtain

$$\begin{aligned} \frac{1}{|I|} \int_I |g(x) - g(b, d)| dx &= (|\ln b|^{1/2} + R_{ab})(|\ln d|^{1/2} + R_{cd}) - |\ln b|^{1/2} |\ln d|^{1/2} \\ &= |\ln b|^{1/2} R_{cd} + |\ln d|^{1/2} R_{ab} + R_{ab}R_{cd} \\ &\leq 2|\ln b|^{1/2}/|\ln h|^{1/2} + 2|\ln d|^{1/2}/|\ln h|^{1/2} + 4/|\ln h| \\ &\leq c < \infty, \end{aligned}$$

where we used that  $b, d \geq h$  and  $0 < h \leq 1/4$ . Thus (A.2) holds true.

Case 3:  $0 \leq a < b$  and  $c \leq 0 < d$ . In this case one proceeds as above using (A.3) and (A.5). We omit the details.

It remains to show that (A.2) holds whenever  $I \not\subset [-\frac{1}{4}, \frac{1}{4}] \times [-1, 1] \cup [-1, 1] \times [-\frac{1}{4}, \frac{1}{4}]$ . If in this case  $h < 1/8$  and  $I \subset [-1, 1]^2$ , then (A.2) is obvious since  $\|g\|_{L^\infty(I)} \leq g(1/8, 1/8) = \ln 8$ .

If  $h > 1/8$ , then  $\int_I g(x) dx \leq \int_{[-1, 1]^2} g(x) dx = c < \infty$  and (A.2) follows.

If  $I^\circ \cap [-1, 1]^2 = \emptyset$ , then (A.2) is again obvious.

Finally, suppose that  $h \leq 1/8$ ,  $I \not\subset [-1, 1]^2$ , and  $I \cap [-1, 1]^2 \neq \emptyset$ . Then evidently there is a square  $J = [\alpha, \beta] \times [\gamma, \delta]$  of the same size as  $I$  such that  $J \subset [-1, 1]^2$  and  $I \cap [-1, 1]^2 \subset J$ . Then by the monotonicity of  $g(x)$ , we have

$$\frac{1}{|I|} \int_I g(x) dx \leq \frac{1}{|J|} \int_J g(x) dx \leq c < \infty,$$

where we used the results of Case 2 or Case 3 above. This yields (A.2) with  $C = 0$ . The prove of the lemma is complete.  $\square$

Armed with this lemma, we proceed to showing that there is an  $LR$ -triangulation  $\mathcal{T}$  of  $E := [-1, 1]^2$  such that  $H_1(E, \mathcal{T}) \neq H_1(E)$ .

From [9] (see the construction in the beginning of §2.1) it follows that there exists an  $LR$ -triangulation  $\mathcal{T}$  of  $E$  with the property: There is a sequence of cells in  $\mathcal{T}$ :  $\theta_1 \supset \theta_2 \supset \dots$  such that

$$[-\lambda_\nu/2, \lambda_\nu/2] \times [-\varepsilon_\nu/2, \varepsilon_\nu/2] \subset \theta_\nu \subset [-\lambda_\nu, \lambda_\nu] \times [-\varepsilon_\nu, \varepsilon_\nu],$$

where  $1/4 \geq \lambda_1 > \lambda_2 > \dots > 0$ ,  $1/4 \geq \varepsilon_1 > \varepsilon_2 > \dots > 0$ ,  $\lim_{\nu \rightarrow \infty} \lambda_\nu = 0$ ,  $\lim_{\nu \rightarrow \infty} \varepsilon_\nu = 0$ , and  $\lim_{\nu \rightarrow \infty} \varepsilon_\nu/\lambda_\nu = 0$ . In addition to this each  $\theta_\nu$  is convex and symmetric with respect to the  $x_1$ -axis and the  $x_2$ -axis.

Let  $\theta_\nu^+ := \{(x_1, x_2) \in \theta_\nu : x_1 \geq 0\}$  and  $\theta_\nu^- := \{(x_1, x_2) \in \theta_\nu : x_1 < 0\}$ . The functions  $a_\nu := |\theta_\nu|^{-1}(\mathbb{1}_{\theta_\nu^+} - \mathbb{1}_{\theta_\nu^-})$  are obviously atoms in  $H_1(E, \mathcal{T})$ .

Now, fix  $\nu \geq 1$  and denote briefly  $\theta := \theta_\nu$ ,  $\lambda := \lambda_\nu$ ,  $\varepsilon := \varepsilon_\nu$ ,  $a := a_\nu$ , etc. We next estimate from below  $\|a\|_{H_1(E)}$ , where  $H_1(E)$  is the regular  $H_1$ -space on  $E$ . By Lemma A.1 the function  $g$  from (A.1) is in  $BMO(\mathbb{R}^2)$  and hence all functions obtained from  $g$  by dilations and shifts also belong to  $BMO(\mathbb{R}^2)$  and have the same BMO norm. In particular,  $g_\lambda(x) := g(x_1/\lambda + 1/2, x_2/\lambda)$  belongs to  $BMO(\mathbb{R}^2)$  and  $\|g_\lambda\|_{BMO} = \|g\|_{BMO}$ . Therefore, the restriction of  $g_\lambda$  on  $E := [-1, 1]$  (that we denote again by  $g_\lambda$ ) belongs to  $BMO(E)$  and  $\|g_\lambda\|_{BMO(E)} \leq \|g\|_{BMO(\mathbb{R}^2)}$ .

Clearly,

$$\|a\|_{H_1(E)} = \sup_{\varphi \in BMO(E)} \frac{\langle a, \varphi \rangle}{\|\varphi\|_{BMO}} \geq \frac{\langle a, g_\lambda \rangle}{\|g_\lambda\|_{BMO}} \geq c \int_\theta a(x) g_\lambda(x) dx.$$

Since  $g(x_1, -x_2) = g(x_1, x_2)$  and  $g(x_1, x_2)$  is monotone decreasing with respect to  $x_1$  on  $(0, 1)$ , we have

$$\|a\|_{H_1(E)} \geq \frac{c}{\lambda \varepsilon} \int_0^{\varepsilon/2} \left[ \int_{-\lambda/2}^0 g\left(\frac{x_1}{\lambda} + \frac{1}{2}, \frac{x_2}{\lambda}\right) dx_1 - \int_0^{\lambda/2} g\left(\frac{x_1}{\lambda} + \frac{1}{2}, \frac{x_2}{\lambda}\right) dx_1 \right] dx_2.$$

Substituting  $y_1 := x_1/\lambda + 1/2$  and  $y_2 := x_2/\lambda$  we infer

$$\begin{aligned} \|a\|_{H_1(E)} &\geq \frac{c\lambda}{\varepsilon} \int_0^{\varepsilon/\lambda} \left( \int_0^{1/2} g(y_1, y_2) dy_1 - \int_{1/2}^1 g(y_1, y_2) dy_1 \right) dy_2 \\ &= \frac{c\lambda}{\varepsilon} \int_0^{\varepsilon/\lambda} |\ln y_2|^{1/2} dy_2 \left( \int_0^{1/2} |\ln y_1|^{1/2} dy_1 - \int_{1/2}^1 |\ln y_1|^{1/2} dy_1 \right) \\ &\geq c_1 (\ln \lambda / \varepsilon)^{1/2}, \end{aligned}$$

where  $c_1 > 0$  is an absolute constant and for the last estimate we used (A.3). Thus there is a sequence of atoms  $(a_\nu)_{\nu=1}^\infty$  in  $H_1(E, \mathcal{T})$  such that  $\|a_\nu\|_{H_1(E)} \rightarrow \infty$  as  $\nu \rightarrow \infty$ , which leads to the conclusion that  $H_1(E, \mathcal{T}) \neq H_1(E)$ .

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