Exterior differentiation in the Regge calculus

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Regge manifolds are piecewise continuous manifolds constructed from a finite number of basic building blocks. On such manifolds piecewise continuous forms can be defined in a way similar to differential forms on a differentiable manifold. Regge manifolds are used extensively in the construction of space-times in numerical general relativity. In this paper a definition of exterior differentiation suitable for use on piecewise continuous forms on a Regge manifold is presented. It is shown that this definition leads to a version of Stokes' theorem and also to the usual result that $d^2 = 0$. This is preceded by a discussion of certain geometrical properties of the Regge manifolds. It is shown that the version of Stokes' theorem presented here coincides with the usual definition when the Regge manifold is refined, by increasing the number of cells while keeping the total volume constant, to a smooth manifold.

I. INTRODUCTION

It is assumed that solutions of Regge's field equations¹⁻⁵, Regge space-times, are approximations, to a degree, of an Einstein space-time, this being a differentiable solution of Einstein's field equations. This assumption is based on two facts. First, the Regge and Einstein manifolds are equivalent under a homeomorphism. Second, both sets of field equations are derived from the same action principle. It is therefore not unreasonable to expect that there should exist a correspondence between certain properties of the Regge and Einstein space-times. In particular the operation of exterior differentiation on an Einstein manifold should lead to a related operation on a Regge manifold.

The main result to be presented here is an operation on forms built on Regge manifolds which mimics the usual operation of exterior differentiation. This result⁶ is presented in Sec. IV. In Secs. II and III the basic notation and formulas are presented. Finally, in Sec. V, it is shown that this definition reduces to the version usually employed on smooth manifolds.

II. SIMPLEXES AND COMPLEXES

The fundamental building blocks for the manifolds to be considered here are known as *n*-simplexes. They may be defined in a recursive fashion as follows.

- (i) A 0-simplex is a single point. This object is also called a vertex.
- (ii) An (n + 1)-simplex is constructed from a n-simplex by first introducing one new vertex and then joining this vertex to each of the (n + 1) vertices of the n-simplex, and second by demanding that any set of m vertices $(1 \le m \le n + 2)$ of the (n + 1)-simplex is an (m 1)-simplex.
- (iii) A Lorentzian n-simplex is obtained by imposing a flat Lorentzian metric throughout the n-simplex. Only those n-simplexes in which the induced metric on each of its m-simplexes (0 < m < n) is also flat will be considered. This has the effect of disallowing any n-simplexes with curved boundaries.

The obvious method of constructing an n-dimensional manifold is to glue together a collection of n-simplexes. The resulting object is referred to as an n-complex. To avoid certain pathological cases the following restrictions are imposed: (i) the region of the n-complex common to two or more adjacent n-simplexes is an m-simplex with $0 \le m < n$, and (ii) any m-simplex of the complex is contained within at least one n-simplex of the n-complex.

The following notation is drawn, primarily, from Seirfert and Threlfall.⁷

A typical *n*-simplex is denoted by $\sigma_n(i)$, with the index *i* being the label which distinguishes this simplex from all other *n*-simplexes. The set of all *n*-simplexes is represented by S_n . Each *n*-simplex contains exactly (n+1) vertices and is represented as follows:

$$\sigma_n(i) = (i_0 i_1 i_2 \cdots i_n), \quad n = 0, 1, 2, \dots,$$
 (2.1)

where each of the i_j is unique and is the label of a vertex of $\sigma_n(i)$. The order in which the vertices are listed is unimportant unless the simplex is oriented. One of the two possible orientations to the simplex is defined by reading the vertices in (2.1) from left to right. The opposite sense of orientation is obtained if any two vertices in the sequence are swapped. This is indicated by writing

$$(i_0 \cdots i_j \cdots i_k \cdots i_n) = -(i_0 \cdots i_k \cdots i_j \cdots i_n),$$
provided that $j \neq k$ and $j, k = 0, 1, 2, ..., n$. (2.2)

An *n*-complex is denoted by ρ_n and represented by the formal sum

$$\rho_n = \sum_i a_i \sigma_n(i), \tag{2.3}$$

with each $a_i = 0$, -1, or +1. The coefficients a_i represent whether the associated *n*-simplex is present or not and what orientation it possesses in the complex. Complexes in which certain *n*-simplexes are absent are referred to as sub- or secondary complexes. The original complex, when required, is referred to as the primary complex.

One particularly important subcomplex is the boundary of the primary complex. It is defined as follows. First suppose that $\sigma_n(i)$ is represented as in (2.1). Then define the operation

$$\sigma_n(i/i_j) = (i_0 i_1 \cdots i_n/i_j) = \begin{cases} (-1)^j (i_0 i_1 \cdots \hat{i_j} \cdots i_n), \\ 0, & \text{if } i_j \text{ is not in } i_0, i_1, \dots, i_n. \end{cases} (2.4)$$

The symbol over i_j indicates that the vertex i_j is excluded from the list. The boundary of an *n*-simplex $\sigma_n(i)$ is defined as

$$\partial \sigma_n(i) = \sum_{i_i} \sigma_n(i/i_j). \tag{2.5}$$

Similarly for an n-complex

$$\rho_n = \sum_i a_i \sigma_n(i),$$

$$\partial \rho_n = \sum_{i,j} a_i \sigma_n(i/j).$$
(2.6)

Obviously $\sigma_n(i/j)$ is also a simplex and thus the operation (2.5) may be applied twice. However from (2.2) and (2.4) it is clear that

$$\sigma_n(i/j/k) = -\sigma_n(i/k/j). \tag{2.7}$$

This leads to the usual result⁷

$$\partial^2 \sigma_n(i) = 0. (2.8)$$

III. THE METRIC FRAME OF A SIMPLEX

A. The natural frame

One of the easiest ways of ensuring that the metric of a simplex is flat is to demand that all of the metric coefficients are constant throughout the simplex. Of course there are other frames in which the coefficients are not constant and yet the metric is flat. For simplicity such frames will be ignored.

It will be convenient to distinguish between the terms "coordinate frame" and "metric frame." The term "coordinate frame" will be used to refer to a frame possessing coordinates but not a metric. In a "metric frame" there are both coordinates and a metric. A very useful metric frame, the natural frame, will now be described.

Choose one *n*-simplex, $\sigma_n(1)$, and label its vertices from 0 to *n*. Adopt the vertex (0) as the origin of the coordinate frame. The basis vectors for this *n*-simplex are chosen as the *n* vectors that join (0) to the remaining vertices, i.e., along those legs connected to the origin. The coordinate frame is chosen such that the coordinates of the vertices are

$$\begin{split} \sigma_0(0) &= (0) & : \quad (v_0^\mu) = (0,0,\dots,0), \\ \sigma_0(1) &= (1) & : \quad (v_1^\mu) = (1,0,\dots,0), \\ \sigma_0(2) &= (2) & : \quad (v_2^\mu) = (0,1,\dots,0), \\ & \vdots & \\ \sigma_0(n) &= (n) & : \quad (v_n^\mu) = (0,0,\dots,1). \end{split}$$

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$$v_n^{\mu} = \begin{cases} 0, & \text{if } n = 0, & \mu = 1, 2, \dots, n, \\ \delta_n^{\mu}, & \text{if } n > 0, & \mu = 1, 2, \dots, n. \end{cases}$$
 (3.1)

Denote the basis vectors by \mathbf{e}_{μ} . Then any point P in the simplex is described by the vector

$$\mathbf{P} = x^{\mu} \mathbf{e}_{\mu}$$
, with each $x^{\mu} \geqslant 0$ and $\sum_{\mu} x^{\mu} \leqslant 1$. (3.2)

The requirement that $\Sigma_{\mu} x^{\mu} \leq 1$ ensures that the vector does not pass through the face opposite the origin. This completes the construction of a coordinate frame for this simplex. A metric frame will now be constructed by introducing the leg lengths L_{ij} and the metric components $g_{\mu\nu}$.

Denote the proper distance between the vertices (i) and (j) by L_{ij} . If all the $g_{\mu\nu}$ are known then, using (3.1), the L_{ij} would be computed as

$$L_{0i}^{2} = g_{ii}, 1 \le i \le n, L_{ij}^{2} = (g_{ii} - 2g_{ij} + g_{jj}), 1 \le i, j \le n \text{ and } i \ne j.$$
 (3.3)

In this instance there is no summation over repeated indices. Solve these equations for the $g_{\mu\nu}$ to obtain

$$g_{ii} = L_{01}^2,$$
 $1 \le i \le n,$ $g_{ij} = (L_{0i}^2 + L_{0j}^2 - L_{ij}^2)/2,$ $1 \le i, j \le n \text{ and } i \ne j.$ (3.4)

So far nothing has been said about the signature of the metric. For a physically realistic space-time, in the sense of general relativity, the signature must be Lorentzian (i.e., -+++). Thus it is clear that not all of the leg lengths can be specified without restriction. However, this restriction is somewhat weak for it is possible, in all but a few exceptions, to make small arbitrary changes in the L_{ij} and yet not change the signature. Therefore assume that the signature is Lorentzian for each and every n-simplex of the primary complex. In a later paper a technique of constructing an n-simplex will be described by which the signature may be guaranteed to be Lorentzian.

A knowledge of the $g_{\mu\nu}$ also enables the computation of areas and volumes of *n*-simplexes. Define the measure of an *n*-simplex, σ_n , as the *n*-fold integral

$$M(\sigma_n) = \int_{n-\text{simplex}} {\binom{n}{g}}^{1/2} d^n x, \qquad (3.5)$$

where "g is the determinant of the $g_{\mu\nu}$. The limits of integration are easily deduced from (3.1). For example, for a four-simplex, the four-dimensional measure is

$$M(\sigma_4) = ({}^4g)^{1/2} \times \int_0^1 \int_0^{1-x^1} \int_0^{1-x^1-x^2} \int_0^{1-x^1-x^2-x^3} dx^4 dx^3 dx^2 dx^1,$$

since 4g is a constant. In this instance the repeated integral has the value $\frac{1}{4}!$, in the general case of (3.5) the value is 1/n!. Since the signature of the metric may be indefinite it is possible to obtain an imaginary value for $M(\sigma_n)$. This is an unnecessary complication and will be avoided by using the absolute value of ng in (3.5). Thus the measure of an n-simplex, σ_n , is

$$M(\sigma_n) = (1/n!) (abs(^ng))^{1/2}.$$
 (3.6)

B. The general frame

On occasions it may be useful to employ a frame other than the natural frame. For example, if a study of the properties of a group of simplexes is to be made then it may be necessary to build a metric frame covering all simplexes of the group. Clearly the natural frame is inappropriate in this example. The construction of a general class of metric frames will now be discussed. Once again assume that the $g_{\mu\nu}$ are constant throughout each n-simplex and that the basis vectors are chosen as $\partial /\partial x^{\mu}$.

Consider a complex with one or more n-simplexes. To each vertex, $\sigma_0(i)$, of this complex, associate n coordinates, $x^{\mu}(i)$. Provided the topology of this complex is not too peculiar, it should be possible to choose these coordinates so that, for each and every n-simplex, the coordinates constitute a coordinate frame for that simplex. This condition simply ensures that locally (i.e., within one n-simplex) the coordinate frame is n-dimensional. Assume that such a choice can and has been made.

The components of the vector joining $\sigma_0(j)$ to $\sigma_0(i)$, denoted by $L^{\mu}(ij)$, have the values

$$L^{\mu}(ij) = x^{\mu}(i) - x^{\mu}(j). \tag{3.7}$$

The values of the $g_{\mu\nu}$ are obtained by solving the equation

$$L_{ij}^{2} = g_{\mu\nu}L^{\mu}(ij)L^{\nu}(ij). \tag{3.8}$$

Since there are n(n + 1)/2 leg lengths in each n-simplex and a similar number of $g_{\mu\nu}$'s there may exist a unique solution of (3.8). That a unique solution does exist is guaranteed by the earlier requirement that the coordinate frame be everywhere n-dimensional.

Notice that the values of the $g_{\mu\nu}$ need not be the same in each n-simplex. Thus there may be discontinuous changes in the $g_{\mu\nu}$ across the interfaces between pairs of *n*-simplexes. Consequently there results a possible ambiguity in the process of raising and lowering indices. For example suppose the leg (i j) is common to two *n*-simplexes. Then the values of $L^{\mu}(i j)$ with the index lowered may depend upon the choice of simplex in which the computation was performed. There is of course no ambiguity in the $L^{\mu}(ij)$. It would therefore be inaccurate to write $L_{\mu}(ij)$ as the lowered version of $L^{\mu}(ij)$, however, in most applications it will be clear which n-simplex is intended.

Consider one *m*-simplex $\sigma_m = (i_0 i_1 ... i_m)$ in an (m + 1)complex and now define the following quantities

$$L_{j}^{\mu} = L^{\mu}(i_{0}i_{j}), \quad \text{for } j = 1, 2, \dots, m,$$

$$L^{\mu_{1}\mu_{2}\cdots\mu_{m}}(\sigma_{m}) = e_{12\cdots m}^{\nu_{1}\nu_{2}\cdots\nu_{m}}L_{\nu_{1}}^{\mu_{1}}L_{\nu_{2}}^{\mu_{2}}\cdots L_{\nu_{m}}^{\mu_{m}},$$
(3.9)

with $e_{ij...k}^{\mu\nu...\rho} = +1, -1, 0$ when $(\mu\nu...\rho)$ is either an even, odd, or a nonpermutation of $(ij \cdots k)$, respectively. For example, for m = 3,

$$\begin{split} L^{\mu\nu\rho}(\sigma_3) &= L^{\mu}_{1}L^{\nu}_{2}L^{\rho}_{3} - L^{\nu}_{1}L^{\mu}_{2}L^{\rho}_{3} + L^{\nu}_{1}L^{\rho}_{2}L^{\mu}_{3} \\ &- L^{\rho}_{1}L^{\nu}_{2}L^{\mu}_{3} + L^{\rho}_{1}L^{\mu}_{2}L^{\nu}_{3} - L^{\mu}_{1}L^{\rho}_{2}L^{\nu}_{3}. \end{split}$$

The following expressions, shown only for m = 3 but easily generalized, are all derived from the definition (3.9):

$$L^{\mu\nu\rho}(\sigma_3) = e_{123}^{ijk} L^{\mu}_{i} L^{\nu}_{i} L^{\rho}_{k}, \qquad (3.10a)$$

$$L^{\mu\nu\rho}(\sigma_3) = e_{123}^{ijk} N_i^{\mu} L_i^{\nu} L_i^{\rho}, \qquad (3.10b)$$

$$L^{\mu\nu\rho}(\sigma_3) = e^{\mu\nu\rho}_{\alpha\beta\gamma}L^{\alpha}(\sigma_1)L^{\beta\gamma}(\sigma_2), \qquad (3.10c)$$

$$L^{\mu\nu\rho}(\sigma_3) = e^{\mu\nu\rho}_{\alpha\beta\gamma} N^{\alpha}(\sigma_1) L^{\beta\gamma}(\sigma_2), \tag{3.10d}$$

with σ_2 = one face of σ_3 , σ_1 = one leg of σ_3 but not of σ_2 , $N^{\alpha}(\sigma_2)$ = the projection of $L^{\alpha}(\sigma_1)$ onto the normal to σ_2 , and N_i^{α} = the projection of L_i^{α} onto the normal of its adjacent

Now let $L(\sigma_n)$ = the measure of the parallel *n*-cube formed from the n legs

$$(i_0i_1),...,(i_0i_n)$$
 of $\sigma_n = (i_0i_1...i_n)$.

It is well known that

$$n!L^{2}(\sigma_{n}) = L_{(\mu)}(\sigma_{n})L^{(\mu)}(\sigma_{n}), \text{ with } (\mu) = (\mu_{1}\mu_{2}\cdots\mu_{n}).$$
(3.11)

Alternatively $L(\sigma_n)$ can be computed by an integration like that in (3.5). In this case the limits of integration must now be chosen to cover an *n*-cube rather than an *n*-simplex. The result of this integration is similar to (3.6) with the exclusion of the n!, thus

$$M(\sigma_n) = (1/n!)L(\sigma_n). \tag{3.12}$$

As the measure of any simplex must be a property of that simplex alone, it follows that any ambiguity in the computation of, for example, $L_{\mu\nu}(\sigma_2)$ must be resolved in the process of computing $L(\sigma_2)$. This circumstance is also evident from the fact that (3.11) is a scalar equation.

For the remaining part of this section it is assumed that the dimension of the complex is 3. After presenting and justifying the definition of exterior differentiation the result will be extended to higher dimensions.

Consider a typical three-simplex σ_3 . Suppose that σ_3 has σ_2 as a base and that σ_1 is a leg of σ_3 but not of σ_2 . Then from (3.10c)

$$L^{\rho\mu\nu}(\sigma_3) = L^{\rho}(\sigma_1)L^{\mu\nu}(\sigma_2) - L^{\mu}(\sigma_1)L^{\rho\nu}(\sigma_2) + L^{\nu}(\sigma_1)L^{\rho\mu}(\sigma_2).$$
(3.13)

Now suppose that $n_{\rho}(\sigma_2)$ is a unit vector normal to the base σ_2 . Then a contraction of (3.13) with n_0 results in

$$n_{\rho}(\sigma_2)L^{\rho\mu\nu}(\sigma_3) = n_{\rho}(\sigma_2)L^{\rho}(\sigma_1)L^{\mu\nu}(\sigma_2).$$

However, $n_{\rho}(\sigma_2)L^{\rho}(\sigma_1)$ is the projection of L^{ρ} in the direction of n_{ρ} , which is just the height of σ_1 above σ_2 , which in turn is just $L(\sigma_3)/L(\sigma_2)$. This leads to

$$\frac{L^{\mu\nu}(\sigma_2)}{L(\sigma_2)} = n_{\rho}(\sigma_2) \frac{L^{\rho\mu\nu}(\sigma_3)}{L(\sigma_3)}.$$
 (3.14)

This expression will be used to obtain a relation between a sum of a two-form over a two surface and a sum of a threeform over a three-surface. A similar relation, on a smooth manifold, will involve the exterior derivative of a two-form. The essence of our definition of exterior differentiation is that it is chosen so as to mimic the usual form of Stokes' theorem.

The expression (3.14) is easily generalized to complexes of dimension greater than 3. Suppose that σ_{m-1} is one face of σ_m and that σ_m is one m-simplex of an m-complex. If the unit inward normal to σ_{m-1} is $n_{\rho}(\sigma_{m-1})$ then

$$\frac{L^{\mu_{1}\mu_{2}\cdots\mu_{m-1}}(\sigma_{m-1})}{L(\sigma_{m-1})} = n_{\mu_{1}}(\sigma_{m-1})\frac{L^{\mu_{1}\mu_{2}\cdots\mu_{m}}(\sigma_{m})}{L(\sigma_{m})}.$$
 (3.15)

This expression can be proved with techniques similar to those that led to (3.14).

IV. EXTERIOR DIFFERENTIATION

Consider a complex ρ_3 that has been subdivided into a set of three-simplexes. The integral of any two-form A * over a two-dimensional subcomplex ρ_2 of ρ_3 is defined, in a coordinate frame, as

$$I(\rho_2, \rho_3) = \sum_{\alpha, \text{in } \rho_2, \int \sigma_2} A^*_{\mu\nu}(x^{\alpha}) dx^{\mu} \wedge dx^{\nu}. \tag{4.1}$$

Similarly, the integral of a three-form B^* over the complex is defined as

$$I(\rho_3) = \sum_{\sigma_3 \text{in } \rho_3} \int_{\sigma_3} B^*_{\mu\nu\tau}(x^{\alpha}) dx^{\mu} \wedge dx^{\nu} \wedge dx^{\tau}. \tag{4.2}$$

Define the quantities $A_{\mu\nu}(\sigma_2)$ and $B_{\mu\nu\tau}(\sigma_3)$ via the equations

$$A_{\mu\nu}(\sigma_2)M(\sigma_2) = \int_{\sigma_2} A_{\mu\nu}^*(x^{\alpha})d^2S,$$
 (4.3a)

$$B_{\mu\nu\tau}(\sigma_3)M(\sigma_3) = \int_{\sigma_2} B^*_{\mu\nu\tau}(x^{\alpha})d^3S,$$
 (4.3b)

with d^iS and $M(\sigma_i)$ being the differential and total measures of the σ_i , respectively. The $A_{\mu\nu}(\sigma_2)$ and $B_{\mu\nu\tau}(\sigma_3)$ are the averages of their associated forms over the simplexes σ_2 and σ_3 . The relations (4.1) and (4.2) may now be rewritten as

$$I(\rho_2, \rho_3) = \sum_{1 = 1}^{n} A_{\mu\nu}(\sigma_2) L^{\mu\nu}(\sigma_2)$$
 (4.4a)

and

$$I(\rho_3) = \sum_{\sigma_3 \text{in } \rho_3} B_{\mu\nu\tau}(\sigma_3) L^{\mu\nu\tau}(\sigma_3). \tag{4.4b}$$

Suppose now that the subcomplex ρ_2 is the boundary of ρ_3 . Our aim is to show that $I(\partial \rho_3, \rho_3)$ may be evaluated either directly from (4.4a) or via an expression similar to (4.4b). The expression (4.4a) may be rewritten as a sum over all σ_3 's of the complex by introducing

$$J(\sigma_3) = \sum_{\sigma_2 \text{in } \partial \sigma_3} A_{\mu\nu}(\sigma_2) L^{\mu\nu}(\sigma_2). \tag{4.5}$$

Then

$$I(\partial \rho_3, \rho_3) = \sum_{\sigma, \text{in } \rho_3} J(\sigma_3), \tag{4.6}$$

since all σ_2 's on the interior of ρ_3 will be counted twice, each with opposite orientations, and will therefore cancel each other. Substitution of (3.14) in (4.15) and the resultant expression in (4.6) leads to

$$I(\partial \rho_{3}, \rho_{3}) = \sum_{\sigma_{2} \text{in } \partial \rho_{3}} A_{\mu\nu}(\sigma_{2}) L^{\mu\nu}(\sigma_{2})$$

$$= \sum_{\sigma_{3} \text{in } \rho_{3}} \sum_{\sigma_{2} \text{in } \partial \sigma_{3}} n_{\rho}(\sigma_{2}) A_{\mu\nu}(\sigma_{2}) \frac{L(\sigma_{2})}{L(\sigma_{3})} L^{\rho\mu\nu}(\sigma_{3}).$$
(4.7)

This expression is greatly simplified by writing

$$A(\sigma_2) = A_{\mu\nu}(\sigma_2)L^{\mu\nu}(\sigma_2), \tag{4.8a}$$

and

$$dA(\sigma_3) = \sum_{\sigma_2 \text{in } \partial \sigma_3} n_{\rho}(\sigma_2) A_{\mu\nu}(\sigma_2) \frac{L(\sigma_2)}{L(\sigma_3)} L^{\rho\mu\nu}(\sigma_3), \qquad (4.8b)$$

for then

$$I(\partial \rho_3, \rho_3) = \sum_{\sigma_2 \text{in } \partial \rho_3} A(\sigma_2) = \sum_{\sigma_2 \text{in } \rho_2} dA(\sigma_3). \tag{4.9}$$

In this form the similarity of this expression with the usual continuum form of Stokes' theorem is quite apparent. The relation (4.8a) defines the value of the two-form A on σ_2 and (4.8b) defines its exterior derivative evaluated on σ_3 .

An analysis similar to that which lead to (4.8a), (4.8b), and (4.9) may be applied to complexes of dimension other than 3. Consider a complex ρ_n of dimension n. Suppose there is defined an m-form $A(\sigma_m)$ on each of the σ_m 's of ρ_n . Thus put

$$A(\sigma_m) = A_{\mu_1 \mu_2 \cdots \mu_m}(\sigma_m) L^{\mu_1 \mu_2 \cdots \mu_m}(\sigma_m). \tag{4.10}$$

Then the exterior derivative of A evaluated on σ_{m+1} is defined as

$$dA(\sigma_{m+1}) = \sum_{\sigma_{m} \in \partial \sigma_{m+1}} n_{\mu_{0}}(\sigma_{m}) A_{\mu_{1}\mu_{2}\cdots\mu_{m}}(\sigma_{m})$$

$$\times \frac{L(\sigma_{m})}{L(\sigma_{m+1})} L^{\mu_{0}\mu_{1}\cdots\mu_{m}}(\sigma_{m+1}), \qquad (4.11)$$

and Stokes' theorem takes the form

$$I(\partial \rho_{m+1}, \rho_{m+1}) = \sum_{\sigma_m \text{in} \partial \rho_{m+1}} A(\sigma_m)$$

$$= \sum_{\sigma_{m+1} \text{in} \rho_{m+1}} dA(\sigma_{m+1}). \tag{4.12}$$

If the complex consists of only one σ_{m+1} then this expression reduces to

$$dA\left(\sigma_{m+1}\right) = \sum_{\sigma_{m} \text{ in } \partial \sigma_{m+1}} A\left(\sigma_{m}\right). \tag{4.13}$$

This provides an alternative yet equivalent method for computing the exterior derivative. In some situations this expression may be more useful than (4.11). As an example it will now be shown that the value of a form, twice exterior differentiated, is zero. Consider a set of numbers $B(\sigma_{m-1})$ on the σ_{m-1} of ρ_{m-1} . Suppose that each number arose as the value of an (m-1)-form B on each of the σ_{m-1} of ρ_{m+1} . The exterior derivative of B, evaluated on each σ_m , gives rise to another set of numbers $A(\sigma_m)$ distributed on the σ_m of ρ_{m+1} . Thus

$$A(\sigma_m) = dB(\sigma_m) = \sum_{\sigma_m \text{ in } \partial \sigma_m} B(\sigma_{m-1}). \tag{4.14}$$

Now the exterior derivative of $A(\sigma_m)$ is

$$\begin{split} dA\left(\sigma_{m+1}\right) &= \sum_{\sigma_{m} \text{in } \partial \sigma_{m+1}} A\left(\sigma_{m}\right) \\ &= \sum_{\sigma_{m} \text{in } \partial \sigma_{m+1}} \sum_{\sigma_{m-1} \text{ in } \partial \sigma_{m}} B\left(\sigma_{m-1}\right). \end{split}$$

However each σ_{m-1} is counted twice, each time with opposite orientations, thus

$$dA\left(\sigma_{m+1}\right) = d\ dB\left(\sigma_{m}\right) = 0.$$

Exactly the same result occurs in the continuum theory of differential forms.

As another example consider the flux of a constant vector A, with components A^{μ} , over the surface of one m-sim-

plex in an m-complex. Clearly this quantity vanishes and is expressed as

$$0 = \sum_{\sigma_{m-1} \text{in } \partial \sigma_m} A^{\mu} n_{\mu} (\sigma_{m-1}) L(\sigma_{m-1}).$$

As this expression is true for any constant field A it follows that

$$0 = \sum_{\sigma_{m-1} \text{in } \partial \sigma_{m}} n_{\mu}(\sigma_{m-1}) L(\sigma_{m-1}).$$

Unfortunately, since n_{μ} need not be continuous across each σ_{m-1} , this expression cannot be applied directly to complexes of more than one m-simplex. However after a contraction with $L^{\mu\mu_1\mu_2\cdots\mu_{m-1}}(\sigma_m)$ and using (3.15) this expression reduces to

$$0 = \sum_{\sigma_{m-1} \text{in } \partial \sigma_m} L^{\mu_1 \mu_2 \cdots \mu_{m-1}} (\sigma_{m-1}),$$

which is easily extended to complexes, thus

$$0 = \sum_{\sigma_{m-1} \text{in } \partial \rho_m} L^{\mu_1 \mu_2 \cdots \mu_{m-1}} (\sigma_{m-1}).$$

This expression can also be proved directly from the definition (3.9).

V. THE CONTINUUM LIMIT

The definitions (4.10) and (4.11) may be extended to complexes built from blocks other than simplexes. For example, an initial manifold could be constructed by piecing together a sequence of three-dimensional cubes. Each such cube could be subdivided, by the addition of extra vertices, legs, and faces, into a set of three-simplexes thus producing a three-complex. To this complex the identity (4.12) would apply. However the terms of this expression may be regrouped so that those terms involving the faces of the threesimplexes are combined into terms involving the faces of the cubes. Similarly the terms involving the three-simplexes would be grouped into terms involving the cubes. In effect the expression (4.12) is unaltered except that the objects in the summation are now parallelograms and parallel cubes instead of triangles and tetrahedrons.

To show that the similarity of (4.9) and (4.12) with Stokes' theorem is not just a consequence of formal algebraic manipulations, the nature of (4.9), over a sequence of complexes, will now be investigated. The following assumptions are necessary.

- (i) The dimension of the complex is 3. A similar analysis may be used for higher dimensions.
- (ii) The sequence of complexes converges, as the number of σ_3 's is increased without limit while keeping the total measure, fixed, to a smooth differentiable manifold.
- (iii) The σ_3 's of each complex are sufficiently small that the values of $A_{\mu\nu}^*$ on the faces σ_2 of σ_3 may be derived from a Taylor series based at some point within σ_3 . Thus

$$A^{*}_{\mu\nu}(x^{\alpha}) = A^{*}_{\mu\nu} + A^{*}_{\mu\nu,\rho} \delta x^{\rho} + O(\delta x)^{2}, \tag{5.1}$$

with $\delta x^{\alpha} = x^{\alpha} - x_0^{\alpha}$, and x_0^{α} is the point, within σ_3 , from which the Taylor series is developed.

(iv) All σ_3 's are three-cubes (parallelepipeds).

For a cube the inward pointing normals for two opposite

faces are equal apart from their directions. Thus the terms in (4.8b) may be regrouped as

$$dA\left(\sigma_{3}\right) = \sum_{\text{three adjacent}} n_{\rho}(\sigma_{2}) \Delta A_{\mu\nu}(\sigma_{2}) \frac{L\left(\sigma_{2}\right)}{L\left(\sigma_{3}\right)} L^{\rho\mu\nu}(\sigma_{3}). \tag{5.2}$$

However, from (4.3a)

$$\Delta A_{\mu\nu}(\sigma_2)L(\sigma_2) = \int_{\sigma_2 - \sigma_{2*}} A_{\mu\nu}^*(x^{\alpha})d^2S, \qquad (5.3)$$

with σ_{2*} being the face opposite σ_{2} . Substitution of (5.1) into (5.3) and noting that $A_{\mu\nu}^*$ and $A_{\mu\nu,\rho}^*$ are constant throughout σ_3 results in

$$\Delta A_{\mu\nu}(\sigma_2)L(\sigma_2) = A^*_{\mu\nu,\rho} \int_{\sigma_2-\sigma_{2*}} \delta x^{\rho} d^2S.$$

By projecting δx^{ρ} onto the normal and tangential vectors of σ_2 it is not hard to show that this last integral equals $n^{\rho}L(\sigma_3)$.

$$\Delta A_{\mu\nu}(\sigma_2)L(\sigma_2) = A^*_{\mu\nu,\,\rho} n^{\rho}(\sigma_2)L(\sigma_3)$$

and consequently (4.8b) becomes

$$dA\left(\sigma_{3}\right) = A *_{\mu\nu,\,\rho} \sum_{\substack{\text{three adjacent} \\ \text{faces}}} n_{\alpha}(\sigma_{2}) n^{\,\rho}(\sigma_{2}) L^{\,\alpha\mu\nu}(\sigma_{3}) \,\cdot$$

But from (3.10a) and (3.10b) the summation reduces to $L^{\rho\mu\nu}(\sigma_3)$. Thus

$$dA(\sigma_3) = A *_{\mu\nu} {}_{\alpha}(\sigma_3) L^{\rho\mu\nu}(\sigma_3)$$

and (4.7) becomes

$$I(\partial \rho_3, \rho_3) = \sum_{\sigma_2 \text{in } \partial \rho_3} A_{\mu\nu}(\sigma_2) L^{\mu\nu}(\sigma_2)$$

$$= \sum_{\sigma_3 \text{in } \rho_3} A_{\mu\nu, \rho}(\sigma_3) L^{\rho\mu\nu}(\sigma_3), \qquad (5.4)$$

with $A_{\mu\nu,\,\rho}(\sigma_3) = A^*_{\mu\nu,\,\rho}$.

This last result shows clearly that this definition of exterior differentiation does reduce to the usual form when the Regge manifold and the forms built on it are made smooth and differentiable.

VI. CONCLUSION

It has been shown that the concept of exterior differentiation has a natural extension to the Regge calculus. Results similar to (4.9) may be found in references 8-10. The motivation for the development of a Regge version of exterior differentiation arises in the attempt to show that certain Regge expressions "converge" to their usual classical counterparts under certain conditions. An example of this process, that our version of Stokes' theorem reduces to its usual form when applied to differtiable forms, has been presented in Sec. V. A more ambitious project would be to prove (or disprove) that the Regge field equations reduce to the Einstein field equations when an appropriate limiting process is applied. This may form the basis of a future investigation.

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