Sample Paper for the aomart Class

By AMERICAN MATHEMATICAL SOCIETY and BORIS VEYTSMAN[®], with Appendix by Frodo Baggins and Afterword by Bilbo Baggins[®]

Abstract

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<u>11</u>	This is a test file for aomart class based on the testmath.tex file from	
<u>12</u>	the amsmath distribution.	
<u>13</u>	It was changed to test the features of the Annals of Mathematics class.	
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<u>38</u> Keywords: Hamiltonian paths, Typesetting

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 $[\]underbrace{\frac{11}{42}}{\text{(C) } 2008-2020 \text{ Boris Veytsman.}}$

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1. Introduction

¹⁹ This paper demonstrates the use of aomart class. It is based on testmath.tex ²⁰ from $\mathcal{A}_{\mathcal{M}}\mathcal{S}$ -IATEX distribution. The text is (slightly) reformatted according to ²¹ the requirements of the aomart style. See also [12], [22], [17], [1], [16], [15], ²² [24], [23], and [6].

It is always a pleasure to cite Knuth [9].

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2. Enumeration of Hamiltonian paths in a graph

Let $\mathbf{A} = (a_{ij})$ be the adjacency matrix of graph G. The corresponding Kirchhoff matrix $\mathbf{K} = (k_{ij})$ is obtained from \mathbf{A} by replacing in $-\mathbf{A}$ each diagonal entry by the degree of its corresponding vertex; i.e., the *i*th diagonal entry is identified with the degree of the *i*th vertex. It is well known that

<u>31</u> (1) det $\mathbf{K}(i|i)$ = the number of spanning trees of G, i = 1, ..., n

 $\frac{32}{2}$ where $\mathbf{K}(i|i)$ is the *i*th principal submatrix of \mathbf{K} .

$\frac{33}{34}$ \det\mathbf{K}(i|i)=\text{ the number of spanning trees of \$G\$},

Let $C_{i(j)}$ be the set of graphs obtained from G by attaching edge $(v_i v_j)$ is each spanning tree of G. Denote by $C_i = \bigcup_j C_{i(j)}$. It is obvious that the collection of Hamiltonian cycles is a subset of C_i . Note that the cardinality of C_i is $k_{ii} \det \mathbf{K}(i|i)$. Let $\widehat{X} = \{\widehat{x}_1, \ldots, \widehat{x}_n\}$.

<u>40</u> Define multiplication for the elements of \widehat{X} by

$$\frac{41}{42} (2) \qquad \hat{x}_i \hat{x}_j = \hat{x}_j \hat{x}_i, \quad \hat{x}_i^2 = 0, \quad i, j = 1, \dots, n.$$

Proof: page numbers may be temporary

Are these quotations necessary? 18

¹ Let $\hat{k}_{ij} = k_{ij}\hat{x}_j$ and $\hat{k}_{ij} = -\sum_{j \neq i} \hat{k}_{ij}$. Then the number of Hamiltonian cycles ² H_c is given by the relation [13]

$$\frac{3}{\frac{4}{5}} \quad (3) \qquad \left(\prod_{j=1}^{n} \hat{x}_{j}\right) H_{c} = \frac{1}{2} \hat{k}_{ij} \det \widehat{\mathbf{K}}(i|i), \qquad i = 1, \dots, n.$$

The task here is to express (3) in a form free of any \hat{x}_i , i = 1, ..., n. The result also leads to the resolution of enumeration of Hamiltonian paths in a graph.

It is well known that the enumeration of Hamiltonian cycles and paths <u>9</u> in a complete graph K_n and in a complete bipartite graph $K_{n_1n_2}$ can only 10 be found from *first combinatorial principles* [7]. One wonders if there exists a <u>11</u> formula which can be used very efficiently to produce K_n and $K_{n_1n_2}$. Recently, $\underline{12}$ using Lagrangian methods, Goulden and Jackson have shown that H_c can be 13expressed in terms of the determinant and permanent of the adjacency matrix 14[5]. However, the formula of Goulden and Jackson determines neither K_n nor <u>15</u> $K_{n_1n_2}$ effectively. In this paper, using an algebraic method, we parametrize 16 the adjacency matrix. The resulting formula also involves the determinant 1718 and permanent, but it can easily be applied to K_n and $K_{n_1n_2}$. In addition, we eliminate the permanent from H_c and show that H_c can be represented by 19a determinantal function of multivariables, each variable with domain $\{0, 1\}$. <u>20</u> Furthermore, we show that H_c can be written by number of spanning trees of $\underline{21}$ subgraphs. Finally, we apply the formulas to a complete multigraph $K_{n_1...n_p}$. 22 The conditions $a_{ij} = a_{ji}$, i, j = 1, ..., n, are not required in this paper.

The conditions $a_{ij} = a_{ji}$, i, j = 1, ..., n, are not required in this paper. All formulas can be extended to a digraph simply by multiplying H_c by 2. Some other discussion can be found in [4] and [3].

3. Main theorem

29 Notation. For $p, q \in P$ and $n \in \omega$ we write $(q, n) \leq (p, n)$ if $q \leq p$ and 30 $A_{q,n} = A_{p,n}$.

\begin{notation} For \$p,q\in P\$ and \$n\in\omega\$

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\end{notation}
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<u>35</u> Let $\mathbf{B} = (b_{ij})$ be an $n \times n$ matrix. Let $\mathbf{n} = \{1, \ldots, n\}$. Using the properties <u>36</u> of (2), it is readily seen that

Lemma 3.1.

$$\frac{-}{39} \quad (4) \qquad \qquad \prod_{i \in \mathbf{n}} \left(\sum_{j \in \mathbf{n}} b_{ij} \hat{x}_i \right) = \left(\prod_{i \in \mathbf{n}} \hat{x}_i \right) \operatorname{per} \mathbf{B}$$

 $\frac{41}{42}$ where per **B** is the permanent of **B**.

n.

Let
$$\widehat{Y} = {\widehat{y}_1, \ldots, \widehat{y}_n}$$
. Define multiplication for the elements of \widehat{Y} by

(5)
$$\hat{y}_i \hat{y}_j + \hat{y}_j \hat{y}_i = 0, \quad i, j = 1, \dots,$$

Then, it follows that

LEMMA 3.2.

(6)
$$\prod_{i \in \mathbf{n}} \left(\sum_{j \in \mathbf{n}} b_{ij} \hat{y}_j \right) = \left(\prod_{i \in \mathbf{n}} \hat{y}_i \right) \det \mathbf{B}.$$

Note that all basic properties of determinants are direct consequences of 10Lemma 3.2. Write 11

$$\frac{12}{13} \quad (7) \qquad \sum_{j \in \mathbf{n}} b_{ij} \hat{y}_j = \sum_{j \in \mathbf{n}} b_{ij}^{(\lambda)} \hat{y}_j + (b_{ii} - \lambda_i) \hat{y}_i \hat{y}_j$$

where 15

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$$\frac{16}{17} \quad (8) \qquad \qquad b_{ii}^{(\lambda)} = \lambda_i, \quad b_{ij}^{(\lambda)} = b_{ij}, \quad i \neq j.$$

Let $\mathbf{B}^{(\lambda)} = (b_{ij}^{(\lambda)})$. By (6) and (7), it is straightforward to show the following $\underline{18}$ result: $\underline{19}$

THEOREM 3.3.

$$\frac{-22}{23} (9) \qquad \det \mathbf{B} = \sum_{l=0}^{n} \sum_{I_l \subseteq n} \prod_{i \in I_l} (b_{ii} - \lambda_i) \det \mathbf{B}^{(\lambda)}(I_l | I_l),$$

where $I_l = \{i_1, \ldots, i_l\}$ and $\mathbf{B}^{(\lambda)}(I_l|I_l)$ is the principal submatrix (obtained from $\underline{25}$ $\mathbf{B}^{(\lambda)}$ by deleting its i_1, \ldots, i_l rows and columns). $\underline{26}$

<u>27</u> *Remark* 3.1 (convention). Let **M** be an $n \times n$ matrix. The convention $\underline{28}$ $\mathbf{M}(\mathbf{n}|\mathbf{n}) = 1$ has been used in (9) and hereafter. <u>29</u>

Before proceeding with our discussion, we pause to note that Theorem 3.3 30 yields immediately a fundamental formula which can be used to compute the $\underline{31}$ coefficients of a characteristic polynomial [14]: 32

$$\begin{array}{l} \frac{33}{34} \\ \frac{34}{35} \\ \frac{35}{36} \end{array} \quad \text{COROLLARY 3.4. Write } \det(\mathbf{B} - x\mathbf{I}) = \sum_{l=0}^{n} (-1)^{l} b_{l} x^{l}. \text{ Then} \\ b_{l} = \sum_{I_{l} \subseteq \mathbf{n}} \det \mathbf{B}(I_{l}|I_{l}). \end{array}$$

Let

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$$\frac{39}{40} (11) \qquad \mathbf{K}(t \ t_1 \ t_2) = \begin{pmatrix} D_1 t & -a_{12} t_2 \ \dots \ -a_{1n} t_n \\ -a_{21} t_1 & D_2 t \ \dots \ -a_{2n} t_n \end{pmatrix}$$

Proof: page numbers may be temporary

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1 $\begin{pmatrix} D_1t\&-a_{12}t_2\&\dots\&-a_{1n}t_n\$ $\underline{2}$ $-a_{21}t_1\&D_2t\&\dots\&-a_{2n}t_n\$ <u>3</u> $hdotsfor[2]{4}\$ 4 $-a_{n1}t_1\&-a_{n2}t_2\&\dots\&D_nt\end{pmatrix}$ <u>5</u> where <u>6</u> $D_i = \sum_{i \in \mathcal{I}} a_{ij} t_j, \quad i = 1, \dots, n.$ (12)7 8 <u>9</u> Set $D(t_1,\ldots,t_n) = \frac{\delta}{\delta t} \det \mathbf{K}(t,t_1,\ldots,t_n)|_{t=1}.$ <u>10</u> <u>11</u> Then 12 $D(t_1,\ldots,t_n) = \sum_{i=1}^{n} D_i \det \mathbf{K}(t=1,t_1,\ldots,t_n;i|i),$ $\underline{13}$ (13)14 $\underline{15}$ where $\mathbf{K}(t = 1, t_1, \dots, t_n; i | i)$ is the *i*th principal submatrix of $\mathbf{K}(t = 1, t_1, \dots, t_n)$. <u>16</u> Theorem 3.3 leads to $\underline{17}$ (14) det $\mathbf{K}(t_1, t_1, \dots, t_n) = \sum_{I \in \mathbf{n}} (-1)^{|I|} t^{n-|I|} \prod_{i \in I} t_i \prod_{j \in I} (D_j + \lambda_j t_j) \det \mathbf{A}^{(\lambda t)}(\overline{I}|\overline{I}).$ 18 <u>19</u> <u>20</u> Note that $\underline{21}$ (15) $\det \mathbf{K}(t=1,t_1,\ldots,t_n) = \sum_{I \in \mathbf{n}} (-1)^{|I|} \prod_{i \in I} t_i \prod_{j \in I} (D_j + \lambda_j t_j) \det \mathbf{A}^{(\lambda)}(\overline{I}|\overline{I}) = 0.$ <u>22</u> $\underline{23}$ $\underline{24}$ Let $t_i = \hat{x}_i, i = 1, \dots, n$. Lemma 3.1 yields <u>25</u> <u>26</u> (16) $\left(\sum a_{l_i} x_i\right) \det \mathbf{K}(t=1,x_1,\ldots,x_n;l|l)$ 27 $\underline{28}$ $= \left(\prod_{i \in \mathbf{n}} \hat{x}_i\right) \sum_{I \subset \mathbf{n} - II} (-1)^{|I|} \operatorname{per} \mathbf{A}^{(\lambda)}(I|I) \det \mathbf{A}^{(\lambda)}(\overline{I} \cup \{l\} | \overline{I} \cup \{l\}).$ 29 <u>30</u> <u>31</u> \begin{multline} <u>32</u> \biggl(\sum_{\,i\in\mathbf{n}}a_{1 _i}x_i\biggr) <u>33</u> $\det\{K\}(t=1,x_1,dots,x_n;1 | 1) \$ <u>34</u> =\biggl(\prod_{\,i\in\mathbf{n}}\hat x_i\biggr) <u>35</u> $\sum_{I \in \mathbb{N}^{1}}$ <u>36</u> $(-1)^{\left(1\right)}\left(1\right)^{\left(1\right)}(1|1)$ <u>37</u> $\det \mathbb{A}^{(\lambda)}$ <u>38</u> (\overline I\cup\{1 \}|\overline I\cup\{1 \}). 39 \label{sum-ali} <u>40</u> \end{multline} <u>41</u> By (3), (6), and (7), we have <u>42</u>

PROPOSITION 3.5. (17) $H_c = \frac{1}{2n} \sum_{l=0}^{n} (-1)^l D_l,$

 $\underline{5}$ where

$$\frac{\frac{6}{7}}{\frac{8}{2}} (18) \qquad D_l = \sum_{I_l \subseteq \mathbf{n}} D(t_1, \dots, t_n) 2|_{t_i = \begin{cases} 0, \text{ if } i \in I_l \\ 1, \text{ otherwise} \end{cases}}, i=1,\dots,n.$$

4. Application

 $\frac{10}{11}$ We consider here the applications of Theorems 5.1 and 5.2 on page 23 to a complete multipartite graph $K_{n_1...n_p}$. It can be shown that the number of spanning trees of $K_{n_1...n_p}$ may be written

$$\frac{14}{15} (19) T = n^{p-2} \prod_{i=1}^{p} (n-n_i)^{n_i-1}$$

 $\frac{16}{17}$ where

 $\underline{19}$

$$\underbrace{18} (20) \qquad \qquad n = n_1 + \dots + n_p.$$

It follows from Theorems 5.1 and 5.2 that

$$\frac{20}{21} \qquad H_c = \frac{1}{2n} \sum_{l=0}^n (-1)^l (n-l)^{p-2} \sum_{l_1 + \dots + l_p = l} \prod_{i=1}^p \binom{n_i}{l_i}$$

$$\frac{23}{24} \qquad \qquad \cdot \left[(n-l) - (n_i - l_i) \right]^{n_i - l_i} \cdot \left[(n-l)^2 - \sum_{j=1}^p (n_i - l_i)^2 \right].$$

26 ... \binom{n_i}{l _i}\\

 $\frac{\underline{27}}{\underline{28}}$ and

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(22)
$$H_c = \frac{1}{2} \sum_{l=0}^{n-1} (-1)^l (n-l)^{p-2} \sum_{l_1 + \dots + l_p = l} \prod_{i=1}^p \binom{n_i}{l_i}$$

$$\cdot [(n-l) - (n_i - l_i)]^{n_i - l_i} \left(1 - \frac{l_p}{n_p}\right) [(n-l) - (n_p - l_p)].$$

The enumeration of H_c in a $K_{n_1\cdots n_p}$ graph can also be carried out by Theorem 7.2 or 7.3 together with the algebraic method of (2). Some elegant representations may be obtained. For example, H_c in a $K_{n_1n_2n_3}$ graph may be written

(23)
$$H_{c} = \frac{n_{1}! n_{2}! n_{3}!}{n_{1} + n_{2} + n_{3}} \sum_{i} \left[\binom{n_{1}}{i} \binom{n_{2}}{n_{3} - n_{1} + i} \binom{n_{3}}{n_{3} - n_{2} + i} \right]$$
$$\binom{n_{3}}{i} = \binom{n_{1}! n_{2}! n_{3}!}{(n_{1} - 1) (n_{2} - 1) (n_{3} - 1)}$$

$$\frac{41}{42} \qquad \qquad + \binom{n_1 - 1}{i} \binom{n_2 - 1}{n_3 - n_1 + i} \binom{n_3 - 1}{n_3 - n_2 + i} \end{bmatrix}.$$

Proof: page numbers may be temporary

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5. Secret key exchanges

Modern cryptography is fundamentally concerned with the problem of secure private communication. A Secret Key Exchange is a protocol where Alice and Bob, having no secret information in common to start, are able to agree on a common secret key, conversing over a public channel. The notion of a Secret Key Exchange protocol was first introduced in the seminal paper of Diffie and Hellman [2]. [2] presented a concrete implementation of a Secret Key Exchange protocol, dependent on a specific assumption (a variant on the discrete log), specially tailored to yield Secret Key Exchange. Secret Key Exchange is of course trivial if trapdoor permutations exist. However, there is no known implementation based on a weaker general assumption.

The concept of an informationally one-way function was introduced in [8]. We give only an informal definition here: 14

<u>15</u> Definition 5.1 (one way). A polynomial time computable function f =<u>16</u> $\{f_k\}$ is informationally one-way if there is no probabilistic polynomial time 17 algorithm which (with probability of the form $1 - k^{-e}$ for some e > 0) returns 18 on input $y \in \{0, 1\}^k$ a random element of $f^{-1}(y)$.

In the non-uniform setting [8] show that these are not weaker than one-way functions:

<u>22</u> THEOREM 5.1 ([8] (non-uniform)). The existence of informationally one- $\underline{23}$ way functions implies the existence of one-way functions.

We will stick to the convention introduced above of saying "non-uniform" before the theorem statement when the theorem makes use of non-uniformity. It should be understood that if nothing is said then the result holds for both the uniform and the non-uniform models.

It now follows from Theorem 5.1 that

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30 THEOREM 5.2 (non-uniform). Weak SKE implies the existence of a oneway function.

<u>32</u> More recently, the polynomial-time, interior point algorithms for linear <u>33</u> programming have been extended to the case of convex quadratic programs [19] 34and [21], certain linear complementarity problems [11] and [18], and the nonlin-35 ear complementarity problem [10]. The connection between these algorithms 36 and the classical Newton method for nonlinear equations is well explained in <u>37</u> [11]. 38

6. Review

We begin our discussion with the following definition:

$$\lim_{v \to 0} \frac{H(z+v) - H(z) - BH(z)v}{\|v\|} = 0.$$

The function H is *B*-differentiable in set S if it is B-differentiable at every point in S. The B-derivative BH(z) is said to be strong if

$$\frac{10}{11} \qquad \qquad \lim_{(v,v')\to(0,0)} \frac{H(z+v) - H(z+v') - BH(z)(v-v')}{\|v-v'\|} = 0.$$

 $\begin{array}{ll} \frac{12}{13} & \text{LEMMA 6.1. There exists a smooth function } \psi_0(z) \ defined \ for \ |z| > 1-2a \\ satisfying \ the \ following \ properties: \end{array}$

(i) $\psi_0(z)$ is bounded above and below by positive constants $c_1 \leq \psi_0(z) \leq c_2$.

(ii) If
$$|z| > 1$$
, then $\psi_0(z) = 1$.

(iii) For all z in the domain of ψ_0 , $\Delta_0 \ln \psi_0 \ge 0$.

 $\frac{1}{18} \text{ (iv) } If 1 - 2a < |z| < 1 - a, then \Delta_0 \ln \psi_0 \ge c_3 > 0.$

19 Proof. We choose $\psi_0(z)$ to be a radial function depending only on r = |z|. 20 Let $h(r) \ge 0$ be a suitable smooth function satisfying $h(r) \ge c_3$ for 1 - 2a < 2121 |z| < 1 - a, and h(r) = 0 for $|z| > 1 - \frac{a}{2}$. The radial Laplacian

$$\Delta_0 \ln \psi_0(r) = \left(\frac{d^2}{dr^2} + \frac{1}{r}\frac{d}{dr}\right) \ln \psi_0(r)$$

has smooth coefficients for r > 1 - 2a. Therefore, we may apply the existence and uniqueness theory for ordinary differential equations. Simply let $\ln \psi_0(r)$ be the solution of the differential equation

$$\frac{28}{29}\qquad \qquad \left(\frac{d^2}{dr^2} + \frac{1}{r}\frac{d}{dr}\right)\ln\psi_0(r) = h(r)$$

³⁰ with initial conditions given by $\ln \psi_0(1) = 0$ and $\ln \psi'_0(1) = 0$.

Next, let D_{ν} be a finite collection of pairwise disjoint disks, all of which are contained in the unit disk centered at the origin in C. We assume that $D_{\nu} = \{z \mid |z - z_{\nu}| < \delta\}$. Suppose that $D_{\nu}(a)$ denotes the smaller concentric disk $D_{\nu}(a) = \{z \mid |z - z_{\nu}| \le (1 - 2a)\delta\}$. We define a smooth weight function $\Phi_0(z)$ for $z \in C - \bigcup_{\nu} D_{\nu}(a)$ by setting $\Phi_0(z) = 1$ when $z \notin \bigcup_{\nu} D_{\nu}$ and $\Phi_0(z) = \psi_0((z - z_{\nu})/\delta)$ when z is an element of D_{ν} . It follows from Lemma 6.1 that Φ_0 satisfies the properties:

(i) $\Phi_0(z)$ is bounded above and below by positive constants $c_1 \leq \Phi_0(z) \leq c_2$.

(ii) $\Delta_0 \ln \Phi_0 \ge 0$ for all $z \in C - \bigcup_{\nu} D_{\nu}(a)$, the domain where the function Φ_0 is defined.

 $\frac{41}{42} \text{ (iii) } \Delta_0 \ln \Phi_0 \ge c_3 \delta^{-2} \text{ when } (1-2a)\delta < |z-z_\nu| < (1-a)\delta.$

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Let A_{ν} denote the annulus $A_{\nu} = \{(1-2a)\delta < |z-z_{\nu}| < (1-a)\delta\}$, and 1 set $A = \bigcup_{\nu} A_{\nu}$. The properties (2) and (3) of Φ_0 may be summarized as $\underline{2}$ <u>3</u> $\Delta_0 \ln \Phi_0 \geq c_3 \delta^{-2} \chi_A$, where χ_A is the characteristic function of A. 4 Suppose that α is a nonnegative real constant. We apply Proposition 3.5 <u>5</u> with $\Phi(z) = \Phi_0(z)e^{\alpha|z|^2}$. If $u \in C_0^{\infty}(\mathbb{R}^2 - \bigcup_{\nu} D_{\nu}(a))$, assume that \mathcal{D} is a <u>6</u> bounded domain containing the support of u and $A \subset \mathcal{D} \subset R^2 - \bigcup_{\nu} D_{\nu}(a)$. A 7 calculation gives 8 <u>9</u> $\int_{\mathcal{D}} \left| \overline{\partial} u \right|^2 \Phi_0(z) e^{\alpha |z|^2} \ge c_4 \alpha \int_{\mathcal{D}} |u|^2 \Phi_0 e^{\alpha |z|^2} + c_5 \delta^{-2} \int_{\mathcal{A}} |u|^2 \Phi_0 e^{\alpha |z|^2}.$ 10 <u>11</u> The boundedness, property (1) of Φ_0 , then yields $\underline{12}$ $\underline{13}$ $\int_{\Sigma} \left| \overline{\partial} u \right|^2 e^{\alpha |z|^2} \ge c_6 \alpha \int_{\Sigma} |u|^2 e^{\alpha |z|^2} + c_7 \delta^{-2} \int_{\Lambda} |u|^2 e^{\alpha |z|^2}.$ 14 $\underline{15}$ Let B(X) be the set of blocks of Λ_X and let b(X) = |B(X)|. If $\phi \in Q_X$ 16 then ϕ is constant on the blocks of Λ_X . 1718 $P_X = \{ \phi \in M \mid \Lambda_\phi = \Lambda_X \},\$ $Q_X = \{ \phi \in M \mid \Lambda_\phi > \Lambda_X \}.$ (24)19If $\Lambda_{\phi} \geq \Lambda_X$ then $\Lambda_{\phi} = \Lambda_Y$ for some $Y \geq X$ so that <u>20</u> $\underline{21}$ $Q_X = \bigcup_{Y > X} P_Y.$ <u>22</u> $\underline{23}$ Thus by Möbius inversion $\underline{24}$ $\underline{25}$ $|P_Y| = \sum_{X > V} \mu(Y, X) |Q_X|.$ <u>26</u> $\underline{27}$ Thus there is a bijection from Q_X to $W^{B(X)}$. In particular $|Q_X| = w^{b(X)}$. $\underline{28}$ Next note that $b(X) = \dim X$. We see this by choosing a basis for X $\underline{29}$ <u>30</u> consisting of vectors v^k defined by <u>31</u> $v_i^k = \begin{cases} 1 & \text{if } i \in \Lambda_k, \\ 0 & \text{otherwise.} \end{cases}$ <u>32</u> <u>33</u> 34 $\left[v^{k}_{i}\right]$ <u>35</u> \begin{cases} 1 & \text{if \$i \in \Lambda_{k}\$},\\ 36 0 &\text{otherwise.} \end{cases} <u>37</u> \backslash] <u>38</u> 39 LEMMA 6.2. Let \mathcal{A} be an arrangement. Then <u>40</u> $\chi(\mathcal{A}, t) = \sum_{\mathcal{B} \subset \mathcal{A}} (-1)^{|\mathcal{B}|} t^{\dim T(\mathcal{B})}.$ <u>41</u> <u>42</u>

In order to compute R'' recall the definition of S(X, Y) from Lemma 3.1. 1 $\underline{2}$ Since $H \in \mathcal{B}$, $\mathcal{A}_H \subseteq \mathcal{B}$. Thus if $T(\mathcal{B}) = Y$ then $\mathcal{B} \in S(H, Y)$. Let $L'' = L(\mathcal{A}'')$. <u>3</u> Then $\underline{4}$ $R'' = \sum_{H \in \mathcal{B} \subset A} (-1)^{|\mathcal{B}|} t^{\dim T(\mathcal{B})}$ $\underline{5}$ <u>6</u> $= \sum_{Y \in L''} \sum_{\mathcal{B} \in S(H|Y)} (-1)^{|\mathcal{B}|} t^{\dim Y}$ $\underline{7}$ 8 $= -\sum_{Y \in L''} \sum_{\mathcal{B} \in S(H,Y)} (-1)^{|\mathcal{B} - \mathcal{A}_H|} t^{\dim Y}$ (25) $\underline{9}$ 10 $= -\sum_{Y \in L''} \mu(H, Y) t^{\dim Y}$ 11 <u>12</u> $\underline{13}$ $=-\chi(\mathcal{A}'',t).$ $\underline{14}$ COROLLARY 6.3. Let $(\mathcal{A}, \mathcal{A}', \mathcal{A}'')$ be a triple of arrangements. Then 1516 $\pi(\mathcal{A}, t) = \pi(\mathcal{A}', t) + t\pi(\mathcal{A}'', t).$ <u>17</u> Definition 6.2. Let $(\mathcal{A}, \mathcal{A}', \mathcal{A}'')$ be a triple with respect to the hyperplane <u>18</u> $H \in \mathcal{A}$. Call H a separator if $T(\mathcal{A}) \notin L(\mathcal{A}')$. 19 $\underline{20}$ COROLLARY 6.4. Let $(\mathcal{A}, \mathcal{A}', \mathcal{A}'')$ be a triple with respect to $H \in \mathcal{A}$. $\underline{21}$ (i) If H is a separator then <u>22</u> $\mu(\mathcal{A}) = -\mu(\mathcal{A}'')$ $\underline{23}$ $\underline{24}$ and hence $\underline{25}$ $|\mu(\mathcal{A})| = |\mu(\mathcal{A}'')|.$ $\underline{26}$ (ii) If H is not a separator then <u>27</u> $\underline{28}$ $\mu(\mathcal{A}) = \mu(\mathcal{A}') - \mu(\mathcal{A}'')$ $\underline{29}$ and 30 $|\mu(\mathcal{A})| = |\mu(\mathcal{A}')| + |\mu(\mathcal{A}'')|.$ <u>31</u> <u>32</u> *Proof.* It follows from Theorem 5.1 that $\pi(\mathcal{A}, t)$ has leading term <u>33</u> $(-1)^{r(\mathcal{A})} \mu(\mathcal{A}) t^{r(\mathcal{A})}.$ 34The conclusion follows by comparing coefficients of the leading terms on both $\underline{35}$ sides of the equation in Corollary 6.3. If H is a separator then $r(\mathcal{A}') < r(\mathcal{A})$ 36 <u>37</u> and there is no contribution from $\pi(\mathcal{A}', t)$. <u>38</u>

The Poincaré polynomial of an arrangement will appear repeatedly in these notes. It will be shown to equal the Poincaré polynomial of the graded algebras which we are going to associate with \mathcal{A} . It is also the Poincaré polynomial of the complement $M(\mathcal{A})$ for a complex arrangement. Here we prove

Figure 1.
$$Q(\mathcal{A}_1) = xyz(x-z)(x+z)(y-z)(y+z)$$

Figure 1. $Q(\mathcal{A}_1) = xyz(x-z)(x+z)(y-z)(y+z)$
Figure 2. $Q(\mathcal{A}_2) = xyz(x+y+z)(x+y-z)(x-y+z)(x-y-z)$
that the Poincaré polynomial is the chamber counting function for a real arrangement. The complement $M(\mathcal{A})$ is a disjoint union of chambers
 $M(\mathcal{A}) = \bigcup_{C \in Cham(\mathcal{A})} C.$
The number of chambers is determined by the Poincaré polynomial as follows.
THEOREM 6.5. Let $\mathcal{A}_{\mathbf{R}}$ be a real arrangement. Then
 $|Cham(\mathcal{A}_{\mathbf{R}})| = \pi(\mathcal{A}_{\mathbf{R}}, 1).$
Proof. We check the properties required in Corollary 6.4: (i) follows from
 $\pi(\Phi_t, t) = 1$, and (ii) is a consequence of Corollary 3.4. \Box
THEOREM 6.6. Let ϕ be a protocol for a random pair (X, Y) . If one of
 $\sigma_{\phi}(x', y)$ and $\sigma_{\phi}(x, y')$ is a prefix of the other and $(x, y) \in S_{X,Y}$, then
 $\langle \sigma_j(x', y) \rangle_{j=1}^{\infty} = \langle \sigma_j(x, y) \rangle_{j=1}^{\infty} = \langle \sigma_j(x, y') \rangle_{j=1}^{\infty}.$

Proof. We show by induction on i that

$$\langle \sigma_j(x',y) \rangle_{j=1}^i = \langle \sigma_j(x,y) \rangle_{j=1}^i = \langle \sigma_j(x,y') \rangle_{j=1}^i.$$

4 The induction hypothesis holds vacuously for i = 0. Assume it holds for i-1, in particular $[\sigma_j(x',y)]_{j=1}^{i-1} = [\sigma_j(x,y')]_{j=1}^{i-1}$. Then one of $[\sigma_j(x',y)]_{j=i}^{\infty}$ $\underline{5}$ and $[\sigma_i(x,y')]_{i=i}^{\infty}$ is a prefix of the other which implies that one of $\sigma_i(x',y)$ <u>6</u> $\underline{7}$ and $\sigma_i(x, y')$ is a prefix of the other. If the *i*th message is transmitted by 8 $P_{\mathcal{X}}$ then, by the separate-transmissions property and the induction hypothe-9 sis, $\sigma_i(x,y) = \sigma_i(x,y')$, hence one of $\sigma_i(x,y)$ and $\sigma_i(x',y)$ is a prefix of the $\underline{10}$ other. By the implicit-termination property, neither $\sigma_i(x, y)$ nor $\sigma_i(x', y)$ can <u>11</u> be a proper prefix of the other, hence they must be the same and $\sigma_i(x', y) =$ $\underline{12}$ $\sigma_i(x,y) = \sigma_i(x,y')$. If the *i*th message is transmitted by P_V then, symmetrically, $\sigma_i(x,y) = \sigma_i(x',y)$ by the induction hypothesis and the separate- $\underline{13}$ 14transmissions property, and, then, $\sigma_i(x, y) = \sigma_i(x, y')$ by the implicit termination 15property, proving the induction step. 16

If ϕ is a protocol for (X, Y), and (x, y), (x', y) are distinct inputs in $S_{X,Y}$,

then, by the correct-decision property, $\langle \sigma_j(x,y) \rangle_{j=1}^{\infty} \neq \langle \sigma_j(x',y) \rangle_{j=1}^{\infty}$. Equation (25) defined $P_{\mathcal{Y}}$'s ambiguity set $S_{X|Y}(y)$ to be the set of possible

 $\frac{19}{20} \qquad \text{Equation (25) defined } P_{\mathcal{Y}} \text{ s ambiguity set } S_{X|Y}(y) \text{ to be the set of possible} \\ \frac{20}{21} \quad X \text{ values when } Y = y. \text{ The last corollary implies that for all } y \in S_Y, \text{ the} \\ \frac{21}{21} \quad \text{multiset}^1 \text{ of codewords } \{\sigma_{\phi}(x,y) : x \in S_{X|Y}(y)\} \text{ is prefix free.} \end{cases}$

7. One-way complexity

 $\begin{array}{l} \frac{24}{25} \qquad \hat{C}_1(X|Y), \text{ the one-way complexity of a random pair } (X,Y), \text{ is the number} \\ \frac{25}{26} \qquad \text{of bits } P_{\mathcal{X}} \text{ must transmit in the worst case when } P_{\mathcal{Y}} \text{ is not permitted to transmit} \\ \frac{26}{26} \qquad \text{any feedback messages. Starting with } S_{X,Y}, \text{ the support set of } (X,Y), \text{ we define} \\ \frac{27}{28} \qquad G(X|Y), \text{ the characteristic hypergraph of } (X,Y), \text{ and show that} \end{array}$

$$\hat{C}_1(X|Y) = \left\lceil \log \chi(G(X|Y)) \right\rceil$$

We can now prove a continuity theorem.

THEOREM 7.1. Let $\Omega \subset \mathbf{R}^n$ be an open set, let $u \in BV(\Omega; \mathbf{R}^m)$, and let $\frac{37}{38}$ (26) $T_x^u = \left\{ y \in \mathbf{R}^m : y = \tilde{u}(x) + \left\langle \frac{Du}{|Du|}(x), z \right\rangle \text{ for some } z \in \mathbf{R}^n \right\}$

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 $[\]frac{41}{42}$ ¹A multiset allows multiplicity of elements. Hence, $\{0, 01, 01\}$ is prefix free as a set, but not as a multiset.

 $\begin{array}{ll} \frac{1}{2} & \text{for every } x \in \Omega \backslash S_u. \ Let \ f \colon \mathbf{R}^m \to \mathbf{R}^k \ be \ a \ Lipschitz \ continuous \ function \ such \\ \frac{2}{2} & \text{that } f(0) = 0, \ and \ let \ v = f(u) \colon \Omega \to \mathbf{R}^k. \ Then \ v \in BV(\Omega; \mathbf{R}^k) \ and \end{array}$

$$\frac{3}{4} \quad (27) \qquad \qquad Jv = \left(f(u^+) - f(u^-)\right) \otimes \nu_u \cdot \mathcal{H}_{n-1}\big|_{S_u}$$

In addition, for $|\widetilde{D}u|$ -almost every $x \in \Omega$ the restriction of the function f to T_x^u is differentiable at $\tilde{u}(x)$ and

$$\frac{\frac{8}{9}}{\frac{10}{2}} \quad (28) \qquad \qquad \widetilde{D}v = \nabla(f|_{T^u_x})(\widetilde{u})\frac{\widetilde{D}u}{\left|\widetilde{D}u\right|} \cdot \left|\widetilde{D}u\right|.$$

Before proving the theorem, we state without proof three elementary remarks which will be useful in the sequel.

<u>14</u> Remark 7.1. Let $\omega: [0, +\infty[\rightarrow]0, +\infty[$ be a continuous function such <u>15</u> that $\omega(t) \to 0$ as $t \to 0$. Then

$$\lim_{h\to 0^+} g(\omega(h)) = L \Leftrightarrow \lim_{h\to 0^+} g(h) = L$$

^{<u>18</u>} for any function $g: [0, +\infty[\rightarrow \mathbf{R}.$

 $\frac{20}{21} \qquad Remark 7.2. \text{ Let } g \colon \mathbf{R}^n \to \mathbf{R} \text{ be a Lipschitz continuous function and as-}$ sume that

$$L(z) = \lim_{h \to 0^+} \frac{g(hz) - g(0)}{h}$$

exists for every $z \in \mathbf{Q}^n$ and that L is a linear function of z. Then g is differentiable at 0.

 $\begin{array}{ll} \frac{26}{27} & Remark \ 7.3. \ \text{Let} \ A \colon \mathbf{R}^n \to \mathbf{R}^m \ \text{be a linear function, and let} \ f \colon \mathbf{R}^m \to \mathbf{R} \\ \frac{27}{28} & \text{be a function. Then the restriction of } f \ \text{to the range of } A \ \text{is differentiable at } 0 \\ \frac{28}{29} & \text{if and only if} \ f(A) \colon \mathbf{R}^n \to \mathbf{R} \ \text{is differentiable at } 0 \ \text{and} \end{array}$

$$\nabla(f|_{\operatorname{Im}(A)})(0)A = \nabla(f(A))(0)$$

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Proof. We begin by showing that
$$v \in BV(\Omega; \mathbf{R}^{k})$$
 and

$$\underbrace{33} (29) \qquad |Dv|(B) \le K |Du|(B) \qquad \forall B \in \mathbf{B}(\Omega),$$

where K > 0 is the Lipschitz constant of f. By (13) and by the approximation result quoted in §3, it is possible to find a sequence $(u_h) \subset C^1(\Omega; \mathbf{R}^m)$ converging to u in $L^1(\Omega; \mathbf{R}^m)$ and such that

$$\lim_{h \to +\infty} \int_{\Omega} |\nabla u_h| \, dx = |Du| \, (\Omega).$$

 $\frac{40}{41}$ The functions $v_h = f(u_h)$ are locally Lipschitz continuous in Ω , and the definition of differential implies that $|\nabla v_h| \leq K |\nabla u_h|$ almost everywhere in Ω . The 1 lower semicontinuity of the total variation and (13) yield $\underline{2}$ ſ

$$\frac{3}{4} \quad (30) \qquad |Dv|(\Omega) \leq \liminf_{h \to +\infty} |Dv_h|(\Omega) = \liminf_{h \to +\infty} \int_{\Omega} |\nabla v_h| \, dx \\ \leq K \liminf_{h \to +\infty} \int_{\Omega} |\nabla u_h| \, dx = K |Du|(\Omega).$$

 $\underline{7}$ Since f(0) = 0, we have also

$$\frac{\frac{8}{9}}{\frac{10}{2}} \int_{\Omega} |v| \ dx \le K \int_{\Omega} |u| \ dx;$$

<u>11</u> therefore $u \in BV(\Omega; \mathbf{R}^k)$. Repeating the same argument for every open set $A \subset \Omega$, we get (29) for every $B \in \mathbf{B}(\Omega)$, because |Dv|, |Du| are Radon mea- $\underline{12}$ $\underline{13}$ sures. To prove Lemma 6.1, first we observe that

$$\frac{\underline{14}}{\underline{15}} (31) \qquad S_v \subset S_u, \qquad \tilde{v}(x) = f(\tilde{u}(x)) \qquad \forall x \in \Omega \backslash S_u.$$

 $\underline{16}$ In fact, for every $\varepsilon > 0$ we have 17

$$\frac{1}{18} \qquad \{y \in B_{\rho}(x) : |v(y) - f(\tilde{u}(x))| > \varepsilon\} \subset \{y \in B_{\rho}(x) : |u(y) - \tilde{u}(x)| > \varepsilon/K\},$$

<u>19</u> hence $\underline{20}$

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$$\lim_{\rho \to 0^+} \frac{|\{y \in B_\rho(x) : |v(y) - f(\tilde{u}(x))| > \varepsilon\}|}{\rho^n} = 0$$

<u>22</u> whenever $x \in \Omega \setminus S_u$. By a similar argument, if $x \in S_u$ is a point such that $\underline{23}$ there exists a triplet (u^+, u^-, ν_u) satisfying (14), (15), then $\underline{24}$

$$(v^+(x) - v^-(x)) \otimes \nu_v = (f(u^+(x)) - f(u^-(x))) \otimes \nu_u \quad \text{if } x \in S_v$$

$$\frac{26}{27} \text{ and } f(u^-(x)) = f(u^+(x)) \text{ if } x \in S_u \backslash S_v. \text{ Hence, by (1.8) we get}$$

and Lemma 6.1 is proved. <u>33</u>

To prove (31), it is not restrictive to assume that k = 1. Moreover, to $\underline{35}$ simplify our notation, from now on we shall assume that $\Omega = \mathbf{R}^n$. The proof 36 of (31) is divided into two steps. In the first step we prove the statement in <u>37</u> the one-dimensional case (n = 1), using Theorem 5.2. In the second step we 38 achieve the general result using Theorem 7.1. <u>39</u>

<u>40</u> Step 1. Assume that n = 1. Since S_u is at most countable, (7) yields <u>41</u> that $\left|\widetilde{D}v\right|(S_u \setminus S_v) = 0$, so that (19) and (21) imply that $Dv = \widetilde{D}v + Jv$ is the <u>42</u>

1 Radon-Nikodým decomposition of Dv in absolutely continuous and singular $\underline{2}$ part with respect to $|\widetilde{D}u|$. By Theorem 5.2, we have <u>3</u>

$$\frac{\widetilde{D}v}{\left|\widetilde{D}u\right|}(t) = \lim_{s \to t^+} \frac{Dv([t,s[))}{\left|\widetilde{D}u\right|([t,s[))}, \qquad \frac{\widetilde{D}u}{\left|\widetilde{D}u\right|}(t) = \lim_{s \to t^+} \frac{Du([t,s[))}{\left|\widetilde{D}u\right|([t,s[))}$$

 $\left|\widetilde{D}u\right|$ -almost everywhere in **R**. It is well known (see, for instance, [20, 2.5.16]) that every one-dimensional function of bounded variation w has a unique left continuous representative, i.e., a function \hat{w} such that $\hat{w} = w$ almost every-<u>10</u> where and $\lim_{s\to t^-} \hat{w}(s) = \hat{w}(t)$ for every $t \in \mathbf{R}$. These conditions imply <u>11</u>

$$\frac{12}{13} \quad (32) \qquad \hat{u}(t) = Du(]-\infty, t[), \qquad \hat{v}(t) = Dv(]-\infty, t[) \qquad \forall t \in \mathbf{R}$$

and 14

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$$\frac{15}{16} \quad (33) \qquad \qquad \hat{v}(t) = f(\hat{u}(t)) \qquad \forall t \in \mathbf{R}.$$

Let $t \in \mathbf{R}$ be such that $\left|\widetilde{D}u\right|([t,s[) > 0 \text{ for every } s > t \text{ and assume that the } \right|$ $\underline{17}$ 18 limits in (22) exist. By (23) and (24) we get <u>19</u>

$$\begin{array}{ll} \frac{20}{21} & \frac{\hat{v}(s) - \hat{v}(t)}{\left|\widetilde{D}u\right| \left([t,s[)\right)} = \frac{f(\hat{u}(s)) - f(\hat{u}(t))}{\left|\widetilde{D}u\right| \left([t,s[)\right)} \\ \end{array}$$

$$\begin{array}{l} \frac{21}{22} & \frac{23}{24} \\ \frac{25}{26} \\ \frac{27}{26} \\ \frac{27}{28} \\ \frac{29}{30} \\ \frac{31}{31} \end{array}$$

$$= \frac{f(\hat{u}(s)) - f(\hat{u}(t) + \frac{\widetilde{D}u}{\left|\widetilde{D}u\right|}(t) \left|\widetilde{D}u\right| \left([t,s[)\right) - f(\hat{u}(t)\right)}{\left|\widetilde{D}u\right| \left([t,s[)\right) - f(\hat{u}(t)\right)} \\ + \frac{f(\hat{u}(t) + \frac{\widetilde{D}u}{\left|\widetilde{D}u\right|}(t) \left|\widetilde{D}u\right| \left([t,s[)\right) - f(\hat{u}(t)\right)}{\left|\widetilde{D}u\right| \left([t,s[)\right)} \\ \end{array}$$

for every s > t. Using the Lipschitz condition on f we find <u>32</u>

$$\begin{array}{c|c} \frac{33}{34} \\ \frac{35}{36} \\ \frac{37}{36} \\ \frac{37}{38} \\ \frac{39}{40} \\ \frac{41}{42} \end{array} & \left| \begin{array}{c} \hat{v}(s) - \hat{v}(t) \\ \left| \widetilde{D}u \right| \left([t, s[) \right) - \frac{f(\hat{u}(t) + \frac{\widetilde{D}u}{\left| \widetilde{D}u \right|}(t) \left| \widetilde{D}u \right| \left([t, s[) \right) - f(\hat{u}(t)) \right|}{\left| \widetilde{D}u \right| \left([t, s[) \right)} \right| \\ \frac{39}{40} \\ \frac{41}{\left| \widetilde{D}u \right|} \left(\frac{\hat{u}(s) - \hat{u}(t)}{\left| \widetilde{D}u \right| \left([t, s[) \right)} - \frac{\widetilde{D}u}{\left| \widetilde{D}u \right|}(t) \right|. \end{array} \right)$$

By (29), the function $s \to \left| \widetilde{D}u \right| ([t, s[) \text{ is continuous and converges to 0 as } s \downarrow t.$ Therefore Remark 7.1 and the previous inequality imply

$$\frac{\widetilde{D}v}{\left|\widetilde{D}u\right|}(t) = \lim_{h \to 0^+} \frac{f(\hat{u}(t) + h\frac{\widetilde{D}u}{\left|\widetilde{D}u\right|}(t)) - f(\hat{u}(t))}{h} \quad \left|\widetilde{D}u\right| \text{-a.e. in } \mathbf{R}.$$

⁸ By (22), $\hat{u}(x) = \tilde{u}(x)$ for every $x \in \mathbf{R} \setminus S_u$; moreover, applying the same argument to the functions u'(t) = u(-t), v'(t) = f(u'(t)) = v(-t), we get

$$\frac{\frac{11}{12}}{\frac{13}{14}} \qquad \frac{\widetilde{D}v}{\left|\widetilde{D}u\right|}(t) = \lim_{h \to 0} \frac{f(\widetilde{u}(t) + h\frac{\widetilde{D}u}{\left|\widetilde{D}u\right|}(t)) - f(\widetilde{u}(t))}{h} \qquad \left|\widetilde{D}u\right| \text{-a.e. in } \mathbf{R}$$

 $\frac{15}{15}$ and our statement is proved.

17 Step 2. Let us consider now the general case n > 1. Let $\nu \in \mathbf{R}^n$ be such 18 that $|\nu| = 1$, and let $\pi_{\nu} = \{y \in \mathbf{R}^n : \langle y, \nu \rangle = 0\}$. In the following, we shall 19 identify \mathbf{R}^n with $\pi_{\nu} \times \mathbf{R}$, and we shall denote by y the variable ranging in π_{ν} 20 and by t the variable ranging in \mathbf{R} . By the just proven one-dimensional result, 21 and by Theorem 3.3, we get

$$\frac{\frac{22}{23}}{\frac{24}{25}} \lim_{h \to 0} \frac{f(\tilde{u}(y+t\nu)+h\frac{\widetilde{D}u_y}{\left|\widetilde{D}u_y\right|}(t)) - f(\tilde{u}(y+t\nu))}{h} = \frac{\widetilde{D}v_y}{\left|\widetilde{D}u_y\right|}(t) \qquad \left|\widetilde{D}u_y\right| \text{-a.e. in } \mathbf{R}$$

 $\frac{1}{27}$ for \mathcal{H}_{n-1} -almost every $y \in \pi_{\nu}$. We claim that

$$\frac{\frac{28}{29}}{\frac{30}{30}} (34) \qquad \qquad \frac{\langle \widetilde{D}u, \nu \rangle}{\left| \langle \widetilde{D}u, \nu \rangle \right|} (y + t\nu) = \frac{\widetilde{D}u_y}{\left| \widetilde{D}u_y \right|} (t) \qquad \left| \widetilde{D}u_y \right| \text{-a.e. in } \mathbf{R}$$

<u>31</u> for \mathcal{H}_{n-1} -almost every $y \in \pi_{\nu}$. In fact, by (16) and (18) we get

$$\frac{\overline{D}u_{y}}{34} = \int_{\pi_{\nu}} \frac{\widetilde{D}u_{y}}{\left|\widetilde{D}u_{y}\right|} \cdot \left|\widetilde{D}u_{y}\right| d\mathcal{H}_{n-1}(y) = \int_{\pi_{\nu}} \widetilde{D}u_{y} d\mathcal{H}_{n-1}(y) \\
\frac{\overline{D}u_{y}}{35} = \langle \widetilde{D}u, \nu \rangle = \frac{\langle \widetilde{D}u, \nu \rangle}{\left|\langle \widetilde{D}u, \nu \rangle\right|} \cdot \left|\langle \widetilde{D}u, \nu \rangle\right| = \int_{\pi_{\nu}} \frac{\langle \widetilde{D}u, \nu \rangle}{\left|\langle \widetilde{D}u, \nu \rangle\right|} (y + \cdot \nu) \cdot \left|\widetilde{D}u_{y}\right| d\mathcal{H}_{n-1}(y)$$

$$\frac{\overline{D}u_{y}}{38} = u_{\nu}(u_{\nu}) \int_{U} \int_$$

and (24) follows from (13). By the same argument it is possible to prove that

$$\frac{40}{41} (35) \qquad \qquad \frac{\langle \widetilde{D}v, \nu \rangle}{\left| \langle \widetilde{D}u, \nu \rangle \right|} (y + t\nu) = \frac{\widetilde{D}v_y}{\left| \widetilde{D}u_y \right|} (t) \qquad \left| \widetilde{D}u_y \right| \text{-a.e. in } \mathbf{R}$$

1

 $\underline{2}$

 $\frac{\underline{3}}{\underline{4}}$ $\underline{5}$ $\underline{6}$ $\underline{7}$

16

1 for \mathcal{H}_{n-1} -almost every $y \in \pi_{\nu}$. By (24) and (25) we get 2 $\int_{\overline{\Omega}} \frac{f(\widetilde{u}(y+t\nu)+h\frac{\langle Du,\nu\rangle}{\left|\langle \widetilde{D}u,\nu\rangle\right|}(y+t\nu)) - f(\widetilde{u}(y+t\nu))}{h} = \frac{\langle \widetilde{D}v,\nu\rangle}{\left|\langle \widetilde{D}u,\nu\rangle\right|}(y+t\nu)$ 3 4 lim -<u>5</u> $h \rightarrow 0$ <u>6</u> 7 for \mathcal{H}_{n-1} -almost every $y \in \pi_{\nu}$, and using again (14), (15) we get 8 <u>9</u> $\lim_{h \to 0} \frac{f(\tilde{u}(x) + h \frac{\langle Du, \nu \rangle}{\left| \langle \widetilde{D}u, \nu \rangle \right|}(x)) - f(\tilde{u}(x))}{h} = \frac{\langle \widetilde{D}v, \nu \rangle}{\left| \langle \widetilde{D}u, \nu \rangle \right|}(x)$ 10 <u>11</u> $\underline{12}$ $\underline{13}$ $\left| \langle \widetilde{D}u, \nu \rangle \right|$ -a.e. in \mathbf{R}^n . 14Since the function $\left|\langle \widetilde{D}u,\nu\rangle\right|/\left|\widetilde{D}u\right|$ is strictly positive $\left|\langle \widetilde{D}u,\nu\rangle\right|$ -almost ev-<u>15</u> <u>16</u> erywhere, we obtain also 17 $\lim_{h \to 0} \frac{f(\widetilde{u}(x) + h \frac{\left| \langle Du, \nu \rangle \right|}{\left| \widetilde{D}u \right|}(x) \frac{\langle \widetilde{D}u, \nu \rangle}{\left| \langle \widetilde{D}u, \nu \rangle \right|}(x)) - f(\widetilde{u}(x))}{h}$ 18 19<u>20</u> $\underline{21}$ <u>22</u> $=\frac{\left|\langle \widetilde{D}u,\nu\rangle\right|}{\left|\widetilde{D}u\right|}(x)\frac{\langle \widetilde{D}v,\nu\rangle}{\left|\langle \widetilde{D}u,\nu\rangle\right|}(x)$ $\underline{23}$ 24 <u>25</u> $\langle \widetilde{D}u, \nu \rangle$ -almost everywhere in \mathbf{R}^n . <u>26</u> $\underline{27}$ Finally, since 28 $\frac{\left|\langle \widetilde{D}u,\nu\rangle\right|}{\left|\widetilde{D}u\right|}\frac{\langle \widetilde{D}u,\nu\rangle}{\left|\langle \widetilde{D}u,\nu\rangle\right|} = \frac{\langle \widetilde{D}u,\nu\rangle}{\left|\widetilde{D}u\right|} = \left\langle\frac{\widetilde{D}u}{\left|\widetilde{D}u\right|},\nu\right\rangle \qquad \left|\widetilde{D}u\right| \text{-a.e. in } \mathbf{R}^{n}$ <u>29</u> <u>30</u> <u>31</u> $\frac{\left|\langle \widetilde{D}u,\nu\rangle\right|}{\left|\widetilde{D}u\right|}\frac{\langle \widetilde{D}v,\nu\rangle}{\left|\langle \widetilde{D}u,\nu\rangle\right|} = \frac{\langle \widetilde{D}v,\nu\rangle}{\left|\widetilde{D}u\right|} = \left\langle \frac{\widetilde{D}v}{\left|\widetilde{D}u\right|},\nu\right\rangle \qquad \left|\widetilde{D}u\right| \text{-a.e. in } \mathbf{R}^{n}$ <u>32</u> <u>33</u> 34<u>35</u> and since both sides of (33) are zero $|\widetilde{D}u|$ -almost everywhere on $|\langle \widetilde{D}u, \nu \rangle|$ -<u>36</u> negligible sets, we conclude that <u>37</u> <u>38</u> $\lim_{t \to 0} \frac{f\left(\tilde{u}(x) + h\left\langle \frac{Du}{\left|\widetilde{D}u\right|}(x), \nu\right\rangle\right) - f(\tilde{u}(x))}{h} = \left\langle \frac{\widetilde{D}v}{\left|\widetilde{D}u\right|}(x), \nu\right\rangle,$ 39 <u>40</u> <u>41</u> <u>42</u>

 $\begin{array}{l} \frac{1}{2} \quad \left| \widetilde{D}u \right| \text{-a.e. in } \mathbf{R}^n. \text{ Since } \nu \text{ is arbitrary, by Remarks 7.2 and 7.3 the restriction of} \\ \frac{1}{3} \quad f \text{ to the affine space } T^u_x \text{ is differentiable at } \widetilde{u}(x) \text{ for } \left| \widetilde{D}u \right| \text{-almost every } x \in \mathbf{R}^n \\ \frac{1}{4} \quad \text{and (26) holds.} \end{array}$

It follows from (13), (14), and (15) that

$$\frac{\frac{6}{7}}{\frac{8}{2}} \quad (36) \qquad D(t_1,\ldots,t_n) = \sum_{I \in \mathbf{n}} (-1)^{|I|-1} |I| \prod_{i \in I} t_i \prod_{j \in I} (D_j + \lambda_j t_j) \det \mathbf{A}^{(\lambda)}(\overline{I}|\overline{I}).$$

<u>9</u> Let $t_i = \hat{x}_i, i = 1, \dots, n$. Lemma 1 leads to

(37)
$$D(\hat{x}_1,\ldots,\hat{x}_n) = \prod_{i \in \mathbf{n}} \hat{x}_i \sum_{I \in \mathbf{n}} (-1)^{|I|-1} |I| \operatorname{per} \mathbf{A}^{(\lambda)}(I|I) \det \mathbf{A}^{(\lambda)}(\overline{I}|\overline{I}).$$

 $\frac{12}{13}$ By (3), (13), and (37), we have the following result:

Theorem 7.2.

$$\frac{15}{16} \quad (38) \qquad \qquad H_c = \frac{1}{2n} \sum_{l=1}^n l(-1)^{l-1} A_l^{(\lambda)},$$

 $\frac{1}{18}$ where

$$\begin{array}{l} \hline 19\\ \underline{19}\\ \underline{20} \end{array} (39) \qquad \qquad A_l^{(\lambda)} = \sum_{I_l \subseteq \mathbf{n}} \operatorname{per} \mathbf{A}^{(\lambda)}(I_l|I_l) \det \mathbf{A}^{(\lambda)}(\overline{I}_l|\overline{I}_l), |I_l| = l. \end{array}$$

²¹ It is worth noting that $A_l^{(\lambda)}$ of (39) is similar to the coefficients b_l of the characteristic polynomial of (10). It is well known in graph theory that the coefficients b_l can be expressed as a sum over certain subgraphs. It is interesting to see whether A_l , $\lambda = 0$, structural properties of a graph.

²⁵ We may call (38) a parametric representation of H_c . In computation, the ²⁶ parameter λ_i plays very important roles. The choice of the parameter usually ²⁷ depends on the properties of the given graph. For a complete graph K_n , let ²⁸ $\lambda_i = 1, i = 1, ..., n$. It follows from (39) that

$$\frac{30}{31} (40) A_l^{(1)} = \begin{cases} n!, & \text{if } l = 1\\ 0, & \text{otherwise.} \end{cases}$$

 $\frac{32}{33}$ By (38)

$$\frac{34}{35}$$
 (41) $H_c = \frac{1}{2}(n-1)!$

₃₆ For a complete bipartite graph $K_{n_1n_2}$, let $\lambda_i = 0, i = 1, \ldots, n$. By (39),

$$\frac{37}{38} (42) A_l = \begin{cases} -n_1! n_2! \delta_{n_1 n_2}, & \text{if } l = 2\\ 0, & \text{otherwise} \end{cases}$$

 $_{40}$ Theorem 7.2 leads to

$$\frac{41}{42} \quad (43) \qquad \qquad H_c = \frac{1}{n_1 + n_2} n_1! n_2! \delta_{n_1 n_2}.$$

 $\underline{5}$

 $\underline{14}$

4 <u>5</u> <u>6</u> 7 8 <u>9</u>

14

<u>15</u>

1

2 3

(44) det
$$\mathbf{K}(t=1,t_1,\ldots,t_n;l|l)$$

$$=\sum_{I\subseteq \mathbf{n}-\{l\}}(-1)^{|I|}\prod_{i\in I}t_i\prod_{j\in I}(D_j+\lambda_jt_j)\det \mathbf{A}^{(\lambda)}(\overline{I}\cup\{l\}|\overline{I}\cup\{l\}).$$

By (3) and (16) we have the following asymmetrical result:

THEOREM 7.3.

$$\frac{10}{11} \quad (45) \quad H_c = \frac{1}{2} \sum_{I \subseteq \mathbf{n} - \{l\}} (-1)^{|I|} \operatorname{per} \mathbf{A}^{(\lambda)}(I|I) \det \mathbf{A}^{(\lambda)}(\overline{I} \cup \{l\} | \overline{I} \cup \{l\})$$

 $\underline{12}$ which reduces to Goulden-Jackson's formula when $\lambda_i = 0, i = 1, ..., n$ [14]. $\underline{13}$

8. Various font features of the amsmath package

8.1. Bold versions of special symbols. In the amsmath package \boldsymbol 16 is used for getting individual bold math symbols and bold Greek letters-17everything in math except for letters of the Latin alphabet, where you'd use 18\mathbf. For example, $\underline{19}$

```
A_\infty + \pi A_0 \sim
<u>20</u>
```

\mathbf{A}_{\boldsymbol{\infty}} \boldsymbol{+} $\underline{21}$

\boldsymbol{\pi} \mathbf{A}_{\boldsymbol{0}} 22

 $\underline{23}$ looks like this:

 $\underline{24}$ <u>25</u>

 $A_{\infty} + \pi A_0 \sim \mathbf{A}_{\infty} + \pi \mathbf{A}_0$

8.2. "Poor man's bold". If a bold version of a particular symbol doesn't <u>26</u> exist in the available fonts, then \boldsymbol can't be used to make that 27symbol bold. At the present time, this means that \boldsymbol can't be used 28with symbols from the msam and msbm fonts, among others. In some cases, $\underline{29}$ poor man's bold (\pmb) can be used instead of \boldsymbol: 30

$$\begin{array}{c|cccc} \frac{31}{32} & & & & & & \\ \hline 32 & & & & & \\ \hline 33 & & & & \\ \hline 34 & & & & \\ \hline 35 & & & & \\ \hline 35 & & & & \\ \hline 36 & & & & \\ \hline 37 & & & & \\ \hline 37 & & & & \\ \hline 30 & & & & \\ \hline 37 & & & & \\ \hline 30 & & & & \\ \hline 31 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 33 & & & & \\ \hline 33 & & & & \\ \hline 34 & & & & \\ \hline 34 & & & & \\ \hline 37 & & & & \\ \hline 37 & & & & \\ \hline 39 & & & & \\ \hline 39 & & & & \\ \hline 31 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 32 & & & & \\ \hline 32 & & & & \\ \hline 33 & & & & \\ \hline 31 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 32 & & & & \\ \hline 32 & & & & \\ \hline 33 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 32 & & & & \\ \hline 32 & & & & \\ \hline 33 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 33 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 31 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 31 & & & & \\ \hline 32 & & & & \\ \hline 31 & & & & \\ 31 & & & & \\ \hline 31 & & & & \\ \hline 31 & & & & \\ 31 & & & & \\ \hline 31 & & & & \\ 31 & & & &$$

$$\frac{\underline{40}}{\underline{41}} \qquad \qquad \sum_{\substack{i < B \\ i \text{ odd}}} \prod_{\kappa} \kappa F(r_i) \qquad \sum_{\substack{i < B \\ i \text{ odd}}} \prod_{\kappa} \kappa(r_i)$$

Proof: page numbers may be temporary

 \sum and \prod require an additional

AMS AND BORIS VEYTSMAN

1 \[\sum_{\substack{i<B\\\text{\$i\$ odd}}}</pre>

2 \prod_\kappa \kappa F(r_i)\qquad

 $\frac{3}{1} \ \$

4 \mathop{\pmb{\prod}}_\kappa \kappa(r_i)

<u>5</u> \]

36

<u>6</u> 7

9. Compound symbols and other features

 $\frac{8}{9}$ 9.1. *Multiple integral signs*. \iint, \iiint, and \iiiint give multiple 10 integral signs with the spacing between them nicely adjusted, in both text and 11 display style. \idotsint gives two integral signs with dots between them.

$$\begin{array}{ccc}
\underline{12}\\
\underline{13}\\
\underline{14}\\
\underline{14}\\$$

$$\underbrace{\frac{11}{15}}_{A} (47) \qquad \iiint_{A} f(w, x, y, z) \, dw \, dx \, dy \, dz \qquad \int \cdots \int f(x_1, \dots, x_k) \, dx \, dy \, dz$$

 $\frac{17}{18}$ 9.2. Over and under arrows. Some extra over and under arrow operations are provided in the amsmath package. (Basic IAT_EX provides \overrightarrow and \overleftarrow).

$$\overrightarrow{\psi_{\delta}(t)E_th} = \underbrace{\psi_{\delta}(t)E_th}_{\overleftarrow{\psi_{\delta}(t)E_th}}$$

$$\underbrace{23}{\psi_{\delta}(t)E_{t}h} = \underbrace{\psi_{\delta}(t)E_{t}h}_{th}$$

$$\overleftarrow{\psi_{\delta}(t)E_th} = \underbrace{\psi_{\delta}(t)E_th}$$

- $\frac{26}{\log \pi}$
- $\frac{27}{1}$ \overrightarrow{\psi_\delta(t) E_t h}&

 $\frac{29}{100}$ \overleftarrow{\psi_\delta(t) E_t h}&

30 =\underleftarrow{\psi_\delta(t) E_t h}\\

```
31 \overleftrightarrow{\psi_\delta(t) E_t h}&
```

```
32 =\underleftrightarrow{\psi_\delta(t) E_t h}
```

```
\frac{33}{\sqrt{2}}
```

 $\frac{34}{35}$ These all scale properly in subscript sizes:

 $\int_{\overrightarrow{AB}} ax \, dx$

<u>38</u> \[\int_{\overrightarrow{AB}} ax\,dx\]

 $\begin{array}{l} \frac{39}{40} \\ \frac{40}{41} \\ \frac{41}{42} \end{array}$ 9.3. Dots. Normally you need only type \dots for ellipsis dots in a math formula. The main exception is when the dots fall at the end of the formula; then you need to specify one of \dotsc (series dots, after a comma), \dotsb

1 (binary dots, for binary relations or operators), \dotsm (multiplication dots), $\underline{2}$ or \dotsi (dots after an integral). For example, the input <u>3</u> Then we have the series \$A_1,A_2,\dotsc\$, 4 the regional sum $A_1+A_2+\$, <u>5</u> the orthogonal product \$A_1A_2\dotsm\$, <u>6</u> and the infinite integral 7 $\left[\left[\left[A_1\right]\right]\right].$ 8 produces 9 Then we have the series A_1, A_2, \ldots , the regional sum $A_1 + A_2 +$ 10 \cdots , the orthogonal product $A_1 A_2 \cdots$, and the infinite integral 11 $\underline{12}$ $\int_{A_1} \int_{A_2} \dots$ $\underline{13}$ 149.4. Accents in math. Double accents: $\underline{15}$ $\hat{\hat{H}} \quad \check{\check{C}} \quad \tilde{\check{T}} \quad \acute{\hat{A}} \quad \grave{\hat{G}} \quad \dot{\check{D}} \quad \ddot{\check{D}} \quad \breve{\check{B}} \quad \vec{\bar{B}} \quad \vec{\bar{V}}$ 16 17\[\Hat{\Hat{H}}\quad\Check{\Check{C}}\quad 18\Tilde{\Tilde{T}}\quad\Acute{\Acute{A}}\quad $\underline{19}$ \Grave{\Grave{G}}\quad\Dot{\Dot{D}}\quad 20\Ddot{\Ddot{D}}\quad\Breve{\Breve{B}}\quad 21 $Bar{Bar{B}}/quadVec{Vec{V}}]$ 2223This double accent operation is complicated and tends to slow down the processing of a LATEX file. $\underline{24}$ <u>25</u> 9.5. Dot accents. \dddot and \dddot are available to produce triple and <u>26</u> quadruple dot accents in addition to the \dot and \ddot accents already avail- $\underline{27}$ able in $\mathbb{A}T_{E}X$: 28 \ddot{O} \overrightarrow{R} <u>29</u> 30 $\left[\ddot{Q} \right]$ <u>31</u> 9.6. Roots. In the amsmath package \leftroot and \uproot allow you to <u>32</u> adjust the position of the root index of a radical: 33 \sqrt[\leftroot{-2}\uproot{2}\beta]{k} 34gives good positioning of the β : <u>35</u> <u>36</u> $\sqrt[\beta]{k}$ <u>37</u> 9.7. Boxed formulas. The command \boxed puts a box around its argu-38ment, like \fbox except that the contents are in math mode: 39 \boxed{W_t-F\subseteq V(P_i)\subseteq W_t} <u>40</u> <u>41</u> $W_t - F \subseteq V(P_i) \subseteq W_t$. <u>42</u> Proof: page numbers may be temporary

1 9.8. Extensible arrows. \xleftarrow and \xrightarrow produce arrows $\underline{2}$ that extend automatically to accommodate unusually wide subscripts or su-<u>3</u> perscripts. The text of the subscript or superscript are given as an optional 4 resp. mandatory argument: Example: $\underline{5}$ $0 \xleftarrow{\alpha}{\zeta} F \times \triangle[n-1] \xrightarrow{\partial_0 \alpha(b)} E^{\partial_0 b}$ <u>6</u> 7 \[0 \xleftarrow[\zeta]{\alpha} F\times\triangle[n-1] 8 \xrightarrow{\partial_0\alpha(b)} E^{\partial_0b}\] 9 109.9. \overset, \underset, and \sideset. Examples: 11 $\begin{array}{cccc} & * & X & X \\ X & X & X \\ & * & & h \end{array}$ <u>12</u> $\underline{13}$ \[\overset{*}{X}\qquad\underset{*}{X}\qquad 14 $\ensuremath{\{underset\{b\}\{X\}\}\}\}$ 1516The command \sideset is for a rather special purpose: putting symbols <u>17</u> at the subscript and superscript corners of a large operator symbol such as \sum <u>18</u> or \prod , without affecting the placement of limits. Examples: $\underline{19}$ $\prod_{k=0}^{*} \prod_{k=0}^{*} \sum_{0 \le i \le m}^{\prime} E_i \beta x$ <u>20</u> $\underline{21}$ <u>22</u> $[\sideset{_*^*}{_*^*}\prod_k\quad$ <u>23</u> \sideset{}{'}\sum_{0\le i\le m} E_i\beta x $\underline{24}$)] $\underline{25}$ 9.10. The \text command. The main use of the command \text is for $\underline{26}$ words or phrases in a display: 27 $\underline{28}$ $\mathbf{y} = \mathbf{y}'$ if and only if $y'_k = \delta_k y_{\tau(k)}$ <u>29</u> \[\mathbf{y}=\mathbf{y}'\quad\text{if and only if}\quad 30 y'_k=\delta_k y_{\tau(k)}\] $\underline{31}$ <u>32</u> 9.11. Operator names. The more common math functions such as log, sin, <u>33</u> and lim have predefined control sequences: \log, \sin, \lim. The amsmath $\underline{34}$ package provides \DeclareMathOperator and \DeclareMathOperator* for $\underline{35}$ producing new function names that will have the same typographical treat-36 ment. Examples: <u>37</u> $||f||_{\infty} = \operatorname{ess\,sup}_{x \in B^n} |f(x)|$ <u>38</u> $\underline{39}$ $\[\norm{f}_\infty=$ $\sum_{x\in \mathbb{R}^n} b\{f(x)\}$ <u>40</u> 41 $\operatorname{meas}_1\{u \in R^1_+ \colon f^*(u) > \alpha\} = \operatorname{meas}_n\{x \in R^n \colon |f(x)| \ge \alpha\} \quad \forall \alpha > 0.$ 42

```
1
     \left[ \sum_1^{u \in R_+^1 \subset n \in \mathbb{N} \right]
\underline{2}
     =\meas_n\{x\in R^n\colon \black{f(x)}\geq\alpha\}
3
     \quad \forall\alpha>0.\]
\underline{4}
     \essup and \meas would be defined in the document preamble as
\underline{5}
     \DeclareMathOperator*{\esssup}{ess\,sup}
6
     \DeclareMathOperator{\meas}{meas}
7
8
           The following special operator names are predefined in the amsmath pack-
<u>9</u>
     age: \varlimsup, \varliminf, \varinjlim, and \varprojlim. Here's what
<u>10</u>
     they look like in use:
11

\lim_{n \to \infty} \mathcal{Q}(u_n, u_n - u^{\#}) \le 0

\lim_{n \to \infty} |a_{n+1}| / |a_n| = 0

     (48)
12
\underline{13}
     (49)
14
                                     \underline{\lim}(m_i^{\lambda} \cdot)^* \le 0
     (50)
\underline{15}
                                     \varprojlim_{p \in S(A)} A_p \le 0
<u>16</u>
     (51)
17
18
     \begin{align}
\underline{19}
     &\varlimsup_{n\rightarrow\infty}
20
               \mathbb{Q}(u_n,u_n-u^{+})\le0
\underline{21}
     &\varliminf_{n\rightarrow\infty}
<u>22</u>
        \left\lvert a_{n+1}\right\rvert/\left\lvert a_n\right\rvert=0\\
23
     &\varinjlim (m_i^\lambda\cdot)^*\le0\\
\underline{24}
     &\varprojlim_{p\in S(A)}A_p\le0
<u>25</u>
     \end{align}
<u>26</u>
27
          9.12. \mod and its relatives. The commands \mod and \pod are variants
\underline{28}
     of \pmod preferred by some authors; \mod omits the parentheses, whereas \pod
<u>29</u>
     omits the 'mod' and retains the parentheses. Examples:
30
                                      x \equiv y + 1 \pmod{m^2}
     (52)
<u>31</u>
                                      x \equiv y + 1 \mod m^2
<u>32</u>
     (53)
<u>33</u>
                                      x \equiv y + 1 \quad (m^2)
     (54)
34
<u>35</u>
     \begin{align}
<u>36</u>
     x&\equiv y+1\pmod{m^2}\\
<u>37</u>
     x&\equiv y+1\mod{m^2}\\
<u>38</u>
     x&\equiv y+1\pod{m^2}
<u>39</u>
     \end{align}
<u>40</u>
          9.13. Fractions and related constructions. The usual notation for binomi-
<u>41</u>
     als is similar to the fraction concept, so it has a similar command \binom with
42
```

```
1
    two arguments. Example:
\underline{2}
                      \sum_{\gamma \in \Gamma_{G}} I_{\gamma} = 2^{k} - \binom{k}{1} 2^{k-1} + \binom{k}{2} 2^{k-2}
<u>3</u>
4
                                 +\cdots + (-1)^{l} \binom{k}{l} 2^{k-l} + \cdots + (-1)^{k}
\underline{5}
    (55)
<u>6</u>
7
                               = (2-1)^k = 1
8
    \begin{equation}
9
   \begin{split}
10
   [\sum_{\gamma\in\Gamma_C} I_\gamma&
11
   =2^k-\binom{k}{1}2^{k-1}+\binom{k}{2}2^{k-2}\\
12
   &\quad+\dots+(-1)^1\binom{k}{1}2^{k-1}
13
14 +\dots+(-1)^k\\
15 &=(2-1)^k=1
16 \end{split}
17 \end{equation}
\underline{18}
   There are also abbreviations
\underline{19}
    \dfrac
                       \dbinom
20
    \tfrac
                       \tbinom
\underline{21}
    for the commonly needed constructions
22
    {\displaystyle\frac ... }
                                          {\displaystyle\binom ... }
23
    {\textstyle\frac ... }
                                           {\textstyle\binom ... }
24
         The generalized fraction command \genfrac provides full access to the
25
    six T<sub>F</sub>X fraction primitives:
\underline{26}
27
                                               \overwithdelims: \left< \frac{n+1}{2} \right>
                 \over: \frac{n+1}{2}
    (56)
\underline{28}
<u>29</u>
                 \texttt{Atop:} \ \frac{n+1}{2}
                                               \atopwithdelims: \binom{n+1}{2}
30
    (57)
31
                                              \abovewithdelims: \left|\frac{n+1}{2}\right|
<u>32</u>
                \above: \frac{n+1}{2}
    (58)
33
\underline{34}
    \text{\cn{over}: }&\genfrac{}{}{}{n+1}{2}&
35
    \text{\cn{overwithdelims}: }&
36
       \langle}{\rangle}{}{n+1}{2}\
<u>37</u>
    \text{\cn{atop}: }&\genfrac{}{}{0pt}{}{n+1}{2}&
38
    \text{\cn{atopwithdelims}: }&
<u>39</u>
       genfrac{(}{)}{0pt}{}{n+1}{2}\
<u>40</u>
    \text{\cn{above}: }&\genfrac{}{}{1pt}{}{n+1}{2}&
41
    \text{\cn{abovewithdelims}: }&
42
```

1 $genfrac{[}{]}{1pt}{}{n+1}{2}$ 2 9.14. Continued fractions. The continued fraction 3 4 $\frac{1}{\sqrt{2} + \frac{1}{\sqrt{2} + \frac{1}{\sqrt{2} + \frac{1}{\sqrt{2} + \frac{1}{\sqrt{2} + \cdots}}}}$ (59)<u>5</u> 6 7 8 <u>9</u> 10<u>11</u> can be obtained by typing $\underline{12}$ $cfrac{1}{sqrt{2}+}$ $\underline{13}$ $cfrac{1}{sqrt{2}+}$ 14 $cfrac{1}{sqrt{2}+}$ <u>15</u> $cfrac{1}{sqrt{2}+}$ 16 $cfrac{1}{\sqrt{2}+\dotsb}$ 17}}}} 18

Left or right placement of any of the numerators is accomplished by using
 <u>20</u> \cfrac[1] or \cfrac[r] instead of \cfrac.

21 9.15. Smash. In amsmath there are optional arguments t and b for the 22 plain T_EX command \smash, because sometimes it is advantageous to be able 23 to 'smash' only the top or only the bottom of something while retaining the 24 natural depth or height. In the formula $X_j = (1/\sqrt{\lambda_j})X'_j$ \smash[b] has been 25 used to limit the size of the radical symbol.

 $\overline{27}$ $X_j=(1/\sqrt{1}) X_j'$

 $\underline{31}$

<u>32</u>

Without the use of \smash[b] the formula would have appeared thus: $X_j = \frac{29}{(1/\sqrt{\lambda_j})}X'_j$, with the radical extending to encompass the depth of the subscript j.

9.16. The 'cases' environment. 'Cases' constructions like the following can be produced using the **cases** environment.

Notice the use of **\text** and the embedded math.

```
\mathbf{2}
            9.17. Matrix. Here are samples of the matrix environments, \matrix,
\underline{3}
     \pmatrix, \bmatrix, \Bmatrix, \vmatrix and \Vmatrix:
\underline{4}
\underline{5}
<u>6</u>
     (61) \begin{array}{c} \vartheta \quad \varrho \\ \varphi \quad \varpi \end{array} \begin{pmatrix} \vartheta \quad \varrho \\ \varphi \quad \varpi \end{pmatrix} = \begin{bmatrix} \vartheta \quad \varrho \\ \varphi \quad \varpi \end{bmatrix} = \begin{bmatrix} \vartheta \quad \varrho \\ \varphi \quad \varpi \end{bmatrix} = \begin{bmatrix} \vartheta \quad \varrho \\ \varphi \quad \varpi \end{bmatrix} = \begin{bmatrix} \vartheta \quad \varrho \\ \varphi \quad \varpi \end{bmatrix} = \begin{bmatrix} \vartheta \quad \varrho \\ \varphi \quad \varpi \end{bmatrix}
7
8
9
10
\frac{11}{\text{begin}}
\frac{12}{\sqrt{vartheta}} \sqrt{varrho}/\sqrt{varphi}
\frac{13}{13} \ \
\frac{14}{} \begin{pmatrix}
\frac{15}{2} \vartheta& \varrho\\\varphi& \varpi
\frac{17}{10} \ge 17
\frac{18}{18} \vartheta& \varrho\\\varphi& \varpi
\frac{19}{2} \ \
\frac{20}{2} \begin{Bmatrix}
\underline{^{21}} \vartheta& \varrho\\\varphi& \varpi
\frac{23}{\text{begin}}
\frac{24}{\sqrt{\pi}} \sqrt{\sqrt{\pi}} \sqrt{\sqrt{\pi}}
\frac{25}{2} \ \
\frac{26}{\text{Vmatrix}}
27 \vartheta& \varrho\\\varphi& \varpi
\underline{28}
     \end{Vmatrix}
\underline{29}
\underline{30}
             To produce a small matrix suitable for use in text, use the smallmatrix
\underline{31}
     environment.
<u>32</u>
<u>33</u>
     \begin{math}
\underline{34}
         \bigl( \begin{smallmatrix}
\underline{35}
                 a&b\\ c&d
36
             \end{smallmatrix} \bigr)
\underline{37}
      \end{math}
38
<u>39</u>
     To show the effect of the matrix on the surrounding lines of a paragraph, we
40
     put it here: \begin{pmatrix} a & b \\ c & d \end{pmatrix} and follow it with enough text to ensure that there will be
41
     at least one full line below the matrix.
42
```

42

1

1 \hdotsfor{number} produces a row of dots in a matrix spanning the $\underline{2}$ given number of columns: 3 $W(\Phi) = \begin{vmatrix} \frac{\varphi}{(\varphi_1, \varepsilon_1)} & 0 & \dots & 0 \\ \frac{\varphi k_{n2}}{(\varphi_2, \varepsilon_1)} & \frac{\varphi}{(\varphi_2, \varepsilon_2)} & \dots & 0 \\ \frac{\varphi k_{n1}}{(\varphi_n, \varepsilon_1)} & \frac{\varphi k_{n2}}{(\varphi_n, \varepsilon_2)} & \dots & \frac{\varphi k_{nn-1}}{(\varphi_n, \varepsilon_{n-1})} & \frac{\varphi}{(\varphi_n, \varepsilon_n)} \end{vmatrix}$ 4<u>5</u> 6 7 8 <u>9</u> 10 \[W(\Phi)= \begin{Vmatrix} 11 \dfrac\varphi{(\varphi_1,\varepsilon_1)}&0&\dots&0\\ $\underline{12}$ \dfrac{\varphi k_{n2}}{(\varphi_2,\varepsilon_1)}& 13\dfrac\varphi{(\varphi_2,\varepsilon_2)}&\dots&0\\ 14 $hdotsfor{5}\$ <u>15</u> \dfrac{\varphi k_{n1}}{(\varphi_n,\varepsilon_1)}& 16 \dfrac{\varphi k_{n2}}{(\varphi_n,\varepsilon_2)}&\dots& 17\dfrac{\varphi k_{n\,n-1}}{(\varphi_n,\varepsilon_{n-1})}& 18\dfrac{\varphi}{(\varphi_n,\varepsilon_n)} $\underline{19}$ \end{Vmatrix}\] 20The spacing of the dots can be varied through use of a square-bracket option, 21for example, \hdotsfor[1.5]{3}. The number in square brackets will be used 22 as a multiplier; the normal value is 1. 23 $\underline{24}$ 9.18. The \substack command. The \substack command can be used <u>25</u> to produce a multiline subscript or superscript: for example <u>26</u> \sum_{\substack{0\le i\le m\\ 0<j<n}} P(i,j)</pre> 2728 produces a two-line subscript underneath the sum: $\underline{29}$ $\sum_{0 \le i \le m} P(i,j)$ (62)30 <u>31</u> 32A slightly more generalized form is the subarray environment which allows 33 you to specify that each line should be left-aligned instead of centered, as here: 34<u>35</u> $\sum_{\substack{0 \le i \le m}} P(i,j)$ (63)36 <u>37</u> 38\sum_{\begin{subarray}{1} 39 $0\leq i\leq m\leq 0\leq j\leq n$ <u>40</u> \end{subarray}} <u>41</u> P(i,j)42

Maybe "... as

below"?

```
1
          9.19. Biq-q-q delimiters. Here are some big delimiters, first in \normalsize:
\underline{2}
                                    \left(\mathbf{E}_y \int_0^{t_{\varepsilon}} L_{x,y^x(s)} \varphi(x) \, ds\right)
<u>3</u>
4
    [\biggl(\mathbf{E}_{y}]
\underline{5}
       \int_0^{t_\varepsilon}L_{x,y^x(s)}\varphi(x)\,ds
\underline{6}
       \biggr)
\underline{7}
    )]
8
\underline{9}
    and now in \Large size:
10
                                  \left(\mathbf{E}_{y}\int_{0}^{t_{\varepsilon}}L_{x,y^{x}(s)}\varphi(x)\,ds\right)
\underline{11}
<u>12</u>
\underline{13}
    {\Large
14
    [\E_{y}]
15
       \int_0^{t_\sqrt{x}(s)}_{x,y^x(s)}
16
       \biggr)
<u>17</u>
    \]}
<u>18</u>
19
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20
    and NIH. They acknowledge the generous help of other agencies.
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<u>22</u>
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<u>40</u>		American Mathematical Society
<u>40</u>		AMS, Providence, Rhode Island
<u>41</u> <u>42</u>		E-mail: tech-support@ams.org
<u>+4</u>		

AMS AND BORIS VEYTSMAN

<u>1</u>	Boris Veytsman
<u>-</u> <u>2</u>	School of Systems Biology, George Mason University, Fairfax, VA 22030, USA
_ <u>3</u>	<i>E-mail</i> : borisv@lk.net
<u>4</u>	https://borisv.lk.net
<u>-</u>	https://orcid.org/0000-0003-4674-8113
<u>6</u>	Frodo Baggins
<u>7</u>	Bilbo Baggins
<u>8</u>	https://orcid.org/000-0000-0000-00000
<u>9</u>	The Unseen University
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Proof: page numbers may be temporary

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