# Interpolation by Piecewise-Linear Radial Basis Functions, I

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In the two-dimensional plane, a set of points  $x_1, x_2, ..., x_n$  (called "nodes") is given. It is desired to interpolate arbitrary data given on the nodes by continuous functions having piecewise-linear (" $\mathscr{P}\mathscr{L}$ ") structure. For this purpose, one can employ the space of all  $\mathscr{P}\mathscr{L}$ -functions on a rectangular grid generated by the nodes. We study this space first. Next, we investigate the special  $\mathscr{P}\mathscr{L}$ -functions that are linear combinations of functions  $h_i(x) = ||x - x_i||_1$ , in which the  $l_1$ -norm on  $\mathbb{R}^2$  is employed. The "dual" case, involving the two-dimensional  $l_\infty$ -norm, is included in our results, as are certain general interpolating functions of the form

$$(s, t) \mapsto F(s-s_i) + G(t-t_i).$$

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#### 1. Introduction

Throughout the paper,  $\mathcal{N}$  denotes a set of n distinct points in  $\mathbb{R}^2$  designated by  $x_1, x_2, ..., x_n$ . These points are called *nodes*. (In Section 9, we consider nodes in  $\mathbb{R}^d$  for  $d \ge 2$ .) The basic problem of two-dimensional

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interpolation addressed here is as follows. A "data-function"  $d: \mathcal{N} \to \mathbb{R}$  is given, and we seek a function  $f: \mathbb{R}^2 \to \mathbb{R}$  such that  $f \mid \mathcal{N} = d$ ; i.e.,  $f(x_i) = d_i$  for i = 1, 2, ..., n. Such a function f is said to *interpolate* d. Usually the search for f is restricted to a class of functions that (a) are easily computed and (b) have some prescribed smoothness.

If the set of nodes has no special structure capable of being exploited, then this problem is called scattered data interpolation. Many methods proposed for this problem are discussed in the surveys of Schumaker [13] and Franke [4, 5]. One method that has been used successfully employs radial basis functions. In the simplest case of this, one seeks an interpolant in the linear space generated by the *n* functions  $h_i(x) = ||x - x_j|| \ (1 \le j \le n)$ , where the norm can be any convenient one on  $\mathbb{R}^2$ . The existence of an interpolant  $f = \sum_{j=1}^{n} c_j h_j$  for arbitrary data depends upon the invertibility of the interpolation matrix A, whose elements are  $A_{ij} = h_j(x_i)$ . For the Euclidean norm, a result of Schoënberg [11] asserts that the matrix A is always nonsingular. (Schoenberg's result holds in any inner-product space.) Micchelli [9] proved a striking generalization of this result in which  $h_j$  can be replaced by  $h_i(x) = F(||x - x_i||^2)$ , where F comes from a certain class of functions. The papers of Powell [10], Jackson [6], Madych and Nelson [7, 8], Dyn [1], and Dyn et al. [2] contain further important contributions to this field. Practical experience with this type of interpolation is reported by Franke [3, 4].

We consider radial basis functions that are generated by the  $l_1$ -norm. Thus, if x = (s, t) and  $x_i = (s_i, t_i)$ , then

$$h_i(x) = ||x - x_i||_1 = |s - s_i| + |t - t_i|$$
  $(1 \le i \le n).$ 

Since these functions are piecewise linear on a rectangular grid, we devote several sections (2–5) to a study of piecewise linear functions in general. Sections 6, 7, 9 concern the space of radial basis functions, and emphasize its role as a linear subspace in the space of piecewise linear functions. Section 8 is devoted to radial basis functions employing the  $l_{\infty}$ -norm. Our results provide a geometric property of  $\mathcal N$  that is necessary and sufficient for the invertibility of the interpolation matrix.

The following notation is adopted. Orthogonal projections onto the coordinate axes are denoted by P and Q. Explicitly,

$$Px = s$$
,  $Qx = t$ ,  $x = (s, t) \in \mathbb{R}^2$ .

The projections of  $\mathcal{N}$  are denoted by

$$P(\mathcal{N}) = \{\sigma_1, \sigma_2, ..., \sigma_m\}, \qquad \sigma_1 < \sigma_2 < \dots < \sigma_m$$

$$Q(\mathcal{N}) = \{\tau_1, \tau_2, ..., \tau_k\}, \qquad \tau_1 < \tau_2 < \dots < \tau_k.$$

The rectangular grid and the rectangular hull determined by the node set  $\mathcal{N}$  are the sets

$$\mathbf{RG} = \{\sigma_1, \sigma_2, ...., \sigma_m\} \times \{\tau_1, \tau_2, ..., \tau_k\}$$

$$\mathbf{RH} = \{(s, t) : \sigma_1 \leqslant s \leqslant \sigma_m \quad \text{and} \quad \tau_1 \leqslant t \leqslant \tau_k\}.$$

It is assumed always that  $\# \mathcal{N} = n$  (i.e., the nodes  $x_1, ..., x_n$  are distinct) and that  $m \ge 2, k \ge 2$ .

### 2. The Space $\mathscr{PL}$ of Piecewise Linear Functions

The horizontal and vertical lines through the points of  $\mathcal{N}$  divide the plane into rectangles, some of which are unbounded. These rectangles are expressible as Cartesian products of intervals. There are m+1 such intervals on the s-axis and there are k+1 intervals on the t-axis. The space  $\mathscr{PL}(\mathcal{N})$ , or simply  $\mathscr{PL}$ , is defined to be the space of all continuous functions  $f: \mathbb{R}^2 \to \mathbb{R}$  such that the restriction of f to each of these rectangles is a linear function of (s, t).

The dimension of the space of all piecewise linear functions is obviously 3(m+1)(k+1) since there are (m+1)(k+1) rectangles, and a linear function has three coefficients. On the other hand, if continuity is imposed, there will be three conditions required at each grid point to ensure that the linear functions in four adjacent rectangles are equal there. There are 3mk conditions of this type. Also, on each of the semi-infinite lines which emanate from RH one continuity condition must be imposed. This provides 2(m+k) further conditions. The number of parameters minus the number of conditions is m+k+3, and one can prove that this is indeed the dimension of  $\mathcal{PL}$ . To this end, we now define  $\sigma_0$ ,  $\tau_0$ ,  $\sigma_{m+1}$ , and  $\tau_{k+1}$  to be any real numbers satisfying

$$\sigma_0 < \sigma_1, \quad \sigma_m < \sigma_{m+1}, \quad \tau_0 < \tau_1, \quad \tau_k < \tau_{k+1}.$$

2.1. THEOREM. A  $\mathscr{PL}$ -function is uniquely determined by assigning to it arbitrary values at the points  $(\sigma_0, \tau_j)$  and  $(\sigma_i, \tau_0)$ , where  $0 \le j \le k+1$  and  $1 \le i \le m+1$ . Consequently,  $\dim(\mathscr{PL}) = m+k+3$ .

*Proof.* Let  $f \in \mathcal{PL}$ . By 5.1, f can be written

$$f(s, t) = u(s) + v(t),$$

where u is a piecewise linear function in  $C(\mathbb{R})$  having knots  $\sigma_1, \sigma_2, ..., \sigma_m$ , and v is a piecewise linear function in  $C(\mathbb{R})$  having knots  $\tau_1, \tau_2, ..., \tau_k$ . By

adding a constant to u and subtracting it from v we can arrange that  $u(\sigma_0) = 0$ . Then v is uniquely determined by the equation

$$v(\tau_i) = f(\sigma_0, \tau_i), \quad 0 \le j \le k+1.$$

After that, u is uniquely specified by the equation

$$u(\sigma_i) = f(\sigma_i, \tau_0) - v(\tau_0), 1 \le i \le m+1.$$

Theorem 2.1 is already known but is included here for completeness. See Schumaker [15].

- 2.2. COROLLARY. Let  $S = \{\sigma_0, \sigma_1, ..., \sigma_{m+1}\}$  and  $T = \{\tau_0, \tau_1, ..., \tau_{k+1}\}$ . Let  $\Pi_0(S)$  and  $\Pi_0(T)$  denote the spaces of constant functions on S and T, respectively. The space  $\mathscr{PL}$ , when interpreted as a subspace of  $C([\sigma_0, \sigma_{m+1}] \times [\tau_0, \tau_{k+1}])$ , is isometrically isomorphic to  $l_\infty(S) \otimes \Pi_0(T) + \Pi_0(S) \otimes l_\infty(T)$  as a subspace of  $l_\infty(S \times T)$ .
- 2.3. THEOREM. When  $\mathscr{PL}$  is restricted to the rectangular hull,  $[\sigma_1, \sigma_m] \times [\tau_1, \tau_k]$ , its dimension is m + k 1.

*Proof.* This follows at once from 2.1 upon changing k to k-2 and m to m-2.

2.4. THEOREM. Every PL-function can be written uniquely in the form

$$f(s,t) = \sum_{i=1}^{m} \alpha_i |s - \sigma_i| + \sum_{i=1}^{k} \beta_i |t - \tau_i| + as + bt + c.$$
 (1)

*Proof.* The function on the right in this equation is a linear combination of m+k+3 functions. By 2.1, dim  $\mathscr{PL}=m+k+3$ . Thus it suffices to prove that every  $\mathscr{PL}$ -function has a representation as claimed. Let f be any  $\mathscr{PL}$ -function. Then f can be written in the form

$$f(s,t) = u(s) + v(t) \tag{2}$$

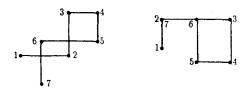
with u and v piecewise linear functions having knots  $\sigma_1 \cdots \sigma_m$  and  $\tau_1 \cdots \tau_k$ , respectively. By a well-known theorem in spline theory, Schumaker [14], we can write

$$u(s) = \sum_{i=1}^{m} \alpha_i |s - \sigma_i| + as + c', \qquad v(t) = \sum_{i=1}^{k} \beta_i |t - \tau_i| + bt + c''.$$

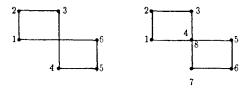
## 3. PATHS AND PATH FUNCTIONALS

A path is a finite ordered set in RG,  $[y_1, y_2, ..., y_r]$ , such that the line segments joining consecutive points are of positive length and are alternately horizontal and vertical. (Repetitions of points are permitted.) The number r is the *length* of the path.

Pictures of paths are shown below.



A path is said to be *closed* if r is even, if  $y_r \neq y_1$ , and if the line segment joining  $y_1$  with  $y_2$  is perpendicular to the line segment joining  $y_r$  with  $y_1$ . Typical closed paths are shown below.



An equivalent formulation of the definition is this: the path  $[y_1, y_2, ..., y_r]$  is closed if r is even and if  $[y_1, y_2, ..., y_r, y_1]$  is also a path. If  $[y_1, y_2, ..., y_r]$  is a closed path, then a linear functional, called a path functional, is associated with it as follows:

$$\phi = \sum_{i=1}^{r} (-1)^{i} \hat{y}_{i}.$$

Here  $\hat{y}$  denotes the point-evaluation functional associated with the point y; i.e., for any function f whose domain includes y,

$$\hat{y}(f) = f(y).$$

If r is 4, the path functional is called "the 4-point rule."

3.1. Lemma. Any path of length n+1 in  $\mathcal N$  contains a closed path.

*Proof.* Select a path of length n+1 in  $\mathcal{N}: [z_0, z_1, ..., z_n]$ . Since

 $\# \mathcal{N} = n$ , there is a first index i such that  $z_i \in \{z_0, z_1, ..., z_{i-1}\}$ . Select j < i such that  $z_i = z_i$ . Now there are several cases.

If  $t_j = t_{j+1}$  and  $t_i = t_{i-1}$  then  $[z_{j+1}, z_{j+2}, ..., z_{i-1}]$  is a closed path. To verify this, notice that the line segment from  $z_j$  to  $z_{j+1}$  is horizontal and the line segment from  $z_{j+1}$  to  $z_{j+2}$  must therefore be vertical. Next, since  $t_{i-1} = t_i = t_j = t_{j+1}$ , we see that the line segment from  $z_{i-1}$  to  $z_{j+1}$  is horizontal. The number of entries in the ordered set is even, by the following reasoning. The segments joining  $z_v$  to  $z_{v+1}$  in the path have vertical and horizontal orientations in the pattern

### $VHVHVH\cdots VHV$ .

Hence there is an odd number of segments and an even number of points  $z_{i+1}, ..., z_{i-1}$ .

If  $t_j = t_{j+1}$  and  $s_i = s_{i-1}$ , a closed path is  $[z_j, z_{j+1}, ..., z_{i-1}]$ . Indeed the segment from  $z_j$  to  $z_{j+1}$  is horizontal, and the segment from  $z_{i-1}$  to  $z_i = z_j$  is vertical. The number of entries in the ordered set is again even.

There are two remaining cases described by the conditions  $(s_i = s_{j+1} \text{ and } s_i = s_{i-1})$  and  $(s_j = s_{j+1} \text{ and } t_i = t_{i-1})$ . These require no further proof since they follow from the first two cases upon interchanging s and t.

3.2. LEMMA. Let A be a nonvoid subset of  $\mathcal{N}$  such that  $\#(A \cap L) \neq 1$  for any horizontal or vertical line, L. Then A contains a closed path.

*Proof.* Assume the hypotheses, and select any point  $z_0$  in A. The vertical line L through  $z_0$  must satisfy  $\#(A \cap L) \ge 2$ , and one can select  $z_1 \in A \cap L$  with  $z_1 \ne z_0$ . Similarly, one can select  $z_2$  on the horizontal line through  $z_1$ , with  $z_2 \ne z_1$ . We continue in this way until we have a list of n+1 points  $[z_0, z_1, ..., z_n]$ . An application of 3.1 completes the proof.

- 3.3. LEMMA. For a subset Z of the rectangular grid the following properties are equivalent:
  - (1) Z contains a closed path,
- (2) There is a nontrivial functional supported on Z that annihilates  $\mathscr{PL}$ .

*Proof.* If (2) is true, let  $\phi \equiv \sum_{i=1}^r a_i \hat{z}_i$  annihilate  $\mathscr{PL}$ , where  $z_i \in Z$  and all coefficients  $a_i$  are nonzero. Put  $Z' = \{z_1, ..., z_r\}$ . We observe that every horizontal and vertical line that intersects Z' does so in at least two points. Indeed, if (for example) the vertical line through  $z_j$  contains no other  $z_i$ , then an element f in  $\mathscr{PL}$  can be constructed that is a function of s only and satisfies  $f(z_i) = \delta_{ij}$ . Then we have the contradiction  $0 = \phi(f) = a_j$ . An application of 3.2 establishes (1).

Now suppose that Z contains a closed path,  $[z_1, ..., z_r]$ . Then the functional  $\phi \equiv \sum_{i=1}^{r} (-1)^i \hat{z}_i$  annihilates C(S) + C(T), as is easily verified by considered functions of s and functions of t separately. By 2.4  $\mathscr{PL}$  is a subspace of C(S) + C(T), and so is annihilated by  $\phi$ .

3.4. Lemma. Any subset of the grid that contains m+k points must contain a closed path.

*Proof.* Let A be such a set. Let S' be the set of all  $\sigma_i$  such that the vertical line through  $\sigma_i$  contains at least two points of A. If #S' = 0, then each vertical grid line contains at most one point of A, and thus  $\#A \leq m$ , contrary to hypotheses. If #S' = 1, then one vertical line contains  $\nu$  points of A, with  $\nu \geq 2$ . Each of the remaining m-1 vertical grid lines contains at most one point of A. Since #T = k we have the contradiction

$$m + k \le \# A \le m - 1 + v \le m - 1 + k$$
.

Thus we conclude that  $\#S' \ge 2$ . Similarly  $\#T' \ge 2$ , where T' is the set of  $\tau_j$  such that the horizontal line through  $\tau_j$  contains at least two points of A. By 3.2, the set  $A \cap (S' \times T')$  contains a closed path.

3.5. LEMMA. Let g and h belong to  $C(\mathbb{R})$ . Let  $\mathcal{N}$  be a set of nodes,  $x_i = (s_i, t_i)$ , in  $\mathbb{R}^2$ . Let

$$H_i(s, t) = g(s - s_i) + h(t - t_i) \qquad (1 \le i \le n).$$

If  $\mathcal N$  contains a closed path, then the functions  $H_i$  form a dependent (indexed) set.

*Proof.* Renumber the nodes, if necessary, so that  $[x_1, ..., x_q]$  is a closed path. It is easily seen that the functional  $\phi = \sum_{i=1}^{q} (-1)^i \hat{x}_i$  annihilates C(S) + C(T). Hence it annihilates the function given by  $(s, t) \mapsto g(\sigma - s) + h(\tau - t)$  where  $(\sigma, \tau)$  is any fixed point in  $\mathbb{R}^2$ . Thus

$$\sum_{i=1}^{q} (-1)^{i} \left\{ g(\sigma - s_{i}) + h(\tau - t_{i}) \right\} = 0$$

and therefore

$$\sum_{i=1}^{q} (-1)^{i} H_{i}(\sigma, \tau) = 0, \qquad (\sigma, \tau) \in \mathbb{R}^{2}. \quad \blacksquare$$

3.6. LEMMA. If a path  $[z_1, z_2, ..., z_r]$  satisfies  $r \ge 4$  and if either  $P(z_1) = P(z_r)$  or  $Q(z_1) = Q(z_r)$ , then it contains a closed subpath.

*Proof.* Select integers p and q in  $\{1, 2, ..., r\}$  to minimize q - p under the constraints

- $(1) \quad q p \geqslant 3,$
- (2)  $[P(z_p) P(z_q)][Q(z_p) Q(z_q)] = 0.$

We shall prove that  $z_p \neq z_q$ . Suppose that  $z_p = z_q$ . Then  $q - p \geqslant 4$ . Since (p, q) is minimal, (p + 1, q) does not satisfy the constraints. Since it satisfies (1), it does not satisfy (2). Hence  $P(z_{p+1}) \neq P(z_q)$  and  $Q(z_{p+1}) \neq Q(z_q)$ . These inequalities violate the definition of a path, and we conclude that  $z_p \neq z_q$ .

Next we prove that  $[z_p, z_{p+1}, ..., z_q]$  is a closed path. Note that  $z_p \neq z_q$  by the preceding paragraph. With no loss of generality, we assume that the first factor in (2) is 0. By the minimality of (p, q), the segment joining  $z_p$  to  $z_{p+1}$  and the segment joining  $z_{q-1}$  to  $z_q$  are both horizontal. Hence the number of elements in the ordered set  $[z_p, ..., z_q]$  is even. Finally, the segent from  $z_p$  to  $z_q$  is vertical and hence perpendicular to the segment from  $z_p$  to  $z_{p+1}$ .

#### 4. Interpolation Properties of $\mathscr{PL}$

The dimension of  $\mathscr{PL}$  when restricted to the rectangular hull of  $\mathscr{N}$  is m+k-1, as shown in 2.3. It is of interest to know what sets of m+k-1 grid points are suitable as nodes for interpolation. In this section we answer this question by exploiting the intuitive geometric idea of a path. Although one would normally expect to compute the coefficients in an interpolating function by inverting a linear system of order m+k-1, a much more economical algorithm is available.

4.1. THEOREM. The space  $\mathscr{PL}$  can interpolate arbitrary data on a set of m+k-1 grid points if and only if that set does not contain a closed path.

**Proof.** Let Y be a set of m+k-1 grid points,  $\{y_1, ..., y_{m+k-1}\}$ . Let  $\{f_j: 1 \le j \le m+k-1\}$  be a basis for  $\mathscr{PL}$  restricted to RG. Interpolation is possible if and only if the matrix  $(f_j(y_i))$  is nonsingular. If Y contains a closed path, then by 3.3, a nontrivial linear combination of the point functionals  $\hat{y}_i$  annihilates  $\mathscr{PL}$ , and the matrix in question is singular. Conversely, if the matrix is singular then a linear combination of its rows is 0, and this gives an annihilating functional supported on Y. By 3.3, Y contains a closed path.

Now let us assume that  $\mathcal{N}$  contains no closed path and that  $n \le m+k-1$ . We shall give an algorithm which produces a continuous  $\mathscr{PL}$  interpolant for arbitrary data given on the node set  $\mathcal{N} =$ 

 $\{x_1, x_2, ..., x_n\}$ . As in classical polynomial interpolation by the Lagrange method, it suffices to construct functions  $l_1, ..., l_n$  in  $\mathscr{PL}(\mathcal{N})$  with the "cardinal property"

$$l_i(x_i) = \delta_{ii}$$
  $(1 \le i, j \le n).$ 

For i=1,2,...,n, define the set  $\mathscr{H}_i$  to consist of  $x_i$  and all other nodes that can be connected to  $x_i$  by a path in  $\mathscr{N}$  starting at  $x_i$  with a horizontal segment. Similarly,  $\mathscr{V}_i$  contains  $x_i$  and all nodes that can be reached along a path in  $\mathscr{N}$  starting with a vertical segment at  $x_i$ . Figure 4.1 shows a typical situation; the nodes labeled v belong to  $\mathscr{V}_i$  and the nodes labeled v belong to  $\mathscr{H}_i$ ,

Now define  $a_i$  on  $P(\mathcal{N})$  and  $b_i$  on  $Q(\mathcal{N})$  as follows:

$$a_i(s) = \begin{cases} +1 & \text{if} \quad s \in P(\mathcal{V}_i) \\ -1 & \text{if} \quad s \in P(\mathcal{H}_i \setminus x_i), \\ 0 & \text{otherwise} \end{cases} \quad b_i(t) = \begin{cases} +1 & \text{if} \quad t \in Q(\mathcal{H}_i) \\ -1 & \text{if} \quad t \in Q(\mathcal{V}_i \setminus x_i). \\ 0 & \text{otherwise}. \end{cases}$$

Finally, define  $l_i(s, t) = \frac{1}{2} [a_i(s) + b_i(t)].$ 

4.2. THEOREM. If  $\mathcal{N}$  contains no closed path, then  $a_i$  and  $b_i$  are well defined, and  $l_i(x_i) = \delta_{ii}$ .

*Proof.* Let  $\sigma \in P(\mathcal{N})$ , and suppose that  $a_i(\sigma)$  is not well defined. Then  $\sigma \in P(\mathcal{V}_i) \cap P(\mathcal{H}_i \setminus x_i)$ . Hence there exist nodes  $x_j$  and  $x_v$  such that  $x_v \in \mathcal{V}_i$ ,  $x_j \in \mathcal{H}_i \setminus x_i$ , and  $P(x_v) = P(x_j) = \sigma$ . (Under some circumstances,  $x_v$  can be  $x_i$ .) There exists a path, starting at  $x_v$ , progressing to  $x_i$  and then to  $x_j$ . By 3.6, such a path contains a closed subpath, contrary to hypotheses.

Now observe that  $s_i = P(x_i) \in P(\mathcal{V}_i)$  and  $t_i = Q(x_i) \in Q(\mathcal{H}_i)$ . Consequently

$$l_i(x_i) = \frac{1}{2} [a_i(s_i) + b_i(t_i)] = \frac{1}{2} (1+1) = 1.$$

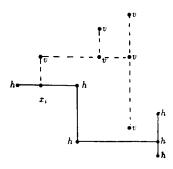


FIGURE 4.1

On the other hand, let x = (s, t) be any node other than  $x_i$ . If there is a path from  $x_i$  to x, then either  $x \in \mathcal{X}_i$  or  $x \in \mathcal{V}_i$ . In these two cases we have

$$2l_i(x) = a_i(s) + b_i(t) = -1 + 1 = 0,$$
  $x \in \mathcal{H}_i$   
 $2l_i(x) = a_i(s) + b_i(t) = 1 - 1 = 0,$   $x \in \mathcal{V}_i$ 

If x is not connected by any path to  $x_i$  then  $s \notin P(\mathcal{V}_i)$ ,  $s \notin P(\mathcal{H}_i)$ ,  $t \notin Q(\mathcal{H}_i)$ , and  $t \notin Q(\mathcal{V}_i)$ . Consequently  $l_i(x) = 0$ .

# 5. Degree of Approximation by $\mathscr{PL}$ -Interpolants

In the preceding section, a necessary and sufficient condition was given on a set of grid points in order that  $\mathscr{PL}$ -interpolation at those points would be possible. The natural question arises of what happens when the grid is refined and the set of interpolation points is chosen to fill out the rectangular hull. It turns out that in general a continuous function on RH cannot be approximated to arbitrary precision by  $\mathscr{PL}$ -interpolants. This is a corollary of a more general result to which we now turn.

In the next theorem we consider functions defined in a piecewise manner on Cartesian products. The setting will be as follows. There are two topological spaces given, S and T. Each is expressed as a union of non-empty sets

$$S = S_1 \cup \cdots \cup S_m$$
,  $T = T_1 \cup \cdots \cup T_n$ .

We assume that  $S_i \cap S_{i+1}$  and  $T_j \cap T_{j+1}$  are nonempty for  $1 \le i \le m-1$ ,  $1 \le j \le n-1$ . The notation C(S) denotes the space of continuous real-valued functins on a topological space S. Subspaces are prescribed as follows

$$G_i \subset C(S_i) \ (1 \leqslant i \leqslant m), \qquad H_j \subset C(T_j) \ (1 \leqslant j \leqslant n).$$

It is assumed that these subspaces contain the constant functions. Three spaces of continuous piecewise-defined functions (or "generalized splines") are given as follows, with | signifying restriction of a function to a subset of its domain:

$$G = \{ g \in C(S) : g \mid S_i \in G_i, \text{ all } i \}$$

$$H = \{ h \in C(T) : h \mid T_j \in H_j, \text{ all } j \}$$

$$K = \{ f \in C(S \times T) : f \mid (S_i \times T_j) \in G_i + H_i, \text{ all } i \text{ and } j \}.$$

5.1. THEOREM. The spaces G, H, and K defined above satisfy the equation K = G + H.

*Proof.* It is clear that if f(s, t) = g(s) + h(t), with  $g \in G$  and  $h \in H$ , then  $f \in K$ , since

$$f | (S_i \times T_i) = g | S_i + h | T_i \in G_i + H_i$$

Thus  $G + H \subset K$ .

Now let  $f \in K$ . We can find  $u_{ij} \in G_i$  and  $v_{ij} \in H_j$  so that

$$f|(S_i \times T_j) = u_{ij} + v_{ij} \qquad (1 \leqslant i \leqslant m, 1 \leqslant j \leqslant n).$$

Suppose j > 1. Select  $t_j \in T_j \cap T_{j-1}$ . Then for all  $s \in S_i$  we have  $(s, t_i) \in (S_i \times T_i) \cap (S_i \times T_{j-1})$ . Hence, for such a point,

$$u_{ij}(s) + v_{ij}(t_j) = u_{i,j-1}(s) + v_{i,j-1}(t_j).$$

This proves that  $u_{ij}(s) - u_{i,j-1}(s)$  is a constant,  $c_{ij}$ , on  $S_i$ . Using this equation repeatedly, we obtain

$$u_{ij} = u_{i,j-1} + c_{ij} = u_{i,j-2} + c_{i,j-1} + c_{ij} = \cdots = u_{i1} + \alpha_{ij}$$

where  $\alpha_{ij}$  is a constant. Because of the symmetry in the situation, we obtain a similar equaton for  $v_{ij}$ :

$$v_{ij} = v_{1j} + \beta_{ij}.$$

From a previous equation we have

$$u_{i1}(s) + \alpha_{ij} + v_{1j}(t_i) + \beta_{ij} = u_{i1}(s) + \alpha_{i,j-1} + v_{1,j-1}(t_i) + \beta_{i,j-1}.$$

Putting  $\gamma_{ij} = \alpha_{ij} + \beta_{ij}$ , we have

$$\gamma_{ij} - \gamma_{i,j-1} = v_{1,j-1}(t_j) - v_{1j}(t_j) = \delta_j.$$

Iterating this equation produces

$$\gamma_{ij} = \gamma_{i1} + \delta_j + \delta_{j-1} + \cdots + \delta_2 = \gamma_{i1} + d_j.$$

Thus on  $S_i \times T_j$  we have

$$f(s, t) = u_{i1}(s) + \alpha_{ij} + v_{1j}(t) + \beta_{ij}$$
  
=  $u_{i1}(s) + v_{1j}(t) + \gamma_{ij}$   
=  $[u_{i1}(s) + \gamma_{i1}] + [v_{1i}(t) + d_i].$ 

The first bracketed expression can be denoted by  $g_i(s)$ , where  $g_i \in G_i$ . The second can be denoted by  $h_i(t)$ , where  $h_i \in H_i$ . Our analysis shows that

$$f(s, t) = g_i(s) + h_j(t), \qquad (s, t) \in S_i \times T_j.$$

If  $s \in S_i \cap S_k$  then  $g_i(s) = g_k(s)$  because for any  $t \in T_j$ 

$$f(s, t) = g_i(s) + h_i(t) = g_k(s) + h_i(t)$$
.

This shows that there are unique functions g and h such that  $g \mid S_i = g_i$  and  $h \mid T_j = h_j$  for all i and j. Since f(s, t) = g(s) + h(t), the functions g and h are continuous.

5.2. COROLLARY. Assume the hypotheses made at the beginning of this section. Then

$$\{f \in C(S \times T): f \mid (S_i \times T_i) \in C(S_i) + C(T_i), \text{ all } i \text{ and } j\} = C(S) + C(T).$$

*Proof.* One inclusion is trivial, and the other is a consequence of the preceding theorem, letting  $G_i = C(S_i)$  and  $H_i = C(T_i)$ .

5.3. THEOREM. It is not possible to approximate with arbitrary precision all functions in  $C(S \times T)$  by use of functions which are piecewise of the form g(s) + h(t), no matter how finely we partition  $S \times T$  into Cartesian-product subsets.

*Proof.* The result follows from 5.2 and the fact that the subspace C(S) + C(T) is not dense in  $C(S \times T)$ . Indeed, it is annihilated by every path functional.

5.4. COROLLARY. There exists a function f in  $C(S \times T)$  such that  $dist(f, \mathcal{PL}) \ge 1$  for all grids.

### 6. THE SPACE RR

As in Section 1, a set of nodes

$$\mathcal{N} = \{x_1, x_2, ..., x_n\}$$

is given in  $\mathbb{R}^2$ , and we define radial basis functions

$$h_i(x) = ||x - x_i||_1 = |s - s_i| + |t - t_i|.$$

The norm symbol is henceforth reserved for the  $l_1$ -norm on  $\mathbb{R}^2$ .

We denote by  $\mathcal{R}\mathcal{B}$  the linear space generated by the radial basis functions  $h_i$   $(1 \le i \le n)$ . This section investigates the structure of  $\mathcal{R}\mathcal{B}$ . Since interpolation by  $h_1, ..., h_n$  at a set of n nodes certainly requires the linear independence of these functions, the first step is to characterize the sets  $\mathcal{N}$  for which  $\mathcal{R}\mathcal{B}$  is of dimension n.

6.1. LEMMA. Let  $\Phi$  be a function in  $C^1(\mathbb{R})$  that satisfies  $\Phi'(s) > 0$  for all s. Let  $\mathcal{N} = \{x_1, ..., x_n\} \subset \mathbb{R}^2$ , with  $\# \mathcal{N} = n$  and  $0 \notin \mathcal{N}$ . If the equation

$$\Phi(\|x\|_1) = \sum_{i=1}^n a_i \Phi(\|x - x_i\|_1)$$
 (1)

is valid in a neighborhood of 0, then  $0 \in P(\mathcal{N})$  and  $0 \in Q(\mathcal{N})$ .

*Proof.* Assume the hypotheses and deny the conclusion. With no loss of generality we suppose that  $0 \notin P(\mathcal{N})$ . Select  $\varepsilon > 0$  so that Eq. (1) is valid for  $||x|| < \varepsilon$ . If necessary, reduce  $\varepsilon$  so that  $(-\varepsilon, \varepsilon)$  contains no element of  $P(\mathcal{N})$ . If  $|s| < \varepsilon$  then Eq. (1) is valid for x = (s, 0), and thus

$$\Phi(|s|) = \sum_{i=1}^{n} a_i \Phi(|s-s_i| + |t_i|), \quad |s| < \varepsilon.$$

The function on the right in this equation is differentiable at s = 0, but the function on the left is not. This contradiction completes the proof.

6.2. LEMMA. Let  $\Phi$  be a function in  $C^1(\mathbb{R})$  that satisfies  $\Phi'(s) > 0$  and  $\Phi(0) = 0$ . Then any set of three functions  $H_i(s) = \Phi(\|x - x_i\|)$  is linearly independent on the corresponding set of three nodes. (In this result, any norm can be used.)

*Proof.* The value of the  $3 \times 3$  determinant  $det(H_i(x_j))$  is

$$2\Phi(\|x_1-x_2\|)\Phi(\|x_1-x_3\|)\Phi(\|x_2-x_3\|)\neq 0.$$

6.3. THEOREM. Let  $\mathcal{N} = \{x_1, ..., x_n\} \subset \mathbb{R}^2$ , with  $\# \mathcal{N} = n$ . Let  $\Phi \in C^1(\mathbb{R})$  and satisfy  $\Phi'(s) > 0$  and  $\Phi(0) = 0$ . Put  $H_i(x) = \Phi(\|x - x_i\|_1)$ . In order that the indexed set  $[H_1, ..., H_n]$  be linearly independent it is sufficient that  $\mathcal{N}$  contain no closed path. If  $\Phi$  is a linear function the condition is also necessary.

*Proof.* Assume that the set is dependent. With no loss of generality we suppose that it has no proper linearly dependent subset. By 6.2,  $n \ge 4$ . Let  $\sum_{i=1}^{n} c_i H_i = 0$  with  $\sum_{i=1}^{n} |c_i| > 0$ . Then  $c_i \ne 0$  for all i, and each  $H_i$  is a linear combination of the others. By 6.1, each node  $x_i$  has the property that the vertical line and the horizontal line through  $x_i$  each contains another node. It follows from 3.2 that  $\mathcal{N}$  contains a closed path.

To complete the proof, assume that  $\Phi$  is linear and that  $\mathcal{N}$  contains a closed path. If  $\Phi(r) = ar + b$ , then  $H_i(x) = a |s - s_i| + a |t - t_i| + b$ . The result now follows from 3.5.

6.4. LEMMA. Let  $x_1, ..., x_n$  be n distinct points in a normed space. Let  $f(x) = \sum_{i=1}^{n} a_i \|x - x_i\|$ . In order that f be bounded, it is necessary and sufficient that  $\sum_{i=1}^{n} a_i = 0$ .

**Proof.** Put 
$$I = \{i: a_i > 0\}, J = \{i: a_i < 0\}$$
. For all  $x$ ,

$$f(x) = \sum_{i} a_{i} \|x - x_{i}\| + \sum_{j} a_{i} \|x - x_{j}\|$$

$$\leq \sum_{i} a_{i} (\|x\| + \|x_{i}\|) + \sum_{i} a_{i} (\|x\| - \|x_{i}\|) = \|x\| \sum_{i} a_{i} + c,$$

where  $c = \sum |a_i| ||x_i||$ . Similarly,

$$f(x) \geqslant ||x|| \sum a_i - c.$$

## 7. Interpolation by RB

This section contains the central result of the paper. It is shown that the interpolation problem

$$\sum_{j=1}^{n} a_{j} \| x_{i} - x_{j} \|_{1} = d_{i} \qquad (1 \le i \le n)$$

has a unique solution for every data function d if and only if the set of nodes  $\{x_1, x_2, ..., x_n\}$  contains no closed path. The work of Micchelli [9] enables us to generalize this problem. We assume throughout this section that F is a function fulfilling five requirements:

- (M1)  $F: [0, \infty) \to [0, \infty),$
- (M2) F is  $C^{\infty}$  on  $(0, \infty)$  and continuous at 0,
- (M3) F(t) > 0 when t > 0,
- (M4) F' is not constant,
- (M5)  $(-1)^{\nu} F^{(\nu+1)}(t) \ge 0$  for  $\nu = 0, 1, 2, ..., \text{ and } t > 0$ .

Suppose that G is a function having the same properties. We consider interpolation at the nodes by a linear combination of these functions:

$$H_i(x) = F((s-s_i)^2) + G((t-t_i)^2)$$
  $(1 \le j \le n)$ .

In this equation, x = (s, t) and  $x_j = (s_j, t_j)$ . The coefficient matrix A that arises in this more general interpolation problem is given by

$$A = B + C$$
 where  $B_{ij} = F((s_i - s_j)^2)$  and  $C_{ij} = G((t_i - t_j)^2)$ .

If F and G are chosen to be the square-root function, then we recover the original problem. Micchelli [9] establishes the following important theorem.

7.1. THEOREM (Micchelli). If F satisfies the five conditions (M1-M5) given above and if  $r_1, ..., r_p$  are distinct reals, then the  $p \times p$  matrix  $D_{ij} = F((r_i - r_j)^2)$  is nonsingular. Also,  $c^T D c < 0$  for every nonzero vector c such that  $\sum_{j=1}^{p} c_j = 0$ .

From Micchelli's theorem, we see immediately that the rank of B is m. Indeed, we can remove from B rows and columns which are duplicates of other rows and columns, arriving at an  $m \times m$  matrix B' whose elements are  $F((\sigma_i - \sigma_i)^2)$ . This matrix is nonsingular, by Micchelli's theorem.

Our next task is to describe a basis for  $\ker(B)$ . If  $1 \le j < i \le n$ , then  $f^{ij}$  will denote a vector in  $\mathbb{R}^n$  having 1 as its *i*th component, -1 as its *j*th component, and 0 components elsewhere. Thus,  $f^{ij}_{\mu} = \delta_{i\mu} - \delta_{j\mu}$ . Define also

$$J = \{(i, j): 1 \le j < i \le n, s_i = s_i, s_u \ne s_i \text{ if } \mu < j\}.$$

Notice that J is a function: different elements of J cannot have the same first component.

7.2. LEMMA. A basis for  $\ker(B)$  is  $\{f^{ij}: (i, j) \in J\}$ .

*Proof.* First we establish that the purported basis is a subset of ker(B). If  $(i, j) \in J$ , then  $Bf^{ij} = 0$  because

$$(Bf^{ij})_{\nu} = \sum_{\mu=1}^{n} B_{\nu\mu} f^{ij}_{\mu} = B_{\nu i} - B_{\nu j} = F((s_{\nu} - s_{i})^{2}) - F((s_{\nu} - s_{j})^{2}) = 0.$$

Next we prove that the purported basis is linearly independent. Suppose that  $\sum \alpha_{ij} f^{ij} = 0$ , the sum being over  $(i, j) \in J$ . Select any  $(v, \mu) \in J$ . We shall show that  $\alpha_{v\mu} = 0$ . This follows from the calculation

$$0 = \sum \alpha_{ij} f_{\nu}^{ij} = \alpha_{\nu\mu} f_{\nu}^{\nu\mu} = \alpha_{\nu\mu}.$$

To justify this, we only have to prove that if  $f_{\nu}^{ij} \neq 0$  then  $(i, j) = (\nu, \mu)$ . If  $f_{\nu}^{ij} \neq 0$ , then either  $\nu = i$  or  $\nu = j$ . If  $\nu = i$ , then  $(i, j) \in J$  and  $(i, \mu) \in J$ . Since J is a function,  $j = \mu$ . If  $\nu = j$  we have  $(i, j) \in J$ ,  $(j, \mu) \in J$ ,  $\mu < j < i$ , and  $s_i = s_j = s_\mu$ , which contradicts  $(i, j) \in J$ . (This case can therefore not arise.) Lastly, we observe that J has the correct cardinality, namely n - m,

$$n = \# \mathcal{N} = \# P(\mathcal{N}) + \sum_{i=1}^{m} [\# P^{-1}(\sigma_i) - 1] = m + \# J.$$

which is the dimension of ker(B). This assertion follows from the equation

7.3. LEMMA. If  $u \in \ker(B)$  then  $\sum_{i=1}^{n} u_i = 0$  and

$$\sum \{u_i : s_i = \sigma_i\} = 0 \qquad (1 \le j \le m).$$

*Proof.* The first equation follows from the second by summing for  $1 \le j \le m$  (or it can be proved directly for the basis vectors in 7.2). In proving the second equation, it suffices to verify it for any one of the basis vectors in  $\ker(B)$  as described above. To this end, fix  $(\mu, \nu) \in J$  and  $j \in \{1, ..., m\}$ . By the definitions of J and  $f^{\mu\nu}$ ,

$$\sum \{f_i^{\mu\nu} : s_i = \sigma_j\} = \sum \{\delta_{i\mu} - \delta_{i\nu} : s_i = \sigma_j\}.$$

This is obviously zero unless  $s_{\mu} = \sigma_j$  or  $s_{\nu} = \sigma_j$ . But these equations imply each other, and if  $s_{\mu} = \sigma_j = s_{\nu}$ , the sum in question reduces to 1 - 1 = 0.

7.4. LEMMA. If 
$$v^T B v = \sum_{i=1}^n v_i = 0$$
, then  $B v = 0$ .

*Proof.* For v=1, 2, ..., m put  $I_v = \{i: 1 \le i \le n, s_i = \sigma_v\}$ . The sets  $I_1, ..., I_m$  form a partition of  $\{1, 2, ..., n\}$ . Hence any sum of the form  $\sum_{i=1}^n$  can be expressed as a double sum  $\sum_{v=1}^m \sum_{i \in I_v}$ . We observe also that if  $i \in I_v$  and  $j \in I_u$  then

$$B_{ii} = F((s_i - s_i)^2) = F((\sigma_v - \sigma_u)^2) \equiv B'_{vu}.$$

Now assume the hypotheses and put  $v'_{\nu} = \sum_{i \in I_{\nu}} v_i$ . Applying the above principles we have

$$0 = v^{T}Bv = \sum_{i=1}^{n} \sum_{j=1}^{n} B_{ij}v_{i}v_{j} = \sum_{v=1}^{m} \sum_{i \in I_{v}} \sum_{\mu=1}^{m} \sum_{j \in I_{\mu}} B'_{\nu\mu}v_{i}v_{j}$$
$$= \sum_{v=1}^{m} \sum_{\mu=1}^{m} B'_{\nu\mu}v'_{\nu}v'_{\mu} = (v')^{T} B'v'.$$

Notice also that  $\sum_{\nu=1}^{m} v'_{\nu} = \sum_{i=1}^{n} v_{i} = 0$ . Since the points  $\sigma_{\nu}$  are distinct, the matrix B' has the properties in Micchelli's theorem (7.1). Hence by 7.1, v' = 0. It follows that Bv = 0 by the calculation

$$(Bv)_{i} = \sum_{j=1}^{n} B_{ij} v_{j} = \sum_{\mu=1}^{m} \sum_{j \in I_{\mu}} B_{ij} v_{j} = \sum_{\mu=1}^{m} \sum_{j \in I_{\mu}} F((s_{i} - \sigma_{\mu})^{2} v_{j})$$

$$= \sum_{\mu=1}^{m} F((s_{i} - \sigma_{\mu})^{2}) v'_{\mu} = 0. \quad \blacksquare$$

7.5. LEMMA. If  $u \in \ker(B)$ , then every vertical line that intersects the set  $\Gamma(u) = \{x_i \in \mathcal{N} : u_i \neq 0\}$  contains at least two points of  $\Gamma(u)$ .

*Proof.* Let  $x_i \in \Gamma(u)$  so that  $u_i \neq 0$ . By 7.3

$$\sum \{u_i: s_i = s_i\} = 0.$$

Thus there must exist at least one index j, different from i, for which  $u_j \neq 0$  and  $s_i = s_i$ . Then  $x_i$  is an element of  $\Gamma(u)$  on the vertical line through  $x_i$ .

7.6. LEMMA. If  $\ker(B) \cap \ker(C) \neq 0$  then  $\mathcal{N}$  contains a closed path.

*Proof.* Suppose that u is a nonzero vector in  $\ker(B) \cap \ker(C)$ . Then  $\Gamma(u)$  is nonvoid. By 7.5, every vertical line that intersects  $\Gamma(u)$  contains at least two points of  $\Gamma(u)$ . Applying the same lemmas to C shows that every horizontal line that intersects  $\Gamma(u)$  contains two points of  $\Gamma(u)$ . By 3.2,  $\Gamma(u)$  contains a path and, a fortiori, so does  $\mathcal{N}$ .

7.7. THEOREM. Let  $\mathcal{N}$  be a set of n distinct points  $x_i = (s_i, t_i) \in \mathbb{R}^2$ . Let F and G be functions satisfying hypotheses (M1)–(M5) above. The  $n \times n$  matrix A defined by

$$A_{ij} = F((s_i - s_j)^2) + G((t_i - t_j)^2),$$

is singular if and only if N contains a closed path.

*Proof.* If  $\mathcal{N}$  contains a closed path, then the functions

$$g_i(s, t) = F((s-s_i)^2) + G((t-t_i)^2),$$

form a dependent set (by 3.5), and thus A is singular.

Now let v be any vector such that  $v \neq 0$  and  $v^T e = 0$ , where  $e = (1, 1, ..., 1)^T$ . If the points  $s_1, s_2, ..., s_n$  were distinct, then 7.1 would imply  $v^T B v < 0$ . Since the points  $s_i$  are not necessarily distinct, a limit argument yields  $v^T B v \leq 0$ . Similarly  $v^T C v \leq 0$ . Hence

$$v^T A v = v^T B v + v^T C v \leqslant 0.$$

Since A is symmetric, its eigenvalues are real and can be ordered  $\lambda_1 \leq \lambda_2 \leq \cdots \leq \lambda_n$ . Since  $e^T A e > 0$ , it follows from the Courant-Fischer ("minimax") theorem that  $\lambda_n > 0$ . Using this result again, we have

$$\lambda_{n-1} = \min_{\dim V = n-1} \max_{\substack{v \in V \\ \|v\| = 1}} v^T A v \leqslant \max_{\substack{v^T e = 0 \\ \|v\| = 1}} v^T A v \leqslant 0.$$

Now assume that A is singular. Then  $\lambda_{n-1} = 0$ , and hence there is a vector v satisfying ||v|| = 1,  $v^T e = 0$ ,  $v^T A v = 0$ . It follows that  $v^T B v = v^T C v = 0$ . By 7.4,  $v \in \ker(B) \cap \ker(C)$ , and by 7.6,  $\mathcal{N}$  contains a closed path.

7.8. COROLLARY. Let  $\mathcal{N}$  be a set of n distinct points  $x_i = (s_i, t_i)$  in  $\mathbb{R}^2$ . Let  $0 < \alpha < 2$ . The  $n \times n$  matrix A given by

$$A_{ij} = |s_i - s_j|^{\alpha} + |t_i - t_j|^{\alpha}$$

is singular if and only if  $\mathcal N$  contains a closed path.

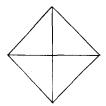
**Proof.** In 7.7, let 
$$F(u) = G(u) = u^{\alpha/2}$$
.

The results of this section are extended in [3] to interpolation by sums of radial functions.

### 8. RADIAL BASIS FUNCTIONS WITH THE MAXIMUM NORM

All of what has been proved for radial basis functions with the  $l_1$ -norm can be proved, *mutatis mutandis*, for the  $l_{\infty}$ -norm. This assertion depends upon the isometry between  $l_1^{(2)}$  and  $l_{\infty}^{(2)}$  that must exist because of the similarity in the unit spheres in these two spaces. (See the figure.)

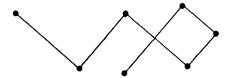




The isometry from  $l_1^{(2)}$  to  $l_{\infty}^{(2)}$  is given by  $(s, t) \to (s + t, s - t)$ . Given  $\mathcal{N} = \{x_1, x_2, ..., x_n\} \subset \mathbb{R}^2$ , our basis functions are now

$$h_j(x) = ||x - x_j||_{\infty} = \max[|s - s_j|, |t - t_j|].$$

The notion of a path must now be modified; we refer to the new concept as an  $l_{\infty}$ -path. It is an ordered set of points  $[z_1, z_2, ..., z_l]$  such that the line segments joining successive points are of positive length and have inclinations alternately 45° and 135°.



The grid lines generated by  $\mathcal{N}$  consist now of lines with inclinations 45° and 135° through points of  $\mathcal{N}$ .

8.1. THEOREM. Let  $\mathcal{N}$  be a set of n points,  $x_i$ , in  $\mathbb{R}^2$ . The functions

 $x \mapsto \|x - x_i\|_{\infty}$   $(1 \le i \le n)$  are capable of interpolating arbitrary data on  $\mathcal{N}$  if and only if  $\mathcal{N}$  contains no closed  $l_{\infty}$ -path.

The preceding considerations provide an example of the following general principle. If functions  $f_1, ..., f_n$  are capable of interpolating arbitrary data at nodes  $x_1, ..., x_m$ , and if L is a nonsingular linear transformation, then the functions  $f_1 \circ L^{-1}, ..., f_n \circ L^{-1}$  are capable of interpolating arbitrary data at nodes  $Lx_1, ..., Lx_m$ .

# 9. GENERALIZATIONS TO HIGHER-DIMENSIONAL SPACES

The basic interpolation results of Section 7 can be generalized to the space  $\mathbb{R}^d$ ,  $d \ge 2$ . To describe the results, a somewhat different formalism from that used in the previous sections is needed.

If x is a point in  $\mathbb{R}^d$ , we write  $x = (\xi_1, \xi_2, ..., \xi_d)$ . Coordinate functionals  $p_i$  are defined by setting  $p_i(x) = \xi_i$ . A function  $f: \mathbb{R}^d \to \mathbb{R}$  is said to be degenerate if it has the form  $f = \sum_i^d g_i \circ p_i$  for suitable  $g_i: \mathbb{R} \to \mathbb{R}$ . Such a function f is a sum of univariate functions. The space of all degenerate functions on  $\mathbb{R}^d$  is denoted by  $\mathcal{D}$ .

If f is a function on a finite set  $S = \{s_1, s_2, ..., s_n\}$ , we write

$$\sum [f(s): s \in S] = f(s_1) + f(s_2) + \dots + f(s_n).$$

Note that if f is not injective, the sum will contain repeated terms. That is why we eschew the notation  $\sum \{f(s): s \in S\}$ , which—if strictly interpreted—means a sum without repetitions.

9.1. LEMMA. In order that a functional of the form  $\phi = \sum_{i=1}^{N} c_i \hat{y}_i$  annihilate  $\mathcal{D}$  it is necessary and sufficient that for all  $t \in \mathbb{R}$  and for all  $v \in \{1, 2, ..., d\}$ ,

$$\sum [c_i: p_v(y_i) = t] = 0.$$

*Proof.* Let  $g: \mathbb{R} \to \mathbb{R}$ , and fix v. Let

$$A = \{ p_{\nu}(y_i) \colon 1 \leqslant i \leqslant N \}.$$

We have then

$$\phi(g \circ p_{\nu}) = \sum_{i=1}^{N} c_{i} g(p_{\nu}(y_{i})) = \sum_{t \in A} \sum [c_{i} g(p_{\nu}(y_{i})) : p_{\nu}(y_{i}) = t]$$

$$= \sum_{t \in A} g(t) \sum [c_{i} : p_{\nu}(y_{i}) = t].$$

The sufficiency of the given condition is now clear. For the necessity, select  $t \in A$ , and construct g so that g(t) = 1 and g(s) = 0 for all  $s \in A \setminus \{t\}$ . The preceding calculation gives us

$$0 = \phi(g \circ p_v) = \sum [c_i : p_v(y_i) = t]. \quad \blacksquare$$

If  $v_1, v_2, ..., v_n$  are elements in a vector space, we adopt the usual meaning for linear dependence of the *indexed* set  $[v_1, v_2, ..., v_n]$ . It can happen that the *unindexed* set  $\{v_1, v_2, ..., v_n\}$  is linearly independent while  $[v_1, v_2, ..., v_n]$  is linearly dependent.

- 9.2. Lemma. The following properties of a set  $\mathcal{N} = \{x_1, x_2, ..., x_n\}$  in  $\mathbb{R}^d$  are equivalent:
  - (a) There is a nonzero functional in  $\mathcal{D}^{\perp}$  that is supported on  $\mathcal{N}$ .
- (b) For every  $f \in \mathcal{D}$ , the indexed set of translates  $x \mapsto f(x x_i)$  is linearly dependent.

*Proof.* Let  $E_u$  denote the translation operator, defined by  $(E_uF)(x) = F(x-u)$ . Define an operator B by putting (BF)(x) = F(-x). Let  $\phi = \sum_{i=1}^{n} c_i \hat{x}_i$ , and assume that  $\phi \in \mathcal{D}^{\perp}$ . If  $f \in \mathcal{D}$ , then  $E_uBf \in \mathcal{D}$ , and consequently

$$0 = \phi(E_u B f) = \sum_i c_i (E_u B f)(x_i) = \sum_i c_i f(u - x_i) = \sum_i c_i (E_{x_i} f)(u).$$

This proves that (a) implies (b). Observe that the proof requires of  $\mathcal{D}$  only its invariance under the operators B and  $E_u$ . For the other half of the proof, assume (b). Let  $f \in \mathcal{D}$ . Then  $Bf \in \mathcal{D}$ , and by (b) there exist coefficients  $c_i$ , not all zero, such that  $\sum c_i E_{x_i} Bf = 0$ . Evaluating at 0, we have

$$0 = \sum c_i (E_{x_i} B f)(0) = \sum c_i f(x_i) = \sum c_i \hat{x}_i(f). \quad \blacksquare$$

Now select functions  $F_1$ ,  $F_2$ , ...,  $F_d$  satisfying the five axioms (M1)–(M5) of Section 7. Define

$$H(x) = \sum_{\nu=1}^{d} F_{\nu}((p_{\nu}(x)^2), \qquad x \in \mathbb{R}^d.$$

As before a set of nodes is given:  $\mathcal{N} = \{x_1, x_2, ..., x_n\}$ , with  $x_i \in \mathbb{R}^d$ . Interpolation at the nodes by the  $x_i$ -translates of H requires the nonsingularity of the interpolation matrix A given by

$$A_{ij} = H(x_i - x_j) \qquad (1 \le i, j \le n).$$

It is clear from the definition of H that A is the sum of matrices  $A^{(v)}$  given by

$$A_{ii}^{(v)} = F_{v} [(p_{v}(x_{i}) - p_{v}(x_{i}))^{2} \qquad (1 \le v \le d).$$

9.3. LEMMA. If the matrix A is singular, then there exists a vector  $u \in \mathbb{R}^n$  such that  $u \neq 0$ ,  $u^T e = 0$ , and  $A^{(v)}u = 0$  for  $1 \leq v \leq d$ .

*Proof.* Proceed exactly as in the proof of 7.7, obtaining thereby a vector u having the desired properties.

- 9.4. Theorem. The following are equivalent properties of the node set  $\mathcal{N}$ :
  - (a) The interpolation matrix A is singular;
- (b) There is a linear dependence among the n basis functions  $x \mapsto H(x-x_i)$ .

*Proof.* That (b) implies (a) is obvious. Assume that (a) is true. By the preceding lemma, there is a nonzero vector u such that  $A^{(v)}u = 0$  for  $1 \le v \le d$ . By 7.3, we have

$$\sum [u_i: p_v(x_i) = t] = 0 \qquad (t \in \mathbb{R}, 1 \le v \le d).$$

By Lemma 9.1, the functional  $\sum u_i \hat{x}_i$  annihilates  $\mathcal{D}$ . By Lemma 9.2, the set of functions  $E_{x_i}H$  is linearly dependent.

The geometrical characteristics of  $\mathcal{N}$  that are equivalent to Properties (a) and (b) in 9.4 will be explored in the second half of this paper.

Notice that if d>2, the theory of radial basis functions using the  $l_{\infty}$ -norm is not a simple consequence of the theory in the case of the  $l_1$ -norm. This is because there is no isometric isomorphism between the spaces  $\mathbb{R}^d$  when these two norms are used. For example, the unit balls in  $\mathbb{R}^3$  are a cube and an octahedron for these two norms.

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